# Math 330 - Additional Material 

Student edition with proofs

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## 1 Before You Start

Errors detected by Math 330 students, Spring 2017:

| Date | Topic |
| :---: | :--- |
| 2017-01-26 | Error in def.4.2. Incorrect version: A relation is symmetric if $x_{1} R x_{2}$ implies |
|  | $x_{1} R x_{2}$ for all $x_{1}, x_{2} \in X$. Correct version: $A$ relation is symmetric if $x_{1} R x_{2}$ |
| implies $x_{2} R x_{1}$ for all $x_{1}, x_{2} \in X$. Detected by Brad Whistance. |  |

History of Updates:

| Date | Topic |
| :---: | :---: |
| 2017-01-30 |  | | Mainly reformatting: increased use of tables. Significant additions to ch.16 |
| :--- |
| (Appendix: Addenda to Beck/Geoghegan's "The Art of Proof"). |
| 2017-01-10 |
|  |
| Many updates during Fall 2016. Significant streamlining and reorg of ch. 8 |
| (liminf, limsup, ...) through ch.13 (Zorn's Lemma). |

### 1.1 About This Document

Remark 1.1 (The purpose of this document). The original version of this document was written in 2005 under the title "Introduction to Abstract Math - A Journey to Approximation Theory". Since then parts of it were discarded and others have been added. It now serves as lecture notes for the course "Math 330: Number systems" which is held at the Department of Mathematical Sciences at Binghamton University. Parts of the remainder of this chapter are specifically addressed to the students of this course.

These notes serve at least two purposes:
a. Here you find background material on topics that cannot found in sufficient detail or not at all in the textbook [1] Beck/Geoghegan: The Art of Proof. That book serves as the primary reference for the first two thirds of the Math 330 course. It is often simply referred to as " $B / G$ " in these notes.
b. This document covers material which is beyond the scope of [1] B/G such as

- material on lim inf and lim sup
- convergence, continuity and compactness in metric spaces
- applications of Zorn's lemma

These topics are usually covered in the last third of my Math 330 class.
Remark 1.2 (Acknowledgements). The early chapters of this document draw on the following chapters of [4] Bryant, Kirby Course Notes for MAD 2104:

Ch.1, section 1: Introduction to Sets
Ch.1, section 2: Introduction to Functions
Ch.2: Logic
Ch.3: Methods of Proofs
Ch.4, section 1: Set Operations
Ch.4, section 2: Properties of Functions

Moreover such a document cannot be written with the intent to supplement the [1] B/G book without strongly borrowing from it.

Remark 1.3 (How to navigate this document I). Scrutinize the table of contents, including the headings for the subchapters. You will find many entries there tagged with a directive.

For example, the reference to ch. 8 (The Completeness of the Real Numbers System) has been tagged with (Understand this!),
the first subchapter ch.8.1 (Minima, Maxima, Infima and Suprema) has been tagged with (Study this!),
the subchapter ch.8.4 (Sequences of Sets and Indicator functions and their liminf and limsup) has no tag.

## Notation Alert:

All directives apply to the entire subtree and a lower level directive overrides the "parent directives". Example: the "Understand this!" directive of subsection 10.2.3: Continuity of Polynomials overrides the "Study this!" directive of subsection 10.2 on Continuity. Accordingly, when you do not see any comment, back up in the table of contents: first to the parent entry, then to its parent entry ... until you find one.
a. "Study this" directive: When you read "Study this", you must understand the material in depth. You will need to do this with paper and pencil in hand and make an effort not only to understand what the definitions and theorems are all about - not a minor undertaking because some of the subject matter is quite abstract - but aso make an effort to follow the proofs at least from a birds eye perspective.
b. "Understand this" directive: When you read "Understand this", you should know the definitions, propositions and theorems without worrying about proofs. Chances are that the material will be referred to from truly important sections of this write-up and is primarily needed for their understanding.
c. "Skim this" directive: When you read "Skim this", you should understand how the material is structured. You may find it easier to do some of your homework. A good example is chapter 3 on logic. It has been marked as "Skim this" but some of the later subchapters override this directive as "Understand this". Understanding the material of this subchapter will be problematic unless you can find your way around in the "Skim this" subchapters.
d. "Skip this" directive: When you read "Skip this", you need not worry about the content.
e. " $(\star)$ " directive: Chapters marked " $(\star)$ " are optional.

## Homework:

You will find almost every week reading assignments as part of your homework. The reading is due prior to when it is needed in class, both for this document and the Beck/Geoghegan text. I assume that you did your reading and I will assume in particular that you have learned the definitions so that I can move along at a fast pace except for some topics that I will focus on in detail. For this document it means that you should do the "Study this" and "Understand this" material as was indicated above.

Remark 1.4 (How to navigate this document II). I believe that, particularly in Math, more words take a lot less time to understand than a skeletal write-up like the one given in the [1] B/G text. Accordingly, almost all of the material in this document comes with quite detailed proofs. Those proofs are there for you to study.

Some of those proofs, notably those in prop. 6.2 on p.104, make use of " $\Leftrightarrow$ " to show that two sets are equal. You should study this technique but, as you will hear me say many times in class, I recommend that you abstain from using " $\Leftrightarrow$ " between statements in your proofs. You very likely lack the experience to use this technique without errors.

Some of the material was written from scratch, other material was pulled in from a document that was written as early as 10 years ago. I have make an attempt to make the entire document more homogeneous but there will be some inconsistencies. Your help in pointing out to me the most notable trouble spots would be deeply appreciated.

There are differences in style: the original document was written in a much more colloquial style as it was addressed to a younger audience of high school students who had expressed a special interest in studying college level math.

This is a living document: material will be added as I find the time to do so. Be sure to check the latest PDF frequently. You certainly should do so when an announcement was made that this document contains new additions and/or corrections.

### 1.2 How to Properly Write a Proof (Study this!)

Study this brief chapter to understand some of the dos and don'ts when submitting your homework.
To prove an equation such as $A=Z$ you are asked to do one of the following:
$a$.

$$
\begin{aligned}
A & =B & \text { (use } \ldots . .) \\
& =C & (\text { use } \ldots .) \\
& =D & (\text { use } \ldots .) \\
& \ldots & \ldots . . . . . . . . . . . . \\
& =Z & \text { (use } \ldots .)
\end{aligned}
$$

You then conclude from the transitivity of equality that $A=Z$ is indeed true.
b. You transform the left side (L.S.) and the right side (R.S.) separately and show that in each case you obtain the same item, say M:

Left side:

$$
\left.\begin{array}{rlrl}
A & =B & (\text { use } \ldots .) \\
& =C & (\text { use } \ldots .) \\
& =D & (\text { use } \ldots .) \\
& \ldots & \ldots . . . . . . . . . . . . . ~
\end{array}\right)
$$

Right side:

$$
\begin{array}{rlr}
Z & =Y & \text { (use } \ldots .) \\
& =X & (\text { use } \ldots .) \\
& =W & \text { (use ....) } \\
& \ldots . . . . . . . . . . . . . . . . ~ \\
& =M & \text { (use ....) }
\end{array}
$$

and you rightfully conclude that the proof is done because it follows from $A=M$ and $Z=M$ that $A=Z$.

Note that you are not allowed to structure your proof as follows.

$$
\begin{aligned}
A=Z & \text { (that's what you want to prove) } \\
B=Y & \text { (you do with both } A \text { and } Z \text { the same operation ...... ) } \\
C=X & \text { (you do with both } B \text { and } Y \text { the same operation ..... ) } \\
D=W & \text { (you do with both } C \text { and } X \text { the same operation ...... ) } \\
& \ldots . . . . . . . . . . . . . . . . ~
\end{aligned}
$$

## What is potentially wrong with that last approach?

In the abstract the issue is that when using method $\boldsymbol{a}$ or $\boldsymbol{b}$ you take in each step an equation that is known to be true or that you assume to be true and you rightfully conclude by the use of transitivity that you have proved what you wanted to be true.

When you use method c you take an equation that you want to be true $(A=Z)$ but have not yet proved to be so. If you are wrong then doing the same thing to both sides may potentially lead to two things that are equal.

Here is a simple example that demonstrates why method c is not allowed for a mathematical proof. This method will be used in two different ways to prove that $-2=2$.

First proof:

$$
\begin{aligned}
-2 & =2 \quad \text { (want to prove) } \\
-2 \cdot 0 & =2 \cdot 0 \quad \text { (multiply both sides from the right } w .0) \\
0 & =0 \quad(B / G \text { ax. } 1.2 \text { (additive neutral element) }
\end{aligned}
$$

We are done.
Second proof:

$$
\begin{aligned}
-2 & =2 \quad \text { (want to prove) } \\
(-2)^{2} & =2^{2} \quad \text { (square both sides) } \\
4 & =4 \quad \text { (obvious) }
\end{aligned}
$$

## We are done.

Now you know why you must never use method c. ${ }^{1}$

[^0]
## 2 Preliminaries about Sets, Numbers and Functions

### 2.1 Sets and Basic Set Operations

Ask a mathematician how her or his Math is different from the kind of Math you learn in high school, in fact, from any kind of Math you find outside textbooks for mathematicians and theoretical physicists. One of the answers you are likely to get is that Math is not so much about numbers but also about other objects, among them sets and functions. Once you know about those, you can tackle sets of functions, set functions, sets of set functions, ...

An entire book can be filled with a mathematically precise theory of sets. ${ }^{2}$ For our purposes the following "naive" definition suffices:

Definition 2.1 (Sets). A set is a collection of stuff called members or elements which satisfies the following rules: The order in which you write the elements does not matter and if you list an element two or more times then it only counts once.

We write a set by enclosing within curly braces the elements of the set. This can be done by listing all those elements or giving instructions that describe those elements. For example, to denote by $X$ the set of all integer numbers between 18 and 24 we can write either of the following:

$$
X:=\{18,19,20,21,22,23,24\} \quad \text { or } \quad X:=\{n: n \text { is an integer and } 18 \leqq n \leqq 24\}
$$

Both formulas clearly define the same collection of all integers between 18 and 24 . On the left the elements of $X$ are given by a complete list, on the right setbuilder notation, i.e., instructions that specify what belongs to the set, is used instead.

It is customary to denote sets by capital letters and their elements by small letters but this is not a hard and fast rule. You will see many exceptions to this rule in this document.

We write $x_{1} \in X$ to denote that an item $x_{1}$ is an element of the set $X$ and $x_{2} \notin X$ to denote that an item $x_{2}$ is not an element of the set $X$

For the above example we have $20 \in X, 27-6 \in X, 38 \notin X,{ }^{\prime}$ Jimmy' $^{\prime} \notin X$.
Example 2.1 (No duplicates in sets). The following collection of alphabetic letters is a set:

$$
S_{1}=\{a, e, i, o, u\}
$$

and so is this one:

$$
S_{2}=\{a, e, e, i, i, i, o, o, o, o, u, u, u, u, u\}
$$

Did you notice that those two sets are equal?
Remark 2.1. The symbol $n$ in the definition of $X=\{n: n$ is an integer and $18 \leqq n \leqq 24\}$ is a dummy variable in the sense that it does not matter what symbol you use. The following sets all are equal to $X$ :
$\{x: x$ is an integer and $18 \leqq x \leqq 24\}$,
$\{\alpha: \alpha$ is an integer and $18 \leqq \alpha \leqq 24\}$,
$\{\mathfrak{Z}: \mathfrak{Z}$ is an integer and $18 \leqq \mathfrak{Z} \leqq 24\}$

[^1]Remark 2.2 (Russell's Antinomy). Care must be taken so that, if you define a set with the use of setbuilder notation, no inconsistencies occur. Here is an example of a definition of a set that leads to contradictions.

$$
\begin{equation*}
A:=\{B: B \text { is a set and } B \notin B\} \tag{2.1}
\end{equation*}
$$

What is wrong with this definition? To answer this question let us find out whether or not this set $A$ is a member of $A$. Assume that $A$ belongs to $A$. The condition to the right of the colon states that $A \notin A$ is required for membership in $A$, so our assumption $A \in A$ must be wrong. In other words, we have established "by contradiction" that $A \notin A$ is true. But this is not the end of it: Now that we know that $A \notin A$ it follows that $A \in A$ because $A$ contains all sets that do not contain themselves.

In other words, we have proved the impossible: both $A \in A$ and $A \notin A$ are true! There is no way out of this logical impossibility other than excluding definitions for sets such as the one given above. It is very important for mathematicians that their theories do not lead to such inconsistencies. Therefore, examples as the one above have spawned very complicated theories about "good sets". It is possible for a mathematician to specialize in the field of axiomatic set theory (actually, there are several set theories) which endeavors to show that the sets are of any relevance in mathematical theories do not lead to any logical contradictions.

The great majority of mathematicians take the "naive" approach to sets which is not to worry about accidentally defining sets that lead to contradictions and we will take that point of view in this document.

Definition 2.2 (empty set). $\emptyset$ or $\}$ denotes the empty set. It is the one set that does not contain any elements.

Remark 2.3 (Elements of the empty set and their properties). You can state anything you like about the elements of the empty sets as there are none. The following statements all are true:
a: If $x \in \emptyset$ then $x$ is a positive number.
b: If $x \in \emptyset$ then $x$ is a negative number.
c: Define $a \sim b$ if and only if both are integers and $a-b$ is an even number.
For any $x, y, z \in \emptyset$ it is true that
c1: $x \sim x$,
c2: if $x \sim y$ then $y \sim x$,
c3: if $x \sim y$ and $y \sim z$ then $x \sim z$.
d: Let $A$ be any set. If $x \in \emptyset$ then $x \in A$.
As you will learn later, $\mathbf{c}$ : means that " $\sim$ " is an equivalence relation (see def.4.3 on p.68) and d: means that the empty set is a subset (see the next definition) of any other set.

Definition 2.3 (subsets and supersets). We say that a set $A$ is a subset of the set $B$ and we write $A \subseteq B$ if any element of $A$ also belongs to $B$. Equivalently we say that $B$ is a superset of the set $A$ and we write $B \supseteq A$. We also say that $B$ includes $A$ or $A$ is included by $B$. Note that $A \subseteq A$ and $\emptyset \subseteq A$ is true for any set $A$.

If $A \subseteq B$ but $A \neq B$, i.e., there is at least one $x \in B$ such that $x \notin A$, then we say that $A$ is a strict subset or a proper subset of $B$. We write " $A \subsetneq B$ " or " $A \subset B$ ". Alternatively we say that $B$ is a strict superset or a proper superset of $A$ and we write " $B \supsetneq A$ ") or " $B \supset A$ ".


Figure 2.1: Set inclusion: $A \subseteq B, B \supseteq A$

Two sets $A$ and $B$ are equal means that they both contain the same elements. In other words, $A=B$ iff $A \subseteq B$ and $B \subseteq A$.
"iff" is a short for "if and only if": P iff $Q$ for two statements $P$ and $Q$ means that if $P$ is valid then $Q$ is valid and vice versa. ${ }^{3}$

To show that two sets $A$ and $B$ are equal you show that
a. if $x \in A$ then $x \in B$,
b. if $x \in B$ then $x \in A$.

Definition 2.4 (unions, intersections and disjoint unions). Given are two arbitrary sets $A$ and $B$. No assumption is made that either one is contained in the other or that either one contains any elements!

The union $A \cup B$ (pronounced "A union B ") is defined as the set of all elements which belong to $A$ or $B$ or both.

The intersection $A \cap B$ (pronounced "A intersection B ") is defined as the set of all elements which belong to both $A$ and $B$.

We call $A$ and $B$ disjoint if $A \cap B=\emptyset$. We then usually write $A \uplus B$ (pronounced "A disjoint union $\mathrm{B}^{\prime \prime}$ ) rather than $A \cup B$.
$A \cup B:$

$A \cup B \cup C:$

$A \cap B:$
$A \cap B \cap C:$


Figure 2.2: Union and intersection of sets

[^2]Definition 2.5 (set differences and symmetric differences). Given are two arbitrary sets $A$ and $B$. No assumption is made that either one is contained in the other or that either one contains any elements!

The difference set or set difference $A \backslash B$ (pronounced "A minus B ") is defined as the set of all elements which belong to $A$ but not to $B$ :

$$
\begin{equation*}
A \backslash B:=\{x \in A: x \notin B\} \tag{2.2}
\end{equation*}
$$

The symmetric difference $A \triangle B$ (pronounced "A delta B ") is defined as the set of all elements which belong to either $A$ or $B$ but not to both $A$ and $B$ :

$$
\begin{equation*}
A \triangle B:=(A \cup B) \backslash(A \cap B) \tag{2.3}
\end{equation*}
$$

Definition 2.6 (Universal set). Usually there always is a big set $\Omega$ that contains everything we are interested in and we then deal with all kinds of subsets $A \subseteq \Omega$. Such a set is called a "universal" set.

For example, in this document, we often deal with real numbers and our universal set will then be $\mathbb{R}$.
If there is a universal set, it makes perfect sense to talk about the complement of a set:
Definition 2.7 (Complement of a set). The complement of a set $A$ consists of all elements of $\Omega$ which do not belong to $A$. We write $A^{\complement}$. or $\complement A$ In other words:

$$
\begin{equation*}
A^{\complement}:=\complement A:=\Omega \backslash A=\{\omega \in \Omega: x \notin A\} \tag{2.4}
\end{equation*}
$$

$$
A \backslash B:
$$

$A \triangle B:$
Univeral set: $A^{\complement}$


Figure 2.3: Difference, symmetric difference, universal set, complement

Remark 2.4 (Complement of empty, all). Note that for any kind of universal set $\Omega$ it is true that

$$
\begin{equation*}
\Omega^{\complement}=\emptyset, \quad \emptyset^{\complement}=\Omega \tag{2.5}
\end{equation*}
$$

Example 2.2 (Complement of a set relative to the unit interval). Assume we are exclusively dealing with the unit interval, i.e., $\Omega=[0,1]=\{x \in \mathbb{R}: 0 \leqq x \leqq 1\} .{ }^{4}$ Let $a \in[0,1]$ and $\delta>0$ and

$$
\begin{equation*}
A=\{x \in[0,1]: a-\delta<x<a+\delta\} \tag{2.6}
\end{equation*}
$$

[^3]the $\delta$-neighborhood ${ }^{5}$ of $a$ (with respect to $[0,1]$ because numbers outside the unit interval are not considered part of our universe). Then the complement of $A$ is
$$
A^{\complement}=\{x \in[0,1]: x \leqq a-\delta \text { or } x \geqq a+\delta\} .
$$

Draw some Venn diagrams to visualize the following formulas.
Proposition 2.1. Let $A, B, X$ be sets and assume $A \subseteq X$. Then

$$
\begin{array}{rlrl}
A \cup \emptyset & =A ; & & A \cap \emptyset=\emptyset \\
A \cup \Omega & =\Omega ; & A \cap \Omega=A \\
A \cup A^{\complement} & =\Omega ; & A \cap A^{\complement}=\emptyset \\
A \triangle B & =(A \backslash B) \uplus(B \backslash A) \\
A \backslash A & =\emptyset & \\
A \triangle \emptyset & =A ; & A \triangle A=\emptyset \\
X \triangle A & =X \backslash A & \\
A \cup B & =(A \triangle B) \uplus(A \cap B) \\
A \cap B & =(A \cup B) \backslash(A \triangle B) \\
A \triangle B & =\emptyset \text { if and only if } B=A \tag{2.7j}
\end{array}
$$

The proof is left as exercise 2.1.
Definition 2.8 (Power set). The power set

$$
2^{\Omega}:=\mathfrak{P}(\Omega):=\{A: A \subseteq \Omega\}
$$

of a set $\Omega$ is the set of all its subsets.
Remark 2.5. Note that $\emptyset \in 2^{\Omega}$ for any set $\Omega$, even if $\Omega=\emptyset: 2^{\emptyset}=\{\emptyset\}$. It follows that the power set of the empty set is not empty.

A lot more will be said about sets once families are defined.

### 2.2 Numbers

Remark 2.6 (Classification of numbers). ${ }^{6}$
We call numbers without decimal points such as $3,-29,0,3000000,3 \cdot 10^{6},-1, \ldots$ integers and we write $\mathbb{Z}$ for the set of all integers.

Numbers in the set $\mathbb{N}=\{1,2,3, \ldots\}$ of all strictly positive integers are called natural numbers.
A number that is an integer or can be written as a fraction is called a rational number and we write $\mathbb{Q}$ for the set of all rational numbers. Examples of rational numbers are

$$
\frac{3}{4},-0.75,-\frac{1}{3}, . \overline{3}, \frac{13}{4},-5,2.99 \overline{9},-37 \frac{2}{7} .
$$

[^4]The bar on top of the rightmost part of a decimal such as ". $\overline{3}$ " means that this part should be repeated over and over again, e.g., $. \overline{3}=0.33333333333 \ldots$ and $1.234 \overline{567}=1.234567567567 \ldots$.

Note that a mathematician does not care whether a rational number is written as a fraction " $\frac{\text { numerator }}{\text { denominator }}$ " or as a decimal. The following all are representations of one third:

$$
\begin{equation*}
0 . \overline{3}=. \overline{3}=0.33333333333 \ldots=\frac{1}{3}=\frac{-1}{-3}=\frac{2}{6} \tag{2.8}
\end{equation*}
$$

and here are several equivalent ways of expressing the number minus four:

$$
\begin{equation*}
-4=-4.000=-3 . \overline{9}=-\frac{12}{3}=\frac{4}{-1}=\frac{-4}{1}=\frac{12}{-3}=-\frac{400}{100} . \tag{2.9}
\end{equation*}
$$

We call the barred portion of the decimal digits the period of the number and we also talk about repeating decimals. The number of digits in the barred portion is called the period length. This period length can be bigger than one. For example, the number $1.234 \overline{567}$ from above has period length 3 and the number $0.1 \overline{45}$ has period length 2 .

If $q \in \mathbb{Q}$ then there are unique integers $n$ and $d$ such that $q=\frac{n}{d}$ and
a. $d \in \mathbb{N}$,
b. $\quad d$ is minimal: there are no numbers $n^{\prime} \in \mathbb{Z}$ and $d^{\prime} \in \mathbb{N}$ such that $q=\frac{n^{\prime}}{d^{\prime}}$ and $d^{\prime}<d$.

We say that this choice of $n$ and $d$ is a representation of $q$ in lowest terms or that $q$ is then written in lowest terms. For example, the representation of.$\overline{3}$ in lowest terms is $\frac{1}{3}$ and the representation of -4 in lowest terms is $\frac{-4}{1}$.

Note that if $q \in \mathbb{Q}$ is strictly positive and if $\frac{d}{n}$ represents $q$ in lowest terms then $d \in \mathbb{N}$.
There are numbers which cannot be expressed as integers or fractions or numbers with a finite amount of decimals to the right of the decimal point. Examples are $\sqrt{2}$ and $\boldsymbol{\pi}$. Those irrational numbers (that is their official name) fill the gaps between the rational numbers. In fact, there is a simple way (but not easy to prove) of characterizing irrational numbers: Rational numbers are those that can be expressed with at most finitely many digits to the right of the decimal point, including repeating decimals. You can find the underlying theory and exact proofs in B/G ch.12. Irrational numbers must then be those with infinitely many decimal digits without any continually repeating patterns.

Example 2.3. To illustrate that repeating decimals are in fact rational numbers we convert $x=0.1 \overline{45}$ into a fraction:

$$
99 x=100 x-x=14.5 \overline{45}-0.1 \overline{45}=14.4
$$

It follows that $x=144 / 990$, and that is certainly a fraction.
Now we can finally give an informal definition of the most important kind of numbers: We call any kind of number, either rational or irrational, a real number and we write $\mathbb{R}$ for the set of all real numbers. It can be shown that there are a lot more irrational numbers than rational numbers, even though $\mathbb{Q}$ is a dense subset in $\mathbb{R}$ in the following sense: No matter how small an interval $(a, b)=\{x \in \mathbb{R}: a<x<b\}$ of real numbers you choose, it will contain infinitely many rational numbers.

Definition 2.9 (Types of numbers). We summarize what was said sofar about the classification of numbers:
$\mathbb{N}:=\{1,2,3, \ldots\}$ denotes the set of natural numbers.
$\mathbb{Z}:=\{0, \pm 1, \pm 2, \pm 3, \ldots\}$ denotes the set of all integers.
$\mathbb{Q}:=\{n / d: n \in \mathbb{Z}, d \in \mathbb{N}\}$ denotes the set of all rational numbers.
$\mathbb{R}:=\{$ all integers or decimal numbers with finitely or infinitely many decimal digits $\}$ denotes the set of all real numbers.
$\mathbb{R} \backslash \mathbb{Q}=$ \{all real numbers which cannot be written as fractions of integers \} denotes the set of all irrational numbers. There is no special symbol for irrational numbers. Example: $\sqrt{2}$ and $\pi$ are irrational.

The following are customary abbreviations of some often referenced sets of numbers:
$\mathbb{N}_{0}:=\mathbb{Z}_{+}:=\mathbb{Z}_{\geq 0}:=\{0,1,2,3, \ldots\}$ denotes the set of non-negative integers,
$\mathbb{R}_{+}:=\mathbb{R}_{\geq 0}:=\{x \in \mathbb{R}: x \geqq 0\}$ denotes the set of all non-negative real numbers,
$\mathbb{R}^{+}:=\mathbb{R}_{>0}:=\{x \in \mathbb{R}: x>0\}$ denotes the set of all positive real numbers,
$\mathbb{R}^{\star}:=\mathbb{R}_{\neq 0}:=\{x \in \mathbb{R}: x \neq 0\}$.
Definition 2.10 (Intervals of real numbers). We use the following notation for intervals of real numbers $a, b$ :
$[a, b]:=\{x \in \mathbb{R}: a \leqq x \leqq b\}$ is called the closed interval with endpoints $a$ and $b$.
$] a, b[:=\{x \in \mathbb{R}: a<x<b\}$ is called the open interval with endpoints $a$ and $b$.
$[a, b[:=\{x \in \mathbb{R}: a \leqq x<b\}$ and $] a, b]:=\{x \in \mathbb{R}: a<x \leqq b\}$ are called half-open intervals with endpoints $a$ and $b$.

We further define the following intervals of "infinite length":

$$
\begin{align*}
]-\infty, a] & :=\{x \in \mathbb{R}: x \leqq a\}, \quad]-\infty, a[:=\{x \in \mathbb{R}: x<a\},  \tag{2.10}\\
] a, \infty[ & :=\{x \in \mathbb{R}: x>a\}, \quad[a, \infty[:=\{x \in \mathbb{R}: x \geqq a\}, \quad[-\infty, \infty[:=\mathbb{R}
\end{align*}
$$

Finally we define $[a, b[:=] a, b[:=] a, b]:=\emptyset$ for $a \geqq b$ and $[a, b]:=\emptyset$ for $a>b$.
Definition 2.11 (Absolute value). For a real number $x$ we define its absolute value as

$$
|x|= \begin{cases}x & \text { if } x \geqq 0 \\ -x & \text { if } x<0\end{cases}
$$

Example 2.4. $|3|=3 ;|-3|=3 ;|-5.38|=5.38$.
Remark 2.7. For any real number $x$ we have

$$
\begin{equation*}
\sqrt{x^{2}}=|x| \tag{2.11}
\end{equation*}
$$

Assumption 2.1 (Square roots are always assumed non-negative). Remember that for any number $a$ it is true that

$$
a \cdot a=(-a)(-a)=a^{2} \quad \text { e.g., } \quad 2^{2}=(-2)^{2}=4
$$

or that, expressed in form of square roots, for any number $b \geqq 0$

$$
(+\sqrt{b})(+\sqrt{b})=(-\sqrt{b})(-\sqrt{b})=b
$$

We will always assume that " $\sqrt{b}$ " is the positive value unless the opposite is explicitly stated.

Example: $\sqrt{9}=+3$, not -3 .
Proposition 2.2 (The Triangle Inequality for real numbers). The following inequality is used all the time in mathematical analysis to show that the size of a certain expression is limited from above:

$$
\begin{equation*}
\text { Triangle Inequality : }|\boldsymbol{a}+\boldsymbol{b}| \leqq|\boldsymbol{a}|+|\boldsymbol{b}| \tag{2.12}
\end{equation*}
$$

This inequality is true for any two real numbers $a$ and $b$.
It is easy to prove this: just look separately at the three cases where both numbers are non-negative, both are negative or where one of each is positive and negative.

Proposition 2.3 (The Triangle Inequality for $n$ real numbers). The above inequality also holds true for more than two real numbers: Let $n \in \mathbb{N}$ such that $n \geqq 2$. Let $a_{1}, a_{2}, \ldots, a_{n} \in \mathbb{N}$. Then

$$
\begin{equation*}
\left|a_{1}+a_{2}+\ldots+a_{n}\right| \leqq\left|a_{1}\right|+\left|a_{2}\right|+\ldots+\left|a_{n}\right| \tag{2.13}
\end{equation*}
$$

The proof will be done by induction, a principle which is defined first:
Definition 2.12 (Principle of proof by mathematical induction). Actually, "definition" is a misnomer. This principle is a mathematical statement that follows from the structure of the natural numbers which have a starting point to the "left" (a smallest element 1 ) and then progress in the well understood sequence ${ }^{7}$

$$
2,3,4, \ldots, k-1, k, k+1, \ldots
$$

This is the principle: Let us assume that we know that some statement can be proved to be true in the following two situations:
A. Base case. The statement is true for some (small) $k_{0}$; usually that means $k_{0}=0$ or $k_{0}=1$
B. Induction Step. We prove the following for all $k \in \mathbb{N}_{0}$ such that $k \geqq k_{0}$ : if the property is true for $k$ ("Induction Assumption") then it will also be true for $k+1$
C. Conclusion: Then the property is true for any $k \in \mathbb{N}_{0}$ such that $k \geqq k_{0}$.

Either you have been explained this principle before and say "Oh, that - what's the big deal?" or you will be mighty confused. So let me explain how it works by walking you through the proof of the triangle inequality for $n$ real numbers (2.13).

## Proof of the triangle inequality for $n$ real numbers:

A. For $k_{0}=2$, inequality 2.13 was already shown (see 2.12), so we found a $k_{0}$ for which the property is true.

[^5]B. Let us assume that 2.13 is true for some $k \geqq 2$. We now must prove the inequality for $k+1$ numbers $a_{1}, a_{2}, \ldots, a_{k}, a_{k+1} \in \mathbb{N}$ : We abbreviate
$$
A:=a_{1}+a_{2}+\ldots+a_{k} ; \quad B:=\left|a_{1}\right|+\left|a_{2}\right|+\ldots+\left|a_{k}\right|
$$
then our induction assumption for $k$ numbers is that $|A| \leqq B$. We know the triangle inequality is valid for the two variables $A$ and $a_{k+1}$ and it follows that $\left|A+a_{k+1}\right| \leqq|A|+\left|a_{k+1}\right|$. Look at both of those inequalities together and you have
\[

$$
\begin{equation*}
\left|A+a_{k+1}\right| \leqq|A|+\left|a_{k+1}\right| \leqq B+\left|a_{k+1}\right| \tag{2.14}
\end{equation*}
$$

\]

In other words,

$$
\begin{equation*}
\left|\left(a_{1}+a_{2}+\ldots+a_{k}\right)+a_{k+1}\right| \leqq B+\left|a_{k+1}\right|=\left(\left|a_{1}\right|+\left|a_{2}\right|+\ldots+\left|a_{k}\right|\right)+\left|a_{k+1}\right| \tag{2.15}
\end{equation*}
$$

and this is (2.13) for $k+1$ rather than $k$ numbers: We have shown the validity of the triangle inequality for $k+1$ items under the assumption that it is valid for $k$ items. It follows from the induction principle that the inequality is valid for any $k \geqq k_{0}=2$.

To summarize what we did in all of part B: We were able to show the validity of the triangle inequality for $k+1$ numbers under the assumption that it was valid for $k$ numbers.

Remark 2.8 (Why induction works). But how can we from all of the above conclude that the triangle inequality works for all $n \in \mathbb{N}$ such that $n \geqq k_{0}=2$ ? That's much simpler to demonstrate than what we just did.

Step 1: We know that it is true for $k_{0}=2$ because that was actually proved in A.
Step 2: $\quad$ But according to B, if it is true for $k_{0}$, it is also true for the successor $k_{0}+1=3$.
Step 3: $\quad$ But according to B , if it is true for $k_{0}+1$, it is also true for the successor $\left(k_{0}+1\right)+1=4$.
Step 4: $\quad$ But according to B , if it is true for $k_{0}+2$, it is also true for the successor $\left(k_{0}+2\right)+1=5$.

Step 53, 920: But according to B, if it is true for $k_{0}+53,918$, it is also true for the successor $\left(k_{0}+53,918\right)+1=53,921$.

And now you understand why it is true for any natural number $n \geqq k_{0}$.

### 2.2.1 Rings and Algebras of Sets ( $\star$ )

This optional chapter is out of place: It should be a subchapter of ch.2.1 (Sets and Basic Set Operations).

## Note that this chapter is starred, hence optional.

You will benefit from examining the proof of prop. 2.4 on p. 21 and learn how to split a proof which involves 3 or 4 sets can be split into easily dealt with cases.

Definition 2.13 (Rings and Algebras of Sets). A subset $\mathscr{R}$ of $2^{\Omega}$ (a set of sets!) is called a ring of sets if it is closed with respect to the operations " $\cup$ " and " $\backslash$ ", i.e.,

$$
\begin{equation*}
R_{1} \cup R_{2} \in \mathscr{R} \text { and } R_{1} \backslash R_{2} \in \mathscr{R} \quad \text { whenever } R_{1}, R_{2} \in \mathscr{R} \tag{2.16}
\end{equation*}
$$

A subset $\mathscr{A}$ of $2^{\Omega}$ is called an algebra of sets if $\Omega \in \mathscr{A}$ and $\mathscr{A}$ is a ring of sets.
Proposition 2.4. 1. Let $\mathscr{R}$ be a ring of sets and $A, B \in \mathscr{R}$. Then $A \triangle B \in \mathscr{R}$ and $A \cap B \in \mathscr{R}$.
2. Let $\mathscr{R}$ be a ring of sets and $A, B, C \in \mathscr{R}$. Then
a. $(A \triangle B) \triangle C=A \triangle(B \triangle C) \quad$ (associativity of $\triangle$ )
b. $A \triangle \emptyset=\emptyset \triangle A=A \quad$ (neutral element $\emptyset$ for $\triangle$ )
c. $A \triangle A=\emptyset \quad$ (inverse element $\emptyset$ for $\triangle)^{8}$
d. $A \triangle B=B \triangle A \quad$ (commutativity of $\triangle$ )

Further we have the following for the intersection operation:
e. $\quad(A \cap B) \cap C=A \cap(B \cap C) \quad$ (associativity of $\cap$ )
f. $A \cap \Omega=\Omega \cap A=A \quad$ (neutral element $\Omega$ for $\cap$ )
g. $A \cap B=B \cap A \quad$ (commutativity of $\cap$ )

And we have the following interrelationship between $\triangle$ and $\cap$ :
h. $A \cap(B \triangle C)=(A \cap B) \triangle(A \cap C)$ (distributivity)

## Solution:

PROOF of $\boldsymbol{a}$. The proof given here is based on with set membership. It is very tedious and there is a much more elegant proof but it requires knowledge of indicator functions ${ }^{9}$ and of base 2 modular arithmetic (see, e.g., [1] B/G (Beck/Geoghegan) ch.6.2).

By definition $x \in U \triangle V$ if and only if either $x \in U$ or $x \in V$,i.e., (either) $[x \in U$ and $x \notin V]$ or $[x \in V$ and $x \notin U]$

Hence $x \in(A \triangle B) \triangle C$ means either $x \in(A \triangle B)$ or $x \in C$, i.e., either $[x \in A, x \notin B$ or $x \in B, x \notin A]$ or $x \in C$, i.e., we have one of the following four combinations:
a. $\quad x \in A \quad x \notin B \quad x \notin C$
b. $x \notin A \quad x \in B \quad x \notin C$
c. $x \in A \quad x \in B \quad x \in C$
d. $x \notin A \quad x \notin B \quad x \in C$
and $x \in A \triangle(B \triangle C)$ means either $x \in A$ or $x \in(B \triangle C)$, i.e., either $x \in A$ or $[x \in B, x \notin C$ or $x \in C, x \notin B]$, i.e., we have one of the following four combinations:

$$
\begin{array}{llll}
\text { 1. } & x \in A & x \in B & x \in C \\
\text { 2. } & x \in A & x \notin B & x \notin C \\
\text { 3. } & x \notin A & x \in B & x \notin C \\
\text { 4. } & x \notin A & x \notin B & x \in C
\end{array}
$$

We have a perfect match $\boldsymbol{a} \leftrightarrow \mathbf{2}, \boldsymbol{b} \leftrightarrow \mathbf{3}, \boldsymbol{c} \leftrightarrow \mathbf{1}, \boldsymbol{d} \leftrightarrow \mathbf{4}$. and this completes the proof of $\boldsymbol{a}$.

[^6]Remark 2.9 (Rings of sets as rings). A set $G=(G, \oplus)$ with a "binary operation" $g_{3}=g_{1} \oplus g_{2}$ - i.e., an operation $\oplus$ which assigns to any two $g_{1}, g_{2} \in G$ a third element $g_{3} \in G$ - such that $\oplus$ satisfies a-d above with $\triangle$ playing the role of $\oplus$ is called an abelian or commutative group (see ch. 13 on p.269.)

A set $R=(R, \oplus, \odot)$ with two binary operations $g_{3}=g_{1} \oplus g_{2}$ and $g_{4}=g_{1} \odot g_{2}$ which satisfies a $\mathbf{h}$ above with $\triangle$ playing the role of $\oplus$ and $\cap$ playing the role of $\odot$ is called a commutative ring with multiplicative unit. $\oplus$ is customarily referred to as "addition" and $\odot$ is customarily referred to as "multiplication". This explains the name "ring of sets" for $\mathscr{R}=(\mathscr{R}, \triangle, \cap)$.

We note that rings of sets satisfy axioms $1.1-1.4$ but not axiom 1.5 (the cancellation axiom) of [1] B/G (Beck/Geoghegan) ch.1.1. This explains the name "ring of sets".

The name "algebra of sets" for a ring of sets which contains $\Omega$ stems from the fact that such systems of subsets of $\Omega$ are "boolean algebras".

### 2.3 Exercises for Ch. 2

### 2.3.1 Exercises for Sets

Exercise 2.1. Prove the set identities of prop.2.1.
Exercise 2.2. Let $X=\{x, y,\{x\},\{x, y\}\}$. True or false?
a. $\{x\} \in X$
b. $\{x\} \subseteq X$
c. $\{\{x\}\} \in X$
d. $\{\{x\}\} \subseteq X$
e. $y \in X$
f. $y \subseteq X$
g. $\{y\} \in X$
h. $\{y\} \subseteq X$

For the subsequent exercises refer to example 4.4 for the preliminary definition of the size $|A|$ of a set $A$ and to def.4.1 (Cartesian Product of two sets) for the definition of Cartesian product. You find both in ch.4.1 (Cartesian products and relations) on p. 67

Exercise 2.3. Find the size of each of the following sets:
a. $A=\{x, y,\{x\},\{x, y\}\}$
b. $B=\{1,\{0\},\{1\}\}$
c. $C=\{u, v, v, v, u\}$
d. $D=\{3 z-10: z \in \mathbb{Z}\}$
e. $E=\{\sin (k \pi / 2): k \in \mathbb{Z}\}$
f. $F=\{\pi x: x \in \mathbb{R}\}$

Exercise 2.4. Let $X=\{x, y,\{x\},\{x, y\}\}$ and $Y=\{x,\{y\}\}$. True or false?
a. $x \in X \cap Y$
b. $\{y\} \in X \cap Y$
c. $x \in X \cup Y$
d. $\{y\} \in X \cup Y$
e. $x \in X \backslash Y$
f. $\{y\} \in X \backslash Y$
g. $x \in X \Delta Y$
h. $\{y\} \in X \Delta Y$

Exercise 2.5. Let $X=\{1,2,3,4\}$ and let $Y=\{x, y\}$.
a. What is $X \times Y$ ?
c. What is $|X \times Y|$ ?
e. Is $(x, 3) \in X \times Y$ ?
g. Is $3 \cdot x \in X \times Y$ ?
b. What is $Y \times X$ ?
d. What is $|X \times Y|$ ?
f. Is $(x, 3) \in Y \times X$ ?
h. Is $2 \cdot y \in Y \times X$ ?

Exercise 2.6. Let $X=\{2\}$. What is $2^{\left(2^{X}\right)}$ ?

## 3 Logic (Skim this!)

This chapter uses material presented in ch. 2 (Logic) and ch. 3 (Methods of Proofs) of [4] Bryant, Kirby Course Notes for MAD 2104.

### 3.1 Prologue: Notation for Functions

The material on functions presented in this section will be discussed again and in greater detail in chapter 4 (unctions and relations) on p.67. It is presented at the beginning of this chapter about logic because the definition of a function is needed to properly discuss statement functions (see ch.3.2 on p.26).

Note 3.1 (Motivation for a good function definition). When discussing logic we deal with statement functions (predicates) (see def.3.4 on p.26) and we are in the same predicament as when discussing some run of the mill functions known from calculus such as $f_{1}(x)=\sqrt{x}$ and $f_{2}(x, y)=\ln (x-y)$ : Sometimes $f_{1}(x)$ means the entire graph, i.e., the entire collection of pairs $(x, \sqrt{x})$ and sometimes it just refers to the function value $\sqrt{x}$ for a "fixed but arbitrary" number $x$. In case of the function $f_{2}(x)$ : Sometimes $f_{2}(x, y)$ means the entire graph, i.e., the entire collection of pairs $((x, y), \ln (x-y))$ and sometimes it just refers to the function value $\ln (x-y)$ for a pair of "fixed but arbitrary" numbers $(x, y)$.

This issue is addressed in the material of ch. 4.2 on p .70 which precedes the mathematically precise definition of a function (def.4.6 on p.72). You are encouraged to look at it once you have read the remainder of this short section as ch.4.2 contains everything you see here.

To get to a usable definition of a function there are several things to consider. In the following $f_{1}(x)$ and $f_{2}(x, y)$ again denote the functions $f_{1}(x)=\sqrt{x}$ and $f_{2}(x, y)=\ln (x-y)$.
a. The source of all allowable arguments ( $x$-values in case of $f_{1}(x)$ and $(x, y)$-values in case of $f_{2}(x, y)$ ) will be called the domain of the function. The domain is explicitly specified as part of a function definition and it may be chosen for whatever reason to be only a subset of all arguments for which the function value is a valid expression. In case of the function $f_{1}(x)$ this means that the domain must be restricted to a subset of the interval $[0, \infty[$ because the square root of a negative number cannot be taken. In case of the function $f_{2}(x, y)$ this means that the domain must be restricted to a subset of $\{(x, y)$ : $x, y \in \mathbb{R}$ and $x-y>0\}$ because logarithms are only defined for strictly positive numbers.
b. The set to which all possible function values belong will be called the codomain of the function. As is the case for the domain, the codomain also is explicitly specified as part of a function definition. It may be chosen as any superset of the set of all function values for which the argument belongs to the domain of the function. In case of the function $f_{1}(x)$ this means that we are OK if the codomain is a superset of the interval $[0, \infty[$. Such a set is big enough because square roots are never negative. It is OK to specify the interval $]-3.5, \infty$ or even the set $\mathbb{R}$ of all real numbers as the codomain. In case of the function $f_{2}(x, y)$ this means that we are OK if the codomain contains $\mathbb{R}$. Not that it would make a lot of sense but the set $\mathbb{R} \cup\{$ all inhabitants of Chicago $\}$ also is an acceptable choice for the codomain.
c. A function $y=f(x)$ is not necessarily something that maps (assigns) numbers or pairs of numbers to numbers but domain and codomain can be a very different kind of animal. In this chapter on logic you will learn about statement functions $A(x)$ which assign arguments $x$ from an arbitrary set $\mathscr{U}$, called the universe of discourse, to statements $A(x)$, i.e., sentences that are either true or false.
d. Considering all that was said so far you can now think of the graph of a function $f(x)$ with domain $D$ and codomain $C$ (see the beginning of this chapter) as the set $\Gamma_{f}:=$ $\{(x, f(x)): x \in D\}$. Alternatively you can characterize this function by the assignment rule which specifies how the function value $f(x)$ depends on any given argument $x \in D$. We write " $x \mapsto f(x)$ to indicate this. You can also write instead $f(x)=$ whatever the actual function value will be. This is possible if you do not write about functions in general but about specific functions such as $f_{1}(x)=\sqrt{x}$ and $f_{2}(x, y)=\ln (x-y)$. We further write " $f: C \rightarrow D$ " as a short way of saying that the function $f(x)$ has domain $C$ and codomain $D$.
In case of the function $f_{1}(x)=\sqrt{x}$ for which we choose the interval $X:=[2.5,7]$ as the domain (small enough because $X \subseteq[0, \infty[$ ) and $Y:=] 1,3[$ as the codomain (big enough because $1<\sqrt{x}<3$ for any $x \in X$ ) we specify this function as

$$
\text { either } f_{1}: X \rightarrow Y, \quad x \mapsto \sqrt{x} \quad \text { or } f_{1}: X \rightarrow Y, \quad f(x)=\sqrt{x}
$$

Let us choose $U:=\{(x, y): x, y \in \mathbb{R} 1 \leqq x \leqq 10$ and $y<-2\}$ as the domain and $V:=\left[0, \infty\left[\right.\right.$ as the codomain for $f_{2}(x, y)=\ln (x-y)$. These choices are OK because $x-y \geqq 1$ for any $(x, y) \in U$ and hence $\ln (x-y) \geqq 0$, i.e., $f_{2}(x, y) \in V$ for all $(x, y \in U$. We specify this function as
either $f_{2}: U \rightarrow V, \quad(x, y) \mapsto \ln (x-y) \quad$ or $f_{2}: U \rightarrow V, \quad f(x, y)=\ln (x-y)$.

We incorporate the above into the following preliminary definition.
Definition 3.1 (Preliminary definition of a function). A function $f$ consists of two nonempty sets $X$ and $Y$ and an assignment rule $x \mapsto f(x)$ which assigns any $x \in X$ uniquely to some $y \in Y$. We write $f(x)$ for this assigned value and call it the function value of the argument $x . X$ is called the domain and $Y$ is called the codomain of $f$. We write

$$
\begin{equation*}
f: X \rightarrow Y, \quad x \mapsto f(x) \tag{3.1}
\end{equation*}
$$

We read " $a \mapsto b$ " as " $a$ is assigned to $b$ " or " $a$ maps to $b$ " and refer to $\mapsto$ as the maps to operator or assignment operator. The graph of such a function is the collection of pairs

$$
\begin{equation*}
\Gamma_{f}:=\{(x, f(x)): x \in X\} . \tag{3.2}
\end{equation*}
$$

Remark 3.1. The name given to the argument variable is irrelevant. Let $f_{1}, f_{2}, X, Y, U, V$ be as defined in $\mathbf{d}$ of note 3.1. The function

$$
g_{1}: X \rightarrow Y, \quad p \mapsto \sqrt{p}
$$

is identical to the function $f_{1}$. The function

$$
g_{2}: U \rightarrow V, \quad(t, s) \mapsto \ln (t-s)
$$

is identical to the function $f_{2}$ and so is the function

$$
g_{3}: U \rightarrow V, \quad(s, t) \mapsto \ln (s-t) .
$$

The last example tells you that you can swap function names as long as you do it consistently in all places.

Note 3.2 (Textual variables). It was mentioned in cabove that the input variables and function values need not necessarily numbers but they can also be textual. For example, the domain of a function may consist of the first names of certain persons.

A note on textual variables: If the variable is the last name of the person James Joice and valid input for the function $F: p \mapsto$ "Each morning $p$ writes two pages.") then we write interchangeably Joyce or 'Joyce'. Quotes are generally avoided unless they add clarity.

In the above example "Each morning 'Joyce' writes two pages." emphasizes that Joyce is the replacement of a parameter whereas $F$ ('Joyce') does not seem to improve the simpler notation $F$ (Joyce) and you will most likely see the expression $F$ (Joyce) = "Each morning 'Joyce' writes two pages."

We also need the definition of a cartesian product. ${ }^{10}$
Definition 3.2 (Preliminary definition: cartesian product). Let $X$ and $Y$ be two sets The set

$$
\begin{equation*}
X \times Y:=\{(x, y): x \in X, y \in Y\} \tag{3.3}
\end{equation*}
$$

is called the cartesian product of $X$ and $Y$.
Note that the order is important: $(x, y)$ and $(y, x)$ are different unless $x=y$.
We write $X^{2}$ as an abbreviation for $X \times X$.
This definition generalizes to more than two sets as follows: Let $X_{1}, X_{2}, \ldots, X_{n}$ be sets. The set

$$
\begin{equation*}
X_{1} \times X_{2} \cdots \times X_{n}:=\left\{\left(x_{1}, x_{2}, \ldots, x_{n}\right): x_{j} \in X_{j} \text { for each } j=1,2, \ldots n\right\} \tag{3.4}
\end{equation*}
$$

is called the cartesian product of $X_{1}, X_{2}, \ldots, X_{n}$.
We write $X^{n}$ as an abbreviation for $X \times X \times \cdots \times X$.
Example 3.1. The graph $\Gamma_{f}$ of a function with domain $X$ and codomain $Y$ (see def.3.2) is a subset of the cartesian product $X \times Y$.

Example 3.2. The domains given in a and $\mathbf{d}$ of note 3.1 are subsets of the cartesian product $\mathbb{R}^{2}=$ $\mathbb{R} \times \mathbb{R}=\{(x, y): x, y \in \mathbb{R}\}$.

[^7]
### 3.2 Statements and Statement Functions

Definition 3.3 (Statements). A statement ${ }^{11}$ is a sentence or collection of sentences that is either true or false. We write T or true for "true" and F or false for "false" and we refer to those constants as truth values

Example 3.3. The following are examples of statements:
a. "Dogs are mammals" (a true statement);
b. "Roses are mammals. 7 is a number." This is a false statement which also could have been written as a single sentence: "Roses are mammals and 7 is a number";
c. "I own 5 houses" (a statement because this sentence is either true or false depending on whether I told the truth or I lied);
d. "The sum of any two even integers is even" (a true statement);
e. "The sum of any two even integers is even and Roses are mammals" (a false statement);
f. "Either the sum of any two even integers is even or Roses are mammals" (a true statement).

Example 3.4. The following are not statements:
a. "Who is invited for dinner?"
b. " $2 x=27$ " (the variable $x$ must be bound (specified) to determine whether this sentence is true or false: It is true for $x=13.5$ and it is false for $x=33$ )
c. " $x^{2}+y^{2}=34$ " (both variables $x$ and $y$ must be bound to determine whether this sentence is true or false It is true for $x=5$ and $y=3$ and it is false for $x=7.8$ and $y=2$ )
d. "Stop bothering me!"

For the remainder of the entire chapter on logic we define

$$
\begin{equation*}
\mathscr{S}:=\text { the set of all statements } \tag{3.5}
\end{equation*}
$$

$\mathscr{S}$ will appear as the codomain of statement functions.
Be sure to understand the material of ch.3.1 (Prologue: Notation for Functions) on p.23) before continuing.

Definition 3.4 (Statement functions (predicates)). We need to discuss some preliminaries before arriving at the definition of a statement function. Let $A$ be a sentence or collection of sentences which contains one or more variables (placeholders) such that, if each of those variables is assigned a specific value, it is either true or false, i.e., it is an element of the set $\mathscr{S}$ of all statements. If $A$ contains $n$ variables $x_{1}, x_{2}, \ldots, x_{n}$ and if they are bound, i.e., assigned to the specific values $x_{1}=x_{10}, x_{2}=x_{20}, \ldots, x_{n}=x_{n 0}$, we write $A\left(x_{10}, x_{20}, \ldots, x_{n 0}\right)$ for the resulting statement.

To illustrate this let $A:=$ " $x$ is green and $y$ and $z$ like each other". If we know the specific values for the variables $x, y, z$ then this sentence will be true or false. For example $A$ (this lime, Tim, Fred) is true or false depending on whether Tim and Fred do or do not like each other.

There are restrictions for the choice of $x_{1}=x_{10}, x_{2}=x_{20}, \ldots, x_{n}=x_{n 0}$ : Associated with each variable $x_{j}$ in $A$ is a set $\mathscr{U}_{j}$ which we call the universe of discourse, in short, $\mathbf{U o D}$, for the $j$ th variable in $A$. Each value $x_{j_{0}}(j=1,2, \ldots n)$ must be chosen in such a way that $x_{j_{0}} \in \mathscr{U}_{j}$. If this is

[^8]not the case then the expression $A\left(x_{10}, x_{20}, \ldots, x_{n 0}\right)$ is called inadmissible and we refuse to deal with it.

What was said can be rephrased as follows: We have an assignment $\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto$ $A\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ which results in a statement, i.e., an element of $\mathscr{S}$ (see (3.5)) just as long as $x_{j_{0}} \in \mathscr{U}_{j}$. In other words we have a function

$$
\begin{equation*}
A: \mathscr{U}_{1} \times \mathscr{U}_{2} \times \cdots \times \mathscr{U}_{n} \rightarrow \mathscr{S}, \quad\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto A\left(x_{1}, x_{2}, \ldots, x_{n}\right) \tag{3.6}
\end{equation*}
$$

in the sense of def. 3.1. with the cartesian product of the UoDs for $x_{1}, \ldots, x_{n}$ as domain and $\mathscr{S}$ as codomain. We call such a function a statement function ${ }^{12}$ or predicate.

Note 3.3 (Relaxed notation for statement functions). You should remember that a statement function is a function in the sense of def.3.1 but we we will often use the simpler notation $A:=$ "some text that contains the placeholders $x_{1}, x_{2}, \ldots, x_{n}$ and evaluates to true or false once all $x_{j}$ are bound"
together with the specification of each $\mathrm{UoD} \mathscr{U}_{j}$ rather than the formal notation

$$
A: \mathscr{U}_{1} \times \mathscr{U}_{2} \times \cdots \times \mathscr{U}_{n} \rightarrow \mathscr{S}, \quad\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto A\left(x_{1}, x_{2}, \ldots, x_{n}\right) .
$$

If $A$ contains two or more variables then the formal notation has an advantage. There is no doubt when looking at an evaluation such as $A(5.5,7,-3,8)$ which placeholder in the string corresponds to 5.5 , which one corresponds to 7 etc. When employing the relaxed notation then we decide this according to the following

Left to right rule for statement functions: If the string $A$ contains $n$ different place holders then the expression $A\left(x_{10}, x_{20}, \ldots x_{n 0}\right)$ implies the following: If the name of the first (leftmost) place holder in $A$ is $x$ then each occurrence of $x$ is bound to the value $x_{10}$. If the name of the first of the remaining place holders in $A$ is $y$ then each occurrence of $y$ is bound to the value $x_{20}, \ldots$. After $n-1$ steps the remaining placeholders all have the same name, say $z$ and each occurrence of $z$ is bound to the value $x_{n 0}$. If there is any confusion about what is first, what is second, ... then this will be indicated when $A$ is specified or when its variables are bound for the first time.

Example 3.5. In def.3.4 $A=$ " $x$ is green and $y$ and $z$ like each other" was used to illustrate the concept of a statement function. We never showed how to write the actual statement function. We must decide the UoDs for $x, y, z$ and we define them as follows.

UoD for $x: \mathscr{U}_{x}:=$ all plants and animals in the U.S.,
UoDs for $y$ and $z: \quad \mathscr{U}_{y}:=\mathscr{U}_{z}:=$ all BU majors in actuarial science.
a. Here is the formal definition: Let $A$ be the statement function

$$
A: \mathscr{U}_{x} \times \mathscr{U}_{y} \times \mathscr{U}_{z} \rightarrow \mathscr{S}, \quad(x, y, z) \mapsto A(x, y, z):=" x \text { is green and } y \text { and } z \text { like each other" }
$$

b. Here is the relaxed definition: Let $A$ be the statement function
$A:=$ " $x$ is green and $y$ and $z$ like each other" with UoDs $\mathscr{U}_{x}$ for $x, \mathscr{U}_{y}$ for $y$ and $\mathscr{U}_{z}$ for $z$.

[^9]The example above and all those below for statement functions of more than a single variable employ the left to right rule.

## Adhering to the left to right rule is not a big deal because of the following convention:

We will restrict ourselves in this document from now on to statement functions of one or two variables.

Example 3.6. Let $A(t)=$ " $t-4.7$ is an integer". Then $A: \mathbb{R} \rightarrow \mathscr{S}, x \mapsto A(x)$ is a one parameter statement function with $\operatorname{UoD} \mathbb{R}$ and $x$ as the variable. Note that it is immaterial that we wrote $t$ in the equation and $x$ in the " $\mapsto$ " expression because with deal with a dummy variable and we have employed its name consistently in both cases. We have
a. $\quad A($ Honda $)=$ "Honda' -4.7 is an integer" is inadmissible because a car brand is not part of our universe of discourse.
b. If $u_{0} \in \mathscr{U}$ then $A\left(u_{0}\right)=$ " $u_{0}-4.7$ is an integer" is a statement which evaluates to true or false depending on that fixed but unknown value of $u_{0}$.
c. If $n \in \mathscr{U}$ then $A(n)$ is the statement(!) " $n-4.7$ is an integer". It does not matter that this expression looks exactly like the original $A$ : The expression $A(n)$ implies that the parameter inside the sentence collection $A$ which happens to be named " $n$ " has been bound to a fixed (but unspecified) value also denoted by $n$.

Example 3.7. Let $B(x, y):=" x^{2}-y+2=11^{\prime \prime}$. Then $\left.B: \mathbb{R} \times\right] 1,100[\rightarrow \mathscr{S},(x, y) \mapsto B(x, y)$ is a two parameter statement function with $\operatorname{UoD} \mathbb{R}$ for $x$ and UoD$] 1,100[$ for the variable $y$. Then
a. $\quad B(4,-2)=$ " $4^{2}-(-2)+2=11^{\prime \prime}$ (a false statement) because $x$ is the leftmost item in $B$.
b. $\quad B(z, 10)=$ " $z^{2}-10+2=11$ " (true or false depending on $z$ ).
c. BE CAREFUL: If $x, y \in \mathbb{R}$ then $B(y, x)=" y^{2}-x+2=11^{\prime}$ and NOT " $x^{2}-y+2=11$ " because the "evaluate left to right" rule matters, not any similarity or even coincidence between the symbols inside the sentence collection and in the evaluation $B(\cdot, \cdot)$

Example 3.8. The following are predicates:
a. $\quad P:=" 2 x=27^{\prime \prime}$ (see example 3.4.b), UoD $\mathscr{U}:=\{x \in \mathbb{R}: x>10\}$
b. $\quad Q:=" x^{2}+y^{2}=34$ " (example 3.4.c), $\mathrm{UoD} \mathscr{V}:=\{(x, y): x, y \in \mathbb{R}$ and $x<y\}$
c. $\quad R:=" x^{2}+y^{2}=34$ and $x y>100$ ", UoDs are $\mathscr{W}_{x}:=\mathscr{W}_{y}:=[-50,25]$.

Note the following for $\mathbf{c}$ : $R(-30,20)$ evaluates to a false statement because $(-30) \cdot 20>100$ is false. $R(30,20)$ does not evaluate to any kind of statement: It is an inadmissible expression because $30 \notin \mathscr{W}_{x}$.
d. The sentence "Stop bothering $x$ !" is not a statement function because this imperative will not be true or false even if $x$ is bound to a specific value.
Example 3.9. Let $B:=" x+7=16$ and $d$ is a $\operatorname{dog}^{\prime \prime}$. Let $\mathscr{U}_{x}:=\mathbb{N}$ and $\mathscr{U}_{d}:=\{d:$ $d$ is a vegetable or animal $\}$.
$B$ becomes a statement function of two variables $x$ and $d$ if we specify that the $\operatorname{UoD}$ for $x$ is $\mathscr{U}_{x}$ and the UoD for $d$ is $\mathscr{U}_{d}$

Assume for the following that Robby is an animal.
a. $\quad B(9$, Robby $)$ is the statement " $9+7=16$ and Robby is a $\operatorname{dog}^{\prime \prime}$. It is true in case Robby is a dog and false in case Robby is not a dog.
b. $B(20$, Robby ) is the statement " $20+7=16$ and Robby is a dog" which is false regardless of what Robby might be because $20+7=16$ by itself is false.
c. $\quad B(d, F)$ is the statement " $d+7=16$ and $F$ is a dog": which is true or false depending on the fixed but unspecified values of $d$ and $F$. Note that $d$ corresponds to the leftmost variable $x$ inside $B$ and not to the second variable $d$ !
d. $B(x)$ is not a valid expression as we do not allow "partial evaluation" of a predicate. ${ }^{13}$

### 3.3 Logic Operations and their Truth Tables

We now resume our discussion of statements.

### 3.3.1 Overview of Logical Operators

Statements can be connected with logical operators, also called connectives, to form another statement, i.e., something that is either true or false.

Here is an overview of the important connectives. ${ }^{14}$ Their meaning will be explained subsequently, once we define compound statements and compound statement functions.

| negation: | $\neg A$ | not $A$ |
| :--- | :--- | :--- |
| conjunction: | $A \wedge B$ | $A$ and $B$ |
| double arrow (biconditional): | $A \leftrightarrow B$ | $A$ double arrow $B$ |
| logical equivalence: | $A \Leftrightarrow B$ | $A$ if and only if $B$ |
| disjunction (inclusive or): | $A \vee B$ | $A$ or $B$ |
| exclusive or: | $A$ xor $B$ | either $A$ or $B$, exactly one of $A$ or $B$ |
| arrow: | $A \rightarrow B$ | $A$ arrow $B$, if $A$ then $B$ |
| implication: | $A \Rightarrow B$ | $A$ implies $B$, if $A$ then $B$ |

Notations 3.1 (use of symbols vs descriptive English).
a. In the entire chapter on logic we generally use for logical operators their symbols like " $\neg$ " or $" \Rightarrow$ " in formulas but we use their corresponding English expressions (not and implies in this case) in connection with constructs which contain English language.

For example we would write $\neg(A \vee \neg B)$ rather than $\operatorname{not}(A$ or not $B)$ but we would write " $d+7=$ 16 and $F$ is a $\operatorname{dog}^{\prime \prime}$ rather than " $(d+7=16) \wedge(F$ is a dog $)$ "

[^10]b. Outside chapter 3 symbols are not used at all for logical operators. We use boldface such as "and" rather than just plain type face only to make it visually easier to understand the structure of a mathematical construct which employs connectives.

Definition 3.5 (Compound statements). A statement which does not contain any logical operators is called a simple statement and one that employs logical operators is called a compound statement.

Similarly statement functions which contain logical operators are called compound statement functions.

Example 3.10. Statements $\mathbf{e}$ and $\mathbf{f}$ of example 3.3 are examples of compound statements.
In $\mathbf{e}$ the two simple statements "The sum of any two even integers is even" and "Roses are mammals" are connected by and.

In $\mathbf{f}$ the two simple statements "The sum of any two even integers is even" and "Roses are mammals" are connected by either ... or.

### 3.3.2 Negation and Conjunction, Truth Tables and Tautologies (Understand this!)

We now give the definition of the first two logical operators which were introduced in the table of section 3.3.1.
Definition 3.6 (Negation). The negation operator is represented by the symbol " $\neg$ " and it reverses the truth value of a statement $A$, i.e., if $A$ is true then $\neg(A)$ is false and if $A$ is false then $\neg(A)$ is true.
(3.7) This is expressed in this "truth table" for $\neg A:^{15}$

| A | $\neg \mathrm{A}$ |
| :---: | :---: |
| F | T |
| T | F |

Example 3.11. Let $A:=$ "Rover is a horse". Then $\neg A=$ "Rover is not a horse" and $\neg \neg A=\neg(\neg A)=$ "Rover is a horse" $=A$.

Let us not quibble here about whether $\neg \neg A$ is not in reality the statement "Rover is not not a horse" which admittedly means the same as "Rover is a horse" but looks different.

There is no question about the fact that the T/F values for $A$ and $\neg \neg A$ are the same. Just compare column 1 with column 3 .

| $A$ | $\neg A$ | $\neg(\neg A)$ |
| :---: | :---: | :---: |
| F | T | F |
| T | F | T |

Note that we did not use any specifics about $A$. We derived the $\mathrm{T} / \mathrm{F}$ values for $\neg \neg A$ from those in the second column by applying the definition of the $\neg$ operator to the statement $B:=\neg A$.

In other words we have proved that the statements $A$ and $\neg \neg A$ are logically equivalent in the sense that one of them is true whenever the other one is true and vice versa.

All operators discussed subsequently are binary operators, i.e., they connect two input parameters (statements) $A, B$ and four rather than two rows are needed to show what will happen for each of the four combinations $A$ : false and B: false, $A$ : false and B: true, $A$ : true and B: false, $A$ : true and B: true.

In contrast, the already discussed negation operator " $\neg$ " is a unary operators, i.e., it has a single input parameter. We will keep referring to " $\neg$ " as a connective even though there are no two or more items that can be connected.

[^11]Definition 3.7 (Conjunction). The conjunction operator is represented by the symbols " $\wedge$ " or "and". The expression $A$ and $B$ is true if and only if both $A$ and $B$ are true.
(3.8) Truth table for $A$ and $B$ :

| $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | F |
| T | F | F |
| T | T | T |

The and connective generalizes to more than two statements $A_{1}, A_{2}, \ldots, A_{n}$ in the obvious manner:
$A_{1} \wedge A_{2} \wedge \cdots \wedge A_{n}$ is true if and only if each one of $A_{1}, A_{2}, \ldots, A_{n}$ is true and false otherwise.
Definition 3.8 (Truth table). A truth table contains the symbols for statements in the header, i.e., the top row and shows in subsequent rows how their truth values relate.

It contains in the leftmost columns statements which you may think of as varying inputs and it contains in the columns to the right compound statements which were built from those inputs by the use of logical operators. We have a row for each possible combination of truth values for the input statements. Such a combination then determines the truth value for each of the other statements.

When we count rows we start with zero for the header which contains the statement names. Row 1 is the first row which contains T/F values.

An example for a truth table is the following table which you encountered in the definition above 3.7 of the conjunction operator:

| $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | F |
| T | F | F |
| T | T | T |

Here the input statements are $A$ and $B$. The compound statement $A \wedge B$ is built from those inputs with the use of the $\wedge$ operator. We have 4 possible T/F combinations for $A$ and $B$ and each one of those determines the truth value of $A \wedge B$. For example, row 2 contains $A: \mathrm{F}$ and $B: \mathrm{T}$ and from this we obtain F as the corresponding truth value of $A \wedge B$.
Some truth tables have more than two inputs. If there are three statements $A, B, C$ from wich the compound statements that interest us are built then there will be $2^{3}=8$ rows to hold all possible combinations of truth values and for $n$ inputs there will be $2^{n}$ rows.

Definition 3.9 (Logically impossible). The statements $A$ and $B$ in the truth table of def.3.8 were of a generally nature and all four $T / F$ combinations had to be considered. If we deal with statements which are more specific but have some variability because they contain place holders ${ }^{16}$ then there may be dependencies that rule out certain combinations as nonsensical. For example let $x$ be some fixed but unspecified number and look at a truth table which has the statements $A:=A(x):=$ " $x>$ $5^{\prime \prime}$ and $B:=B(x):=$ " $x>7$ " as input. It is clearly impossible that $A$ is false and $B$ is true, no matter what value $x$ may have.

We call such combinations logically impossible or contradictory. We abbreviate "logically impossible" with L/I.

[^12]Both truth tables indicate that the combination $A: \mathrm{F}$ and $B: \mathrm{T}$ is logically impossible for $A=" x>5$ " and $B=" x>7$ ".

$\left.$| $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | $\mathrm{~L} / \mathrm{I}$ |
| T | F | F |
| T | T | T |$\quad \right\rvert\,$| $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: |
| F | F | F |
| T | F | F |
| T | T | T |

Remark 3.2. It was mentioned in the definition of logically impossible T/F combinations that there had to be some relationship between the inputs, i.e., some placeholders or some fixed but unspecified constants to make this an interesting definitions.

Consider what happens if you have two statements $A$ and $B$ for which this is not the case. For example, let $A:=$ "All tomatoes are blue" (obviously false) and $B:=$ "Arkansas is a state of the U.S.A." (obviously true).

For those two specific statements we know upfront that we have $A: \mathrm{F}$ and $B: \mathrm{T}$, so why bother with the other three cases? In other words, the appropriate truth table is either of those two:

| $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: |
| F | F | $\mathrm{~L} / \mathrm{I}$ |
| F | T | F |
| T | F | $\mathrm{L} / \mathrm{I}$ |
| T | T | $\mathrm{L} / \mathrm{I}$ |


\section*{| $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: |
| F | T | F |}

Remark 3.3. We chose for a more compact notation to place " $\mathrm{L} / \mathrm{I}$ " into one of the statement columns but be aware that the L/I attribute really belongs to certain combinations of the T/F values of the inputs. In other words,
the $\mathrm{L} / \mathrm{I}$ attribute belongs to certain rows of the truth table. A more accurate way would be to place L/I into a separate status column and place "N/A" or "-" or nothing into all columns other than those for the inputs:

| Status | $A$ | $B$ | $A \wedge B$ |
| :---: | :---: | :---: | :---: |
| $\mathrm{~L} / \mathrm{I}$ | F | F | F |
|  | F | T | - |
|  | T | F | F |
|  | T | T | T |

Of course more than two input statements can be involved when discussing logical impossibility. The following example will show this.

Example 3.12. Let $U, V, W, Z$ be the statement functions
$U:=x \mapsto U(x):=$ " $x \in[0,4]$ ",
$V:=x \mapsto V(x):=$ " $x \notin \emptyset^{\prime \prime}$,
$W:=x \mapsto W(x):=$ " $x<-1$ ",
$Z:=x \mapsto Z(x):=" x>2^{\prime \prime}$
with $\mathrm{UoD} \mathbb{R}$ in each case. Let $Q$ be a statement function that is built from $U, V, W, Z$ with the help of logical operators.

We observe the following:
a. $\quad V(x)$ is always true because the empty set does not contain any elements.
$\mathbf{a}^{\prime}$. In other words, there is no $x$ in the UoD for which $V(x)$ is false.
b. There is no $x$ in the UoD for which $W(x)$ and $Z(x)$ can both be true.

The following rows in the resulting truth table yield an $L / I$ regardless whether we enter a truth value of $T$ or F into anyone of the " $\bullet$ " entries.

| $U(x)$ | $V(x)$ | $W(y)$ | $Z(x)$ | $Q(x)$ |
| :---: | :---: | :---: | :---: | :---: |
| $\bullet$ | F | $\bullet$ | $\bullet$ | $\mathrm{L} / \mathrm{I}$ |
| $\bullet$ | $\bullet$ | T | T | $\mathrm{L} / \mathrm{I}$ |

Remark 3.4. As in example 3.12 above let
$U:=U(x):=" x \in[0,4] ", V:=V(x):=" x \notin \emptyset ", W:=W(x):=" x<-1 ", Z:=Z(x):=" x>$ $2^{\prime \prime}$.
a. The statement ${ }^{17} Q(x):=\neg(U(x) \wedge V(x)) \wedge W(x) \wedge Z(x)$ can never be true, regardless of $x$.

To see this directly note again that $V(x)$ is trivially true for any $x$ because the emptyset by definition does not contain any elements. It follows that $U(x) \wedge V(x)$ means " $x \in[0,4]$ " and $Q(x)$ means " $x<0^{\prime \prime} \wedge$ " $x>4^{\prime \prime} \wedge$ " $x<-1^{\prime \prime} \wedge$ " $x>2^{\prime \prime}$ which is equivalent to " $x<-1^{\prime \prime} \wedge$ " $x>4$ " and certainly false for any $x$ in the UoD, i.e., $x \in \mathbb{R}$.

Alternatively we can use the results from example 3.12 where we found out that $W(x)$ and $Z(x)$ cannot both be true at the same time.

The remaining rows in the resulting truth table yield an F for $Q(x)$ regardless of the truth values of $U(x)$ and $V(x)$ because $W(x) \wedge Z(x)$ is false, hence $Q(x)=$ whatever $\wedge(W(x) \wedge Z(x))$ is false for those remaining rows.

| $U(x)$ | $V(x)$ | $W(y)$ | $Z(x)$ | $Q(x)$ |
| :---: | :---: | :---: | :---: | :---: |
| $\bullet$ | $\bullet$ | F | F | F |
| $\bullet$ | $\bullet$ | F | T | F |
| $\bullet$ | $\bullet$ | T | F | F |

b. Let $R: x \mapsto R(x):=\neg Q(x)$ be the statement function with $\operatorname{UoD} \mathbb{R}$ which represents for each $x$ in the UoD the opposite of $Q$. Because $Q(x)$ is false for all $x, R(x)$ is true for all $x$ in the universe of discourse for $x$.

Statements which are true or false under all circumstances like the statements $R(x)$ and $Q(x)$ from the remark above deserve special names.

Definition 3.10 (Tautologies and contradictions). A tautology is a statement which is true under all circumstances, i.e., under all combinations of truth values which are not logically impossible.

A contradiction is a statement which is false under all circumstances.
We write $T_{0}$ for the tautology " $1=1$ " and $F_{0}$ for the contradiction " $1=0$ ". This gives us a convenient way to incorporate statements which are true or false under all circumstances into formulas that build compound statements.

Example 3.13. Here are some examples of tautologies.
a. The statements $R(x)$ of remark 3.4 are tautologies.
b. $T_{0}$ is a boring example of a tautology. So is any true statement without any variables such as " $9+12=21$ " and "a cat is not a cow".
c. There are formulas involving arbitrary statements which are tautologies. We will show that for any two statements $A$ and $B$ the statement $P:=\neg(A \wedge \neg A)$ is a tautology.

Here are some examples of contradictions.
d. The statements $Q(x)$ of remark 3.4 are contradictions.
e. $F_{0}$ is a boring example of a contradiction. So is any false statement without any variables such as " $9+12=50$ " and "a dog is a whale".

[^13]f. There are formulas involving arbitrary statements which are contradictions. We will show that for any two statements $A$ and $B$ the statement $Q:=(A \wedge \neg A) \wedge B$ is a contradiction.

PROOF of $c$ and $f$ :
$P=\neg(A \wedge \neg A)$ (last column) has entries all $T$, hence $P$ is a tautology.
$Q=(A \wedge \neg A) \wedge B$ (next to last column) has entries all $F$, hence $Q$ is a contradiction.

| $A$ | $B$ | $\neg A$ | $A \wedge \neg A$ | $(A \wedge \neg A) \wedge B$ | $\neg(A \wedge \neg A)$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $F$ | $F$ | $T$ | $F$ | $F$ | $T$ |
| $F$ | $T$ | $T$ | $F$ | $F$ | $T$ |
| $T$ | $F$ | $F$ | $F$ | $F$ | $T$ |
| $T$ | $T$ | $F$ | $F$ | $F$ | $T$ |

We now continue with the conjunction operator.
Example 3.14. In the following let $x, y$ be two (fixed but arbitrary) integers and let $A(x):=" x \in \mathbb{N}$ " and $B(y):=" y \in \mathbb{Z}$ and $y>0$ ". Be sure to understand that $A(x)$ and $B(y)$ are in fact statements and not predicates, because the symbols $x, y$ are bound from the start and hence cannot be considered variables of the predicates $A:=" x \in \mathbb{N}^{\prime \prime}$ and $B:=" y \in \mathbb{Z}$ and $y>0$ ".

We will reuse the statements $A(x)$ and $B(y)$ in examples for the subsequently defined logical operators.
a. If no assumptions are made about a relationship between $x$ and $y$ then all four T/F combinations are possible and, to explore conjunction, we must deal with the full truth table

| $A(x)$ | $B(y)$ | $A(x) \wedge B(y)$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | F |
| T | F | F |
| T | T | T |

b. On the other hand, if $x<y$ then the truth of $A(x)$ implies that of $B(y)$ because if $y$ is an integer which dominates some natural number $x$ then we have $y>x \geqq 1>0$, i.e., $y$ is an integer bigger than zero, i.e., truth of $A(x)$ and falseness of $B(y)$ are incompatible.

It follows that the combination T/F is L/I. We discard the corresponding row and restrict ourselves to the truth table

| $A(x)$ | $B(y)$ | $A(x) \wedge B(y)$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | F |
| T | T | T |

c. Even better, if $x=y$, i.e., we compare truth/falsehood of $A(x)$ with that of $B(x)$, we only need to worry about the two combinations $\mathrm{F} / \mathrm{F}$ and $\mathrm{T} / \mathrm{T}$ for the following reason: The set of positive integers is the set $\{1,2, \ldots\}$ and this is, by definition, the set $\mathbb{N}$ of all natural numbers. This means that the statements " $x \in \mathbb{N}$ " and " $y \in \mathbb{Z}$ and $y>0$ " are just two different ways of expressing the same thing.

It follows that either both $A(x)$ and $B(x)$ are true or both are false. We discard the logically impossible combinations F/T and T/F and restrict ourselves to the truth table

| $A(x)$ | $B(x)$ | $A(x) \wedge B(x)$ |
| :---: | :---: | :---: |
| F | F | F |
| T | T | T |

### 3.3.3 Biconditional and Logical Equivalence Operators - Part 1

Definition 3.11 (Double arrow operator (biconditional)). The double arrow operator ${ }^{18}$ is represented by the symbol " $\leftrightarrow$ " and read " $A$ double arrow $B$ ". $A \leftrightarrow B$ is true if and only if either both $A$ and $B$ are true or both $A$ and $B$ are false.
(3.9) Truth table for $A \leftrightarrow B$ :

| $A$ | $B$ | $A \leftrightarrow B$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | F |
| T | F | F |
| T | T | T |

Definition 3.12 (Logical equivalence operator). Two statements $A$ and $B$ are logically equivalent if the statement $A \leftrightarrow B$ is a tautology, i.e., if the combinations $A$ :true, $B$ :false and $A$ :false, $B$ :true both are logically impossible.

We write $A \Leftrightarrow B$ and we say " $A$ if and only if $B$ " to indicate that $A$ and $B$ are logically equivalent.
(3.10) Truth table for $A \Leftrightarrow B$ :

| A | B | $A \Leftrightarrow B$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | $\mathrm{~L} / \mathrm{I}$ |
| T | F | $\mathrm{L} / \mathrm{I}$ |
| T | T | T |

The discussion of the $\leftrightarrow$ and $\Leftrightarrow$ operators will be continued in ch.3.3.6 (Biconditional and Logical Equivalence Operators - Part 2) on p. 42

### 3.3.4 Inclusive and Exclusive Or

Definition 3.13 (Disjunction). The disjunction operator is represented by the symbols " $\vee$ " or "or". The expression $A$ or $B$ is true if and only if either $A$ or $B$ is true.

| $A$ | $B$ | $A \vee B$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | T |
| T | F | T |
| T | T | T |

The or connective generalizes to more than two statements $A_{1}, A_{2}, \ldots, A_{n}$ in the obvious manner:
$A_{1} \vee A_{2} \vee \cdots \vee A_{n}$ is true if and only if at least one of $A_{1}, A_{2}, \ldots, A_{n}$ is true and false otherwise, i.e., if each of the $A_{k}$ is false.

Example 3.15. As in example 3.14 let $x, y \in \mathbb{Z}$ and let
$A(x):=" x \in \mathbb{N}^{\prime}$ and $B(y):=" y \in \mathbb{Z}$ and $y>0$ "

[^14]a. If no assumptions are made about a relationship between $x$ and $y$ then all four T/F combinations are possible and, to explore conjunction, we must deal with the full truth table
b. Let $x<y$. We have seen in example 3.14.b that the combination T/F is impossible and we can restrict ourselves to the simplified truth table
c. Now let $x=y$. We have seen in example 3.14.c that either both $A(x)$ and $B(y)=B(x)$ are true or both are false. Because the combinations F/T and T/F are impossible we can restrict ourselves to the simplified truth table

| $A(x)$ | $B(y)$ | $A(x) \vee B(y)$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | T |
| T | F | T |
| T | T | T |
| $A(x)$ | $B(y)$ | $A(x) \vee B(y)$ |
| F | F | F |
| F | T | T |
| T | T | T |
| $A(x)$ | $B(x)$ | $A(x) \wedge B(x)$ |
| F | F | F |
| T | T | T |

Definition 3.14 (Exclusive or). The exclusive or operator is represented by the symbol "xor". ${ }^{19} \mathrm{~A}$ xor $B$ is true if and only if either $A$ or $B$ is true (but not both as is the case for the inclusive or).
(3.12) Truth table for $A$ xor $B$ :

| A | B | A xor B |
| :---: | :---: | :---: |
| F | F | F |
| F | T | T |
| T | F | T |
| T | T | F |

Example 3.16. As in example 3.14 let $x, y \in \mathbb{Z}$ and let $A(x):=" x \in \mathbb{N}^{\prime}$ and $B(y):=" y \in \mathbb{Z}$ and $y>0 "$
a. If no assumptions are made about a relationship between $x$ and $y$ then all four T/F combinations are possible and, to explore conjunction, we must deal with the full truth table
b. Let $x<y$. We have seen in example 3.14.b that the combination T/F is impossible and we can restrict ourselves to the simplified truth table
c. Now let $x=y$. We have seen in example 3.14.c that either both $A(x)$ and $B(y)=B(x)$ are true or both are false. Because the combinations F/T and T/F are impossible we can restrict ourselves to the simplified truth table

| $A(x)$ | $B(y)$ | $A(x)$ xor $B(y)$ |
| :---: | :---: | :---: |
| F | F | F |
| F | T | T |
| T | F | T |
| T | T | F |
| $A(x)$ | $B(y)$ | $A(x)$ xor $B(y)$ |
| F | F | F |
| F | T | T |
| T | T | F |
| $A(x)$ | $B(x)$ | $A(x)$ xor $B(x)$ |
| F | F | F |
| T | T | F |

This last truth table is remarkable. The truth values for $A(x)$ xor $B(x)$ are false in each row, hence it is a contradiction as defined in def.3.10 on p.33.
Remark 3.5. Note that the truth values for $A \leftrightarrow B$ are the exact opposites of those for $A$ xor $B$ :
$A \leftrightarrow B$ is true exactly when both $A$ and $B$ have the same truth value whereas $A$ xor $B$ is true exactly when $A$ and $B$ have opposite truth values. In other words,

[^15]$A \leftrightarrow B$ is true whenever $\neg[A$ xor $B]$ is true and false whenever $\neg[A$ xor $B]$ is false.
Exercise 3.1. use that last remark to prove that for any two statements $A$ and $B$ the compound statement
$$
[A \leftrightarrow B] \leftrightarrow \neg[A \text { xor } B]
$$
is a tautology.

### 3.3.5 Arrow and Implication Operators

Definition 3.15 (Arrow operator). The arrow operator ${ }^{20}$ is represented by the symbol " $\rightarrow$ ". We read $A \rightarrow B$ as " $A$ arrow $B$ " but see remark 3.7 below for the interpretation "if $A$ then $B$ ".
(3.13) Truth table for $A \rightarrow B$ :

| A | B | $A \rightarrow B$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | T |
| T | F | F |
| T | T | T |

In other words, $A \rightarrow B$ is false if and only if $A$ is true and $B$ is false.
Definition 3.16 (Implication operator). We say that $A$ implies $B$ and we write

$$
\begin{equation*}
A \Rightarrow B \tag{3.14}
\end{equation*}
$$

for two statements $A$ and $B$ if the statement $A \rightarrow B$ is a tautology, i.e., if the combination $A$ : true, $B$ : false is logically impossible.
(3.15) Truth table for $A \Rightarrow B$ :

| A | B | $A \Rightarrow B$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | T |
| T | F | $\mathrm{~L} / \mathrm{I}$ |
| T | T | T |

Remark 3.6. There are several ways to express $A \Rightarrow B$ in plain english:

Short form:
$A$ implies $B$
if $A$ then $B$
$A$ only if $B$
$B$ if $A$
$B$ whenever $A$
$A$ is sufficient for $B$
$B$ is necessary for $A$

Interpret this as:
The truth of $A$ implies the truth of $B$
if $A$ is true then $B$ is true
$A$ is true only if $B$ is true
$B$ is true if $A$ is true
$B$ is true whenever $A$ is true
The truth of $A$ is sufficient for the truth of $B$ The truth of $B$ is necessary for the truth of $A$

Theorem 3.1 (Transitivity of " $\Rightarrow$ "). Let $A, B, C$ be three statements such that $A \Rightarrow B$ and $B \Rightarrow C$. Then $A \Rightarrow C$.

[^16]PROOF:
$A \Rightarrow B$ means that the combination $A: T, B: F$ is logically impossible because otherwise $A \rightarrow B$ would have a truth value of $F$ and we would not have a tautology. Hence we can drop row 5 from the truth table on the right. Similarly we can drop row 7 because it contains the combination B:T, C:F which contradicts our assumption that $B \Rightarrow C$. But those are the only rows for which $A \rightarrow C$ yields false because only they contain the combination A:T, C:F. It follows that $A \rightarrow C$ is a tautology, i.e., $A \Rightarrow C$.

|  | $A$ | $B$ | $C$ |
| :---: | :---: | :---: | :---: |
| 1 | $F$ | $F$ | $F$ |
| 2 | $F$ | $F$ | $T$ |
| 3 | $F$ | $T$ | $F$ |
| 4 | $F$ | $T$ | $T$ |
| 5 | $T$ | $F$ | $F$ |
| 6 | $T$ | $F$ | $T$ |
| 7 | $T$ | $T$ | $F$ |
| 8 | $T$ | $T$ | $T$ |

Theorem 3.2 (Transitivity of " $\rightarrow$ "). Let $A, B, C$ be three statements.
Then $[(A \rightarrow B) \wedge(B \rightarrow C)] \Rightarrow(A \rightarrow C)$.
PROOF: We must show that $[(A \rightarrow B) \wedge(B \rightarrow C)] \rightarrow(A \rightarrow C)$ is a tautology. We do this by brute force and compute the truth table.

|  |  |  |  | $P:=$ <br> $A$ |  | $B$ | $C$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $A \rightarrow B$ | $B \rightarrow C$ | $(A \rightarrow B) \wedge(B \rightarrow C)$ | $A \rightarrow C$ | $P \rightarrow(A \rightarrow C)$ |  |  |  |
| $F$ | $F$ | $F$ | $T$ | $T$ | $T$ | $T$ | $T$ |
| $F$ | $F$ | $T$ | $T$ | $T$ | $T$ | $T$ | $T$ |
| $F$ | $T$ | $F$ | $T$ | $F$ | $F$ | $T$ | $T$ |
| $F$ | $T$ | $T$ | $T$ | $T$ | $T$ | $T$ | $T$ |
| $T$ | $F$ | $F$ | $F$ | $T$ | $F$ | $F$ | $T$ |
| $T$ | $F$ | $T$ | $F$ | $T$ | $F$ | $T$ | $T$ |
| $T$ | $T$ | $F$ | $T$ | $F$ | $F$ | $F$ | $T$ |
| $T$ | $T$ | $T$ | $T$ | $T$ | $T$ | $T$ | $T$ |

We see that the last column with the truth values for $[(A \rightarrow B) \wedge(B \rightarrow C)] \rightarrow(A \rightarrow C)$ contains true everywhere and we have proved that this statement is a tautology.
Definition 3.17. In the context of $A \rightarrow B$ and $A \Rightarrow B$ we call $A$ the premise or the hypothesis ${ }^{21}$ and we call $B$ the conclusion. ${ }^{22}$

We call $B \rightarrow A$ the converse of $A \rightarrow B$ and we call $\neg B \rightarrow \neg A$ the contrapositive of $A \rightarrow B$.
We call $B \Rightarrow A$ the converse of $A \Rightarrow B$ and we call $\neg B \Rightarrow \neg A$ the contrapositive of $A \Rightarrow B$.

## Remark 3.7.

a. The difference between $A \rightarrow B$ and $A \Rightarrow B$ is that $A \Rightarrow B$ implies a relation between the premise $A$ and the conclusion $B$ which renders the T/F combination $A: T, B: \mathrm{F}$ logically impossible, i.e., the pared down truth table has only true entries in the $A \Rightarrow B$ column. In other words, $A \Rightarrow B$ is the statement $A \rightarrow B$ in case the latter is a tautology as defined in def.3.10 on p.33.
b. Both $A \rightarrow B$ and $A \Rightarrow B$ are interpreted as "if $A$ then $B$ " but we prefer in general to say " $A$ arrow $B$ " for $A \rightarrow B$ because outside the realm of logic $A \Rightarrow B$ is what mathematicians use when they refer to "If ... then " constructs to state and prove theorems.

[^17]Example 3.17. The converse of "if $x$ is a dog then $x$ is a mammal" is "if $x$ is a mammal then $x$ is a dog". You see that, regardless whether you look at it in the context of $\rightarrow$ or $\Rightarrow$, a "if ...then" statement can be true whereas its converse will be false and vice versa.

The contrapositive of "if $x$ is a dog then $x$ is a mammal" is "if $x$ is not a mammal then $x$ is a not a dog". Switching to the contrapositive did not switch the truth value of the "if ... then" statement. This is not an accident: see the Contrapositive Law (3.41) on p.46.

## Remark 3.8. What is the connection between the truth tables for $A \rightarrow B, A \Rightarrow B$ and modeling

 "if $A$ then $B$ "?We answer this question as follows:
a. If the premise A is guaranteed to be false, you should be allowed to conclude from it anything you like:

Consider the following statements which are obviously false:
$F_{1}$ : "The average weight of a 30 year old person is 7 ounces",
$F_{2}$ : "The number 12.7 is an integer",
$F_{3}$ : "There are two odd integers $m$ and $n$ such that $m+n$ is odd",
$F_{4}$ : "All continuous functions are differentiable" ${ }^{23}$
and some that are known to be true:
$T_{1}$ : "The moon orbits the earth",
$T_{2}$ : "The number 12.7 is not an integer",
$T_{3}$ : "If $m$ and $n$ are even integers then $m+n$ is even",
$T_{4}$ : "All differentiable functions are continuous"
a1. What about the statement "if $F_{3}$ then $T_{1}$ ": "If There are two odd integers $m$ and $n$ such that $m+n$ is odd then the moon orbits the earth"? This may not make a lot of sense to you, but consider this:

The truth of "if $F_{3}$ then $T_{1}$ " is not the same as the truth of just $F_{1}$. No absolute claim is made that the moon orbits the earth. You are only asked to concede such is the case under the assumption that two odd integers can be found whose sum is odd. But we know that no such integers exist, i.e., we are dealing with a vacuous premise and there is no obligation on our part to show that the moon indeed orbits the earth! Because of this we should have no problem to accept the validity of "if $F_{3}$ then $T_{1}$ ". Keep in mind though that knowing that if $F_{3}$ then $T_{1}$ will not help to establish the truth or falseness of $T_{1}$ !
a2. Now what about the statement "if $F_{3}$ then $F_{2}$ ": "If There are two odd integers $m$ and $n$ such that $m+n$ is odd then the number 12.7 is an integer"? The truth of this implication should be much easier to understand than allowing to conclude something false from something false:

When was the last time that someone bragged "Yesterday I did xyz" and you responded with something like "If you did xyz then I am the queen of Sheba" in the serene knowledge that there is no way that this person could have possibly done xyz? You know that you have no burden of proof to show that you are the queen of Sheba because you did not make this an absolute claim: You hedged that such is only the case if it is true that the other person in fact did xyz yesterday.

So, yes, the argument "if $F_{3}$ then $F_{2}$ ". sounds OK and we should accept it as true but, as in the case of "if $F_{3}$ then $T_{1}$ ". this has no bearing on the truth or falseness of $F_{2}$.

[^18]To summarize, "if F then $B^{\prime \prime}$. should be true, no matter what you plug in for $B$. We thus have obtained the first two rows of a sensible truth table for $\mathrm{A} \rightarrow \mathrm{B}$ :

| A | B | $\mathrm{A} \rightarrow \mathrm{B}$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | T |

b. Is it OK to say that if the premise $A$ is true then we may infer that the conclusion $B$ is also true? Definitely! There is nothing wrong with "if $T_{2}$ then $T_{4}$ ", i.e., the statement "If The number 12.7 is not an integer then all differentiable functions are continuous"

We can add the fourth row but we do not have \#3 yet:

| A | B | $\mathrm{A} \rightarrow \mathrm{B}$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | T |
| T | F | $? ?$ |
| T | T | T |

c. Is it OK to say that, if the premise A is true, we may say in parallel that A implies B even if the conclusion B is false? No way! Let's assume that Jane is a goldfish. Then A: "Jane is a fish" is true and B: "Jane is a rocket scientist" is false. It is definitely NOT OK to say, under those circumstances, "If Jane is a fish" then Jane is a rocket scientist". Contrast that with this modification that fits case b: "If Jane is a fish' then Jane is not a rocket scientist". No one should have a problem with that! We now can complete row \#3: $\mathrm{T} \rightarrow \mathrm{F}$ is false.

We now have the complete truth table for $A \rightarrow B$ and it matches the one in def.3.15:

The truth table (3.15) for $A \Rightarrow B$ is then derived from that for $A \rightarrow B$ by demanding that $A$ and $B$ be such that $A \rightarrow B$ cannot be false, i.e, the combination $A: F, B: T$ must be logically impossible:

| $A$ | $B$ | $A \rightarrow B$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | T |
| T | F | F |
| T | T | T |
| $A$ | $B$ | $A \Rightarrow B$ |
| F | F | T |
| F | T | T |
| T | F | $\mathrm{~L} / \mathrm{I}$ |
| T | T | T |

We arrived in this remark at the truth tables for $A \rightarrow B$ and $A \Rightarrow B$ based on what seems to be reasonable. But the discipline of logic is as exacting a subject as abstract math and the process had to be done in reverse: We first had to define $A \rightarrow B$ and $A \Rightarrow B$ by means of the truth tables given in def.3.15 and def.3.16 and from there we justified why these operators appropriately model "if $A$ then $B^{\prime \prime}$.

Example 3.18. As in example 3.14 let $x, y \in \mathbb{Z}$ and let
$A(x):=" x \in \mathbb{N}^{\prime}$ and $B(y):=" y \in \mathbb{Z}$ and $y>0 "$
a. If no assumptions are made about a relationship between $x$ and $y$ then all four T/F combinations are possible and, to explore conjunction, we must deal with the full truth table
b. Let $x<y$. We have seen in example 3.14.b that the combination T/F is impossible and we can restrict ourselves to the simplified truth table

| $A(x)$ | $B(y)$ | $A(x) \rightarrow B(y)$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | T |
| T | F | F |
| T | T | T |
| $A(x)$ | $B(y)$ | $A(x) \rightarrow B(y)$ |
| F | F | T |
| F | T | T |
| T | T | T |

c. Now let $x=y$. We have seen in example 3.14.c that either both $A(x)$ and $B(y)=B(x)$ are true or both are false. Because the combinations F/T and T/F are impossible we can restrict ourselves to the simplified truth table

| $A(x)$ | $B(x)$ | $A(x) \rightarrow B(x)$ |
| :---: | :---: | :---: |
| F | F | T |
| T | T | T |

We see that $A(x) \rightarrow B(y)$ is a tautology in case that $x<y$ or $x=y$.
We have seen that some work was involved to show that the " $A(x) \rightarrow B(y)$ " statement of the last example is a tautology. How do we interpret this?

If you show that a "if P then $Q$ " statement is a tautology then you have demonstrated that a true premise necessarily results in a true conclusion. You have "proved" the validity of the conclusion $Q$ from the validity of the hypothesis $P$.

The next example is a modification of the previous one. We replace the statements $A(x)$ and $B(y)$ with statement functions $x \mapsto A(x), y \mapsto B(y),(x, y) \mapsto C(x, y)$. and replace $A(x) \rightarrow B(y)$ with an equivalent $\rightarrow$ statement which involves those three statement functions. Our goal is now to show that this new if . . . then statement is a tautology for all $x$ and $y$ which belong to their universes of discourse.

Example 3.19. Let $\mathscr{U}_{x}:=\mathscr{U}_{y}:=\mathbb{Z}$ be the UoDs for the variables $x$ and $y$.

$$
\text { Let } \begin{array}{ll}
A: \mathscr{U}_{x} \rightarrow \mathscr{S} \text { with } x \mapsto " x \in \mathbb{N}^{\prime \prime}, \\
& B: \mathscr{U}_{y} \rightarrow \mathscr{S} \text { with } y \mapsto " y \in \mathbb{Z} \text { and } y>0^{\prime \prime}, \\
& C: \mathscr{U}_{x} \times \mathscr{U}_{y} \rightarrow \mathscr{S} \text { with }(x, y) \mapsto " x<y^{\prime \prime} .
\end{array}
$$

Let us try to show that for any $x$ in the $\operatorname{UoD}$ of $x$ and $y$ in the $\operatorname{UoD}$ of $y$, i.e., for any two integers $x$ and $y$, the function value $T(x, y)$ of the statement function

$$
\begin{equation*}
T: \mathscr{U}_{x} \times \mathscr{U}_{y} \rightarrow \mathscr{S} \text { with }(x, y) \mapsto T(x, y):=[(A(x) \wedge C(x, y)) \rightarrow B(y)] \text { is a tautology. } \tag{3.16}
\end{equation*}
$$

Note that
a. The last arrow in (3.16) is the arrow operator $\rightarrow$, not the function assignment operator $\mapsto$.
b. if we can demonstrate that (3.16) is correct then we can replace $(A(x) \wedge C(x, y)) \rightarrow B(y)$ with $(A(x) \wedge C(x, y)) \Rightarrow B(y)$. We interpret this as having proved the (trivial)
Theorem: It is true for all integers $x$ and $y$ that if $x \in \mathbb{N}$ and $x<y$ then $y \in \mathbb{Z}$ and $y>0$.
The trick is of course to think of $x$ and $y$ not as placeholders but as fixed but unspecified integers. Then $A(x), B(y)$ and $C(x, y)$ are ordinary statements and we can build truth tables just as always. Observe that we now have three "inputs" $A(x), B(y)$ and $C(x, y)$ and the full truth table contains nine entries.

We need not worry about numbers $x$ and $y$ whose combination $(x, y)$ results in the falseness of the premise $A(x) \wedge C(x, y)$ because false $\rightarrow B(y)$ always results in true. In other words we do not worry about any combination of $x$ and $y$ for which at least one of $A(x), C(x, y)$ is false. To phrase it differently we focus on such $x$ and $y$ for which we have that both $A(x), C(x, y)$ are true and eliminate all other rows from the truth table. There are only two cases to consider: either $B(y)$ is false or $B(y)$ is true:

| $A(x)$ | $C(x, y)$ | $B(y)$ | $A(x) \wedge C(x, y)$ | $(A(x) \wedge C(x, y)) \rightarrow B(y)$ |
| :---: | :---: | :---: | :---: | :---: |
| T | T | F | T | F |
| T | T | T | T | T |

The proof is done if it can be shown that the first row is a logically impossible. We now look at the components $A(x), C(x, y), B(y)$ in context. We have seen in example 3.14b. that the assumed truth of $C(x, y)$ together with that of $A(x)$ is incompatible with $B(y)$ being false. This eliminates the first row from that last truth table and what remains is

$$
\begin{array}{|c|c|c|c|c|}
A(x) & C(x, y) & B(y) & A(x) \wedge C(x, y) & (A(x) \wedge C(x, y)) \rightarrow B(y) \\
\hline \mathrm{T} & \mathrm{~T} & \mathrm{~T} & \mathrm{~T} & \mathrm{~T}
\end{array}
$$

In other words we obtain the value true for all non-contradictory combinations in the last column of the truth table and this proves (3.16).

Remark 3.9. Let us compare example 3.18.b with example 3.19. Besides using statements in the former and predicates in the latter a more subtle difference is that, because $x$ and $y$ were assumed to be known from the outset,
example 3.18.b allowed us to formulate a truth table in which none of the statements had to explicitly refer to the condition $x<y$.

In contrast to this we had to introduce in example 3.19 the predicate $C=$ " $x<y$ " to bring this condition into the truth tables

Was there any advantage of switching from statements to predicates and adding a significant amount of complexity in doing so? The answer is yes but it will only become clear when we introduce quantifiers for statement functions.

We will come back to the subject of proofs in chapter 3.7.1 (Building blocks of mathematical theories) on p.56.

### 3.3.6 Biconditional and Logical Equivalence Operators - Part 2 (Understand this!)

This chapter continues the discussion of the $\leftrightarrow$ and $\Leftrightarrow$ operators from ch.3.3.3 (Biconditional and Logical Equivalence Operators - Part 1) on p. 35 .

## Remark 3.10.

a. Equivalence $A \Leftrightarrow B$ provides a "replacement principle for statements": Logically equivalent statements are not "semantically identical" but they cannot be distinguished as far as their "logic content", i.e., the circumstances under which they are true or false are concerned.
b. Note that $A \Leftrightarrow B$ means the same as the following: $A$ is true whenever $B$ is true and $A$ is false whenever $B$ is false because this is the same as saying that, in a truth table that contains entries for $A$ and $B$, each row either has the value T in both columns or the value F in both columns. This in turn is the same as saying that the column for $A \leftrightarrow B$ has T in each row, i.e., $A \leftrightarrow B$ is a tautology.
$\mathbf{b}^{\prime}$. There is not much value to $\mathbf{b}$ if $A$ and $B$ are simple statements but things become a lot more interesting if compound statements like $A:=\neg(P \wedge Q)$ and $B:=\neg P \vee \neg Q$ are looked at.

We illustrate the above remark with the following theorem.
Theorem 3.3 (De Morgan's laws for statements). Let $A$ and $B$ be statements. Then we have the following logical equivalences:

$$
\begin{align*}
& \neg(A \wedge B) \Leftrightarrow \neg A \vee \neg B,  \tag{3.17}\\
& \neg(A \vee B) \Leftrightarrow \neg A \wedge \neg B . \tag{3.18}
\end{align*}
$$

Those formulas generalize to $n$ statements $A_{1}, A_{2}, \ldots, A_{n}$ as follows:

$$
\begin{align*}
& \neg\left(A_{1} \wedge A_{2} \wedge \cdots \wedge A_{n}\right) \Leftrightarrow \neg A_{1} \vee \neg A_{2} \vee \cdots \vee \neg A_{n},  \tag{3.19}\\
& \neg\left(A_{1} \vee A_{2} \vee \cdots \vee A_{n}\right) \Leftrightarrow \neg A_{1} \wedge \neg A_{2} \wedge \cdots \wedge \neg A_{n} . \tag{3.20}
\end{align*}
$$

PROOF of 3.17: Here is the truth table for both $\neg(A \wedge B)$ and $\neg A \vee \neg B$ depending on the truth values of $A$ and $B$.

| $A$ | $B$ | $A \wedge B$ | $\neg(A \wedge B)$ | $\neg A$ | $\neg B$ | $\neg A \vee \neg B$ | $[\neg(A \wedge B)] \leftrightarrow[\neg A \vee \neg B]$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $F$ | $F$ | $F$ | $T$ | $T$ | $T$ | $T$ | $T$ |
| $F$ | $T$ | $F$ | $T$ | $T$ | $F$ | $T$ | $T$ |
| $T$ | $F$ | $F$ | $T$ | $F$ | $T$ | $T$ | $T$ |
| $T$ | $T$ | $T$ | $F$ | $F$ | $F$ | $F$ | $T$ |

This proves the validity of 3.17. Note that the last column of the truth table is superfluous because getting $T$ in each row follows from the fact that the rows of the statement to the left and the one to the right of " $\leftrightarrow$ " both contain the same entries T-T-T-F. The column has been included because it illustrates what was said in remark 3.10.
PROOF of 3.18: Left as an exercise.
Example 3.20. As in example 3.14 let $x, y \in \mathbb{Z}$ and let $A(x):=" x \in \mathbb{N}^{\prime}$ and $B(y):=" y \in \mathbb{Z}$ and $y>0$ "
a. If no assumptions are made about a relationship between $x$ and $y$ then the full truth table needs all four entries and we obtain

| $A(x)$ | $B(y)$ | $A(x) \leftrightarrow B(y)$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | F |
| T | F | F |
| T | T | T |

b. Let $x<y$. We have seen in example 3.14 that the combination $T / F$ is impossible and we can restrict ourselves to the simplified truth table

| $A(x)$ | $B(y)$ | $A(x) \rightarrow B(y)$ |
| :---: | :---: | :---: |
| F | F | T |
| F | T | F |
| T | T | T |

c. Now let $x=y$. We have seen in example 3.14.c that then either $A(x)$ and $B(y)=B(x)$ must both be true or they must both be false. Because the combinations F/T and T/F are impossible we can restrict ourselves to the simplified truth table

| $A(x)$ | $B(x)$ | $A(x) \leftrightarrow B(x)$ |
| :---: | :---: | :---: |
| F | F | T |
| T | T | T |

It follows that for any given number $x$ the statement $A(x) \leftrightarrow B(x)$ is always true, irrespective of the truth values of $A(x)$ and $B(x)$. Hence $A(x) \leftrightarrow B(x)$ is a tautology and we can write $A(x) \Leftrightarrow$ $B(x)$ for all $x$.

### 3.3.7 More Examples of Tautologies and Contradictions (Understand this!)

Now that we have all logical operators at our disposal we can give additional examples of tautologies and contradictions.

Example 3.21. In the following let $P, Q, R$ be three arbitrary statements, let $x, y$ be two (fixed but arbitrary) integers and let $A(x):=" x \in \mathbb{N}^{\prime \prime}$ and $B(y):=" y \in \mathbb{Z}$ and $y>0^{\prime \prime}$. (see example 3.14 on p. 34).
a. Tautologies:
$T_{0}$,
$A_{1}:=" 5+7=12$ ",
$A_{2}:=$ "Any integer is even or odd",
$A_{3}:=P \vee \neg P$ (Tertium non datur or law of the excluded middle),
$A_{4}:=P \vee T_{0}$,
$A_{5}:=(P \wedge Q) \vee(P \wedge \neg Q)$,
$A_{6}:=(P \rightarrow Q) \leftrightarrow(\neg P \vee Q)$ (Implication is logically equivalent to an or statement),
$A_{7}:=\left[" x<y^{\prime \prime} \wedge A(x)\right] \rightarrow B(y)$ (see 3.18.b on p.40),
$A_{8}:=A(x) \leftrightarrow B(x)$ (see 3.18.c).
Note that we can express the fact that $A_{6}, A_{7}, A_{8}$ are tautologies as follows:

$$
(P \rightarrow Q) \Leftrightarrow(\neg P \vee Q), \quad\left[" x<y^{\prime \prime} \wedge A(x)\right] \Rightarrow B(y), \quad A(x) \Leftrightarrow B(x) .
$$

b. Contradictions:
$F_{0}$,
$B_{1}:=" 5+7=15 "$,
$B_{2}:=$ "There are some non-zero numbers $x$ such that $x=2 x$ ",
$B_{3}:=P \wedge \neg P$,
$B_{3}:=P \wedge F_{0}$,
$B_{4}:=F_{0} \wedge(P \vee \neg P)$,
$B_{5}:=[\neg P \vee \neg Q] \wedge[P \wedge Q]$,
$B_{6}:=A(x)$ xor $B(x)$ (see 3.16.c on p. 36).
Proof that $A_{3}$ is a tautology:

| $P$ | $\neg P$ | $P \vee \neg P$ |
| :---: | :---: | :---: |
| $F$ | $T$ | $T$ |
| $T$ | $F$ | $T$ |

Proof that $A_{4}$ is a tautology:

| $P$ | $T_{0}$ | $P \vee T_{0}$ |
| :---: | :---: | :---: |
| $F$ | $T$ | $T$ |
| $T$ | $T$ | $T$ |

Note that even though there are two inputs, $P$ and $T_{0}$, there are only two valid combinations of truth values because the only choice for $T_{0}$ is true.

Proof that $A_{6}$ is a tautology:

| $P$ | $Q$ | $P \rightarrow Q$ | $\neg P$ | $\neg P \vee Q$ | $(P \rightarrow Q) \leftrightarrow(\neg P \vee Q)$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $F$ | $F$ | $T$ | $T$ | $T$ | $T$ |
| $F$ | $T$ | $T$ | $T$ | $T$ | $T$ |
| $T$ | $F$ | $F$ | $F$ | $F$ | $T$ |
| $T$ | $T$ | $T$ | $F$ | $T$ | $T$ |

Remark 3.11. The interesting tautologies and contradictions are not those involving only specific statements such as $T_{0}, F_{0}, A_{1}, A_{2}, B_{1}, B_{2}$, from above but those statements like $A_{5}, A_{6}, B_{4}$ and $B_{5}$ which specify formulas relating the general statements $P, Q$ and $R$.

### 3.4 Statement Equivalences (Understand this!)

Symbolic logic has a collection of very useful statement equivalences which are given here. They were taken from ch. 2 on logic, subchapter 2.4 (Important Logical Equivalences) of [4] Bryant, Kirby Course Notes for MAD 2104.

Theorem 3.4. Let $P, Q, R$ be statements.
a. Identity Laws:
b. Domination Laws:
c. Idempotent Laws:
d. Double Negation Law:
e. Commutative Laws:
f. Associative Laws:


$$
\begin{gathered}
P \wedge T_{0} \Leftrightarrow P \\
P \vee F_{0} \Leftrightarrow P \\
P \vee T_{0} \Leftrightarrow T_{0} \\
P \wedge F_{0} \Leftrightarrow F_{0} \\
P \vee P \Leftrightarrow P \\
P \wedge P \Leftrightarrow P \\
\neg(\neg P) \Leftrightarrow P \\
P \vee Q \Leftrightarrow Q \vee P \\
P \wedge Q \Leftrightarrow Q \wedge P \\
\text { hence }(P \vee Q) \vee R \Leftrightarrow P \vee Q \vee R \\
(P \wedge Q) \wedge R \Leftrightarrow P \wedge(Q \wedge R) \\
\text { hence }(P \wedge Q) \wedge R \Leftrightarrow P \wedge Q \wedge R
\end{gathered}
$$

g. Distributive Laws:
h. De Morgan's Laws: ${ }^{24}$
i. Absorption Laws:
j. Implication Law:

$$
\begin{align*}
& P \vee(Q \wedge R) \Leftrightarrow(P \vee Q) \wedge(P \vee R)  \tag{3.34}\\
& P \wedge(Q \vee R) \Leftrightarrow(P \wedge Q) \vee(P \wedge R) \tag{3.35}
\end{align*}
$$

You should remember this formula because the fact that implication can be expressed as an OR statement is often extremely useful when showing that two statements are logically equivalent.
k. Contrapositive Laws:
l. Tautology:
m. Contradiction:
n. Equivalence:

$$
\begin{equation*}
(P \rightarrow Q) \wedge(Q \rightarrow P) \Leftrightarrow(P \leftrightarrow Q) \tag{3.45}
\end{equation*}
$$

The proof for only some of the laws stated above are given here. You can prove all others by writing out the truth tables to show that left and right sides of the $\ldots \Leftrightarrow \ldots$ statements are indeed logically equivalent. PROOF of $\boldsymbol{h}$ (De Morgan's laws):
See theorem 3.3 on p. 42.
PROOF of $\boldsymbol{j}$ (implication law):
We prove (3.40) using a truth table:

[^19]We see that the entries T-T-F-T in the $\neg P \vee Q$ column match those given for $P \rightarrow Q$ in def.3.15 on p. 37 of the arrow operator. This proves the logical equivalence of those statements.

| $P$ | $Q$ | $\neg P$ | $\neg P \vee Q$ |
| :---: | :---: | :---: | :---: |
| $F$ | $F$ | $T$ | $T$ |
| $F$ | $T$ | $T$ | $T$ |
| $T$ | $F$ | $F$ | $F$ |
| $T$ | $T$ | $F$ | $T$ |

PROOF of $\boldsymbol{k}$ (contrapositive law for $\rightarrow$ ):
We prove (3.41) with the help of the previously given laws a through $\boldsymbol{j}$ :

$$
(P \rightarrow Q) \stackrel{(j)}{\Leftrightarrow}(\neg P \vee Q) \stackrel{(e)}{\Leftrightarrow}(Q \vee \neg P) \stackrel{(d)}{\Leftrightarrow}(\neg(\neg Q) \vee \neg P) \stackrel{(j)}{\Leftrightarrow}(\neg Q \rightarrow \neg P)
$$

Example 3.22. Use the logical equivalences of thm.3.4 to prove that $\neg(\neg A \wedge(A \wedge B))$ is a tautology.

## Solution:

$$
\begin{array}{lll} 
& \neg(\neg A \wedge(A \wedge B)) & \\
\Leftrightarrow & \neg(\neg A) \vee \neg(A \wedge B) & \text { De Morgan's Law (3.36) } \\
\Leftrightarrow & A \vee(\neg A \vee \neg B) & \text { De Morgan (3.36) + Double negation (3.27) } \\
\Leftrightarrow & (A \vee \neg A) \vee \neg B & \text { Associative law (3.30) } \\
\Leftrightarrow & T_{0} \vee \neg B & \text { Tautology (3.43) } \\
\Leftrightarrow & T_{0} & \text { Commutative Law (3.28) + Domination Law (3.23) }
\end{array}
$$

Example 3.23. Find a simple expression for the negation of the statement "if you come before 6:00 then I'll take you to the movies".

Solution: Let $A:=$ "You come before 6:00" and B $:=$ "I'll take you to the movies". Our task is to find a simple logical equivalent to $\neg(A \rightarrow B)$. We proceed as follows:

$$
\neg(A \rightarrow B) \stackrel{(j)}{\Leftrightarrow} \neg(\neg A \vee B) \stackrel{(h)}{\Leftrightarrow}(\neg(\neg A) \wedge \neg B) \stackrel{(d)}{\Leftrightarrow}(A \wedge \neg B)
$$

This translates into the statement "you come before 6:00 and I won't take you to the movies".

Remark 3.12. Now that we accept that such logical expressions are DEFINED by their truth tables, we must accept the following: if two logical expressions with two statements A and B as input have the same truth table, then they are logically equivalent and we may interchangeably use one or the other in a proof.

### 3.5 The Connection Between Formulas for Statements and for Sets (Understand this!)

Given statements $a, b$ and sets $A, B$ you may have the impression that there are connections between $a \wedge b$ and $A \cap B$, between $a \vee b$ and $A \cup B$, between $\neg a$ and $A^{\complement}$, etc. We will briefly explore this.

In this chapter we switch to small letters for statements and statement functions and use capital letters to denote sets. You have already seen an example in the introduction.

We assume the existence of a universal set $\mathscr{U}$ of which all sets are subsets.
All statements will be of the form $a(x)=$ " $x \in A$ " for some set $A \subset \mathscr{U}$. In other words we associate with such a set $A$ the following statement function:

$$
\begin{equation*}
a: \mathscr{U} \rightarrow \mathscr{S}, \quad x \mapsto a(x)=: " x \in A^{\prime \prime} \tag{3.46}
\end{equation*}
$$

This relationship establishes a correspondence between the subset $A$ of $\mathscr{U}$ and the predicate $a=$ " $x \in A$ " with UoD $\mathscr{U}$. We write $a \cong A$ for this correspondence.

Example 3.24. Let $a \cong A$ and $b \cong B$.
We have
a. $T_{0} \cong \mathscr{U}, F_{0} \cong \emptyset$
b. $\neg a: x \mapsto \neg a(x)=\neg " x \in A^{\prime \prime}$ evaluates to a true statement if and only if $x \notin A$, i.e. $x \in A^{\complement}$. Hence $\neg a \cong A^{\complement}$.
c. $a \wedge b: x \mapsto a(x) \wedge b(x)=" x \in A$ and $x \in B^{\prime \prime}$ evaluates to a true statement if and only if $x \in A \cap B$. Hence $a \wedge b \cong A \cap B$.
d. $a \vee b: x \mapsto a(x) \vee b(x)=" x \in A$ or $x \in B^{\prime \prime}$ evaluates to a true statement if and only if $x \in A \cup B$. Hence $a \vee b \cong A \cup B$.

We expand the table of formulas for statements given in thm 3.4 on p. 45 of ch.3.4 (Statement equivalences) with a third column which shows the corresponding relation for sets. Having a translation of statement relations to set relations allows you to use Venn diagrams as a visualization aid.

Theorem 3.5. For a set $\mathscr{U}$ Let $p, q, r$ be statement functions and let $P, Q, R \subseteq \mathscr{U}$ such that $p \cong P, q \cong Q$, $r \cong R$. Then we have the following:
a. Identity:
(3.48) $p \vee F_{0} \Leftrightarrow p$
(3.49) $p \vee T_{0} \Leftrightarrow T_{0}$
$P \cup \mathscr{U}=\mathscr{U}$
(3.50) $p \wedge F_{0} \Leftrightarrow F_{0}$
(3.51) $\quad p \vee p \Leftrightarrow p$
$p \wedge p \Leftrightarrow p \quad P \cap P=P$

$$
\begin{equation*}
\neg(\neg p) \Leftrightarrow p \tag{3.53}
\end{equation*}
$$

d. Double Negation:

$$
\left(P^{\complement}\right)^{\complement}=P
$$

$$
\begin{array}{ll}
p \vee q \Leftrightarrow q \vee p & P \cup Q=Q \cup P \\
p \wedge q \Leftrightarrow q \wedge p & P \cap Q=Q \cap P
\end{array}
$$

f. Associative:

$$
\begin{align*}
(p & \vee q) & \vee r & (P \cup Q) \cup R  \tag{3.56}\\
& \Leftrightarrow p \vee(q \vee r) & & \\
(p & \wedge q) & \wedge r & (P \cap Q) \cap R=P \cap(Q \cap R) \\
& \Leftrightarrow p \wedge(q \wedge r) & & \tag{3.57}
\end{align*}
$$

$$
P \cup(Q \cap R)=(P \cup Q) \cap(P \cup R)
$$

g. Distributive:

$$
\begin{aligned}
& p \vee(q \wedge r) \\
& \quad \Leftrightarrow(p \vee q) \wedge(p \vee r)
\end{aligned}
$$

$$
p \wedge(q \vee r)
$$

$$
P \cap(Q \cup R)=(P \cap Q) \cup(P \cap R)
$$

$$
\Leftrightarrow(p \wedge q) \vee(p \wedge r)
$$

h. De Morgan:

$$
\begin{align*}
& \neg(p \wedge q) \Leftrightarrow \neg p \vee \neg q  \tag{3.60}\\
& \neg(p \vee q) \Leftrightarrow \neg p \wedge \neg q \tag{3.61}
\end{align*}
$$

$$
(P \cap Q)^{\complement}=P^{\complement} \cup Q^{\complement}
$$

$$
(P \cup Q)^{\complement}=P^{\complement} \cap Q^{\complement}
$$

i. Absorption:

$$
\begin{array}{ll}
p \wedge(p \vee q) \Leftrightarrow p & P \cap(P \cup Q)=P \\
p \vee(p \wedge q) \Leftrightarrow p & P \cup(P \cap Q)=P
\end{array}
$$

Interpretation: $p(x) \rightarrow q(x)$, i.e., " $x \in P^{\prime \prime} \rightarrow$ " $x \in Q^{\prime \prime}$ is true if and only if $p(x): T, q(x):$ : is L/I., i.e., if and only if $x \notin P \cap Q^{\complement}=$ $P \backslash Q$, i.e., $x \in(P \backslash Q)^{\complement}$.

$$
\begin{equation*}
p \Rightarrow q \tag{3.65}
\end{equation*}
$$

j2. Implication 2:
k. Contrapositive:

$$
P \backslash Q=\emptyset \text {, i.e., } P \subseteq Q
$$

Note that we are not dealing with $p \rightarrow q$ but with $p \Rightarrow q$ where we assume for all $x$ a relation between $p$ and $q$ which renders $p(x): T, q(x):$ F logically impossible.
(3.67) $\quad(P \Rightarrow Q) \Leftrightarrow(\neg Q \Rightarrow \neg P)$

$$
\begin{equation*}
P^{\complement} \cup Q=Q \cup P^{\complement} \tag{3.66}
\end{equation*}
$$

$P \subseteq Q \Leftrightarrow Q^{\complement} \subseteq P^{\complement}$
l. Tautology:

$$
\begin{equation*}
(P \vee \neg P) \Leftrightarrow T_{0} \tag{3.68}
\end{equation*}
$$

$$
P \cup P^{\complement}=\mathscr{U}
$$

m. Contradiction:

$$
\begin{equation*}
(P \wedge \neg P) \Leftrightarrow F_{0} \tag{3.69}
\end{equation*}
$$

$$
P \cap P^{\complement}=\emptyset
$$

n1. Equivalence 1:

$$
\begin{array}{rc}
(p \rightarrow q) \wedge(q \rightarrow p) & \left(P^{\complement} \cup Q\right) \cap\left(Q^{\complement} \cup P\right) \\
\Leftrightarrow(p \leftrightarrow q) & =\{x: x \text { both in } P, Q \text { or }  \tag{3.70}\\
& x \text { neither in } P \text { nor in } Q\}
\end{array}
$$

n2. Equivalence 2:

$$
\begin{array}{rr}
(p \Rightarrow q) \wedge(q \Rightarrow p) & (P \subseteq Q) \text { and }(Q \subseteq P)  \tag{3.71}\\
\Leftrightarrow(p \Leftrightarrow q) & \Leftrightarrow(P=Q)
\end{array}
$$

PROOF: The set equalities are evident except for the following:
PROOF of Equivalence 1:

$$
\begin{aligned}
\left(P^{\complement} \cup Q\right) \cap\left(Q^{\complement} \cup P\right) & =\left[\left(P^{\complement} \cup Q\right) \cap Q^{\complement}\right] \cup\left[\left(P^{\complement} \cup Q\right) \cap P\right] \\
& =\left(P^{\complement} \cap Q^{\complement}\right) \cup\left(Q \cap Q^{\complement}\right) \cup\left(P^{\complement} \cap P\right) \cup(Q \cap P) \\
& =\left(P^{\complement} \cap Q^{\complement}\right) \cup(Q \cap P) \\
& =\{x: x \text { neither in } P \text { nor in } Q \text { or } x \text { both in } P, Q\} .
\end{aligned}
$$

### 3.6 Quantifiers for Statement Functions

This chapter has been kept rather brief. You can find more about quantifiers in ch. 2 on logic, subchapter ch.2.3 (Predicates and Quantifiers) of [4] Bryant, Kirby Course Notes for MAD 2104.

### 3.6.1 Quantifiers for One-Variable Statement Functions

Definition 3.18 (Quantifiers). Let $A: \mathscr{U} \rightarrow \mathscr{S}, x \mapsto A(x)$ be a statement function of a single variable $x$ with $\operatorname{UoD} \mathscr{U}$ for $x$.
a. The universal quantification of the predicate $A$ is the statement

$$
\begin{equation*}
\text { "For all } x A(x)^{\prime} \text { ", written } \forall x A(x) \text {. } \tag{3.72}
\end{equation*}
$$

The above is a short for " $A(x)$ is true for each $x \in \mathscr{U}$ ". We call the symbol $\forall$ the universal quantifier symbol.
b. The existential quantification of the predicate $A$ is the statement
"For some $x A(x)^{\prime \prime}$, written $\exists x A(x)$.
The above is a short for "There exists $x \in \mathscr{U}$ such that $A(x)$ is true". ${ }^{25}$ We call the symbol $\exists$ the existential quantifier symbol.

[^20]c. The unique existential quantification of the predicate $A$ is the statement
"There exists unique $x$ such that $A(x)^{\prime \prime}$, written $\exists!x A(x)$.
The above is a short for "There exists a unique $x \in \mathscr{U}$ such that $A(x)$ is true". ${ }^{26}$ We call the symbol $\exists$ ! the unique existential quantifier symbol.

Example 3.25. Let $A:[-3,3] \rightarrow \mathscr{S}$ be the statement function $x \mapsto$ " $x^{2}-4=0$ ".
Let $C:=\forall x A(x) D:=\exists x A(x)$ and $E:=\exists!x A(x)$. Then
$C=$ "for all $x \in[-3,3]$ it is true that $x^{2}-4=0$ "
$D=$ "there is at least one $x \in[-3,3]$ such that $x^{2}-4=0$ "
$E=$ "there is exactly one $x \in[-3,3]$ such that $x^{2}-4=0$ "
Note that each of $C, D, E$ is in fact a statement because each one is either true or false: Clearly the zeroes of the function $f(x)=x^{2}-4$ in the interval $-3 \leqq x \leqq 3$ are $x= \pm 2$. It follows that $D$ is a true statement and $A$ and $C$ are false statements.

Example 3.26. Let $\mathscr{U}:=\{$ all human beings $\}$ be the UoD for the following three predicates:

$$
\begin{aligned}
& S(x):=\text { " } x \text { is a student at NYU", } \\
& C(x):=\text { " } x \text { cheats when taking tests", } \\
& H(x):=\text { " } x \text { is honest", }
\end{aligned}
$$

Let us translate the following three english verbiage statements into formulas:

$$
\begin{aligned}
& A_{1}:=\text { "All humans are NYU students", } \\
& A_{2}:=\text { "All NYU students cheat on tests", } \\
& A_{3}:=\text { "Any NYU student who cheats on tests is not honest". }
\end{aligned}
$$

Solution:

$$
\begin{aligned}
& A_{1}=\forall x S(x), \\
& A_{2}=\forall x[S(x) \rightarrow C(x)], \\
& A_{3}=\forall x[(S(x) \wedge C(x)) \rightarrow \neg H(x)] .
\end{aligned}
$$

## Example 3.27. We continue example 3.26.

Let us simplify $A_{3}=\forall x[(S(x) \wedge C(x)) \rightarrow \neg H(x)]$.
It is clear that " $A(x)$ is true for all $x$ " is equivalent to "There is no $x$ such that $A(x)$ is false". In other words, we have for any statement function $A$ the following:

$$
\forall x A(x) \Leftrightarrow \neg[\exists x(\neg A(x))] .
$$

But $A_{3}$ is the form $\forall x A(x)$ : replace $A(x)$ with $(S(x) \wedge C(x)) \rightarrow \neg H(x)$.
It follows that

$$
\left.A_{3} \Leftrightarrow \neg[\exists x(\neg(S(x) \wedge C(x)) \rightarrow \neg H(x)))\right] .
$$

What a mess! let us drop the " $(x)$ " everywhere and the above becomes

$$
\left.A_{3} \Leftrightarrow \neg[\exists x(\neg(S \wedge C) \rightarrow \neg H))\right] .
$$

We have seen in example 3.23 on p. 47 that for any two statements $P$ and $Q$ the equivalence $\neg(P \rightarrow Q) \Leftrightarrow(P \wedge \neg Q)$ is true.

[^21]Let us apply this with $P:=S \wedge C$ and $Q:=\neg H$. We obtain

$$
A_{3} \Leftrightarrow \neg[\exists x((S \wedge C) \wedge \neg(\neg H))] . \Leftrightarrow \neg[\exists x(S \wedge C \wedge H)] .
$$

where we obtained the last equivalence by applying the double negation law to $\neg(\neg H)$ and the associative law for $\wedge$ to remove the parentheses from $(S \wedge C) \wedge H$.

As a last step we bring back the " $(x)$ " terms and obtain

$$
A_{3} \Leftrightarrow \neg \exists x[S(x) \wedge C(x) \wedge H(x)] .
$$

In other words, $A_{3}$ means "There is no one who is an NYU student and who cheats on tests and is honest". This should make sense if you remember the original meaning of $A_{3}$ : "Any NYU student who cheats on tests is not honest".

### 3.6.2 Quantifiers for Two-Variable Statement Functions

We now discuss quantifiers for statement functions of two variables. Things become a lot more interesting because we can mix up $\forall, \exists$ and $\exists$ !.

Unless mentioned otherwise $B$ denotes the statement function of two variables

$$
\begin{equation*}
B: \mathscr{U}_{x} \times \mathscr{U}_{y} \rightarrow \mathscr{S}, \quad x \mapsto B(x, y) \tag{3.75}
\end{equation*}
$$

It follows that the unverses of discourse are $\mathscr{U}_{x}$ for $x$ and $\mathscr{U}_{y}$ for $y$.
We need a quantifier for each variable to bind the expression $B(x, y)$ with placeholders $x$ and $y$ into a statement, i.e., into something that will be true or false. This done by example as follows:

Definition 3.19 (Doubly quantified expressions). Here is a table of statements involving two quantifiers and their meanings.
a. $\quad \forall x \forall y B(x, y) \quad$ "for all $x \in \mathscr{U}_{x}$ and for all $y \in \mathscr{U}_{y}$ (we have the truth of) $B(x, y)$ ",
b. $\forall x \exists y B(x, y) \quad$ "for all $x \in \mathscr{U}_{x}$ there exists (at least one) $y \in \mathscr{U}_{y}$ such that $B(x, y)$ ",
c. $\exists x \forall y B(x, y)$ "there exists (at least one) $x \in \mathscr{U}_{x}$ such that for all $y \in \mathscr{U}_{y} B(x, y)$ ",
d. $\exists!x \forall y B(x, y)$ "there exists exactly one $x \in \mathscr{U}_{x}$ such that for all $y \in \mathscr{U}_{y} B(x, y)$ ",
e. $\exists x \exists y B(x, y)$ "there exists (at least one) $x \in \mathscr{U}_{x}$ and (at least one) $y \in \mathscr{U}_{y}$ such that $B(x, y)$ ".

Example 3.28. Let $\mathscr{U}_{x}:=\mathbb{N}, \mathscr{U}_{y}:=\mathbb{Z}$ and $B: \mathscr{U}_{x} \times \mathscr{U}_{y} \rightarrow \mathscr{S}, \quad(x, y) \mapsto B(x, y):=" x+y=1$ ". Then
a. $\forall x \forall y B(x, y) \quad$ false
b. $\quad \forall x \exists y B(x, y) \quad$ true: for the given $x$ choose $y:=1-x$.
c. $\exists y \forall x B(x, y)$ false
d. $\forall y \exists x B(x, y) \quad$ false: If you choose $y>0$ then the only $x$ that satisfies the equation $x+y=1$ is $x=1-y \leq 0$, i.e., $x \notin \mathbb{N}$, the $\operatorname{UoD}$ for $x$.
e. $\exists!x \forall y B(x, y)$ false
f. $\quad \exists x \exists y B(x, y) \quad$ true: choose $x:=10$ and $y:=-9$.

Understand the different outcomes of $\mathbf{b}, \mathbf{c}$ and $\mathbf{d}$ and remember this:
a. The order in which the qualifiers are applied is important. $\forall x \exists y$ generally does not mean the same as $\exists y \forall x$.
b. Interchanging variable names in the qualifiers is not OK. $\forall x \exists y$ generally does not mean the same as $\forall y \exists x$.

Proposition 3.1. Note the following:

$$
\begin{align*}
& \forall x \forall y B(x, y) \Leftrightarrow \forall y \forall x B(x, y)  \tag{3.76}\\
& \exists x \exists y B(x, y) \Leftrightarrow \exists y \exists x B(x, y)  \tag{3.77}\\
& \forall x \exists y B(x, y) \Leftrightarrow \exists y \forall x B(x, y)  \tag{3.78}\\
& \exists y \forall x B(x, y) \Rightarrow \forall x \exists y B(x, y) \tag{3.79}
\end{align*}
$$

PROOF: (3.76) and (3.77) follow from $\boldsymbol{a}$ and $\boldsymbol{e}$ in def. 3.19 and we saw an example for (3.76) in the previous example.

The last item is not so obvious. We argue as follows: Assume that $\exists y \forall x B(x, y)$ is true. Then there is some $y_{0} \in \mathscr{U}_{y}$ such that $B\left(x, y_{0}\right.$ is true for all $x \in \mathscr{U}_{x}$.

Why does that imply the truth of $\forall x \exists y B(x, y)$, i.e., for all $x \in \mathscr{U}_{x}$ you can pick some $y \in \mathscr{U}_{y}$ such that $B(x, y)$ is true? Here is the answer: Pick $y_{0}$. This works because, by assumption, $B\left(x, y_{0}\right)$ is true for all $x \in \mathscr{U}_{x}$.

Remark 3.13. The last part of the proof of (3.79) is worth a closer look:
" $\forall x \exists y \ldots$ " only tells you that for all $x$ there will be some $y$ which generally depends on $x$, something we sometimes emphasize using "functional notation" $y=y(x)$.
" $\exists y \forall x \ldots$ " does more: it postulates the existence of some $y_{0}$ which is suitable for each $x$ in its UoD. The assignment $y(x)=y_{0}$ is constant in $x$ !

Remark $\mathbf{3 . 1 4}$ (Partially quantified statement functions). Given a statement function

$$
B: \mathscr{U}_{x} \times \mathscr{U}_{y} \rightarrow \mathscr{S}, \quad x \mapsto B(x, y)
$$

with two place holders $x$ and $y$, we can elect to use only one quantifier for either $x$ or $y$. If we only quantify $x$ then we only bind $x$ and $y$ still remains a placeholder and if we only quantify $y$ then we only bind $y$ and $x$ still remains a placeholder.

Example 3.29. Let $\mathscr{U}_{x}:=\{$ all students at this party $\}$ and $\mathscr{U}_{y}:=\{$ "Linear Algebra", Discrete Mathematics", "Multivariable Calculus", "Ordinary Differential Equations", "Complex Variables", "Graph Theory", "Real Analysis" \}.

Let $A:=$ " $x$ studies $y$ " be the two-variable statement function with $\operatorname{UoD} \mathscr{U}_{x}$ for $x$ and $\operatorname{UoD} \mathscr{U}_{y}$ for $y$, i.e.,
$A: \mathscr{U}_{x} \times \mathscr{U}_{y} \rightarrow \mathscr{S}, \quad(x, y) \mapsto A(x, y)=" x$ studies $y "$.
Then $B:=\forall x A(x, y)$ is the one-variable predicate
$B: \mathscr{U}_{y} \rightarrow \mathscr{S}, \quad y \mapsto B(y)=$ "all students at this party study $y$ "
and $C:=\exists!y A(x, y)$ is the one-variable predicate
$C: \mathscr{U}_{x} \rightarrow \mathscr{S}, \quad x \mapsto C(x)=$ " $x$ studies exactly one of the courses listed in $\mathscr{U}_{y} "$.

### 3.6.3 Quantifiers for Statement Functions of more than Two Variables

Remark 3.15. Although this document limits its scope to statement functions of one or two variables (see the note before remark 3.6 in ch.3.2 (Statements and statement functions)) we discuss briefly
the use of quantifiers for predicates

$$
A: \mathscr{U}_{1} \times \mathscr{U}_{2} \times \cdots \times \mathscr{U}_{n} \rightarrow \mathscr{S}, \quad\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto A\left(x_{1}, x_{2}, \ldots, x_{n}\right) .
$$

with $n$ place holders.
Each one of those variables needs to be bound by one of the quantifiers $\forall, \exists, \exists$ ! in order to obtain a statement, i.e., something that is either true or false.
Example 3.30 (Continuity vs uniform continuity). This example demonstrates the effect of switching a $\forall$ quantifier with an $\exists$ quantifier for a predicate with four variables. You will learn later that one quantification corresponds to ordinary continuity and the other corresponds to uniform continuity of a function. Do not worry if you do not understand how this example relates to continuity. The only point of interest here is the use of the quantifiers.

Let $a<b$ be two real numbers and let $f:] a, b[\rightarrow \mathbb{R}$ be a function which maps each $x$ in its domain $] a, b$ to a real number $y=f(x)$.

Let $\left.\mathscr{U}_{\varepsilon}:=\mathscr{U}_{\delta}:=\right] 0, \infty\left[\right.$ and $\left.\mathscr{U}_{x}:=\mathscr{U}_{x^{\prime}}:=\right] a, b\left[\right.$. Let $P: \mathscr{U}_{x} \times \mathscr{U}_{x^{\prime}} \times \mathscr{U}_{\delta} \times \mathscr{U}_{\varepsilon} \rightarrow \mathscr{S}$ be the predicate

$$
\left(x, x^{\prime}, \delta, \varepsilon\right) \mapsto P\left(x, x^{\prime}, \delta, \varepsilon\right):=\text { "if }\left|x-x^{\prime}\right|<\delta \text { then }\left|f(x)-f\left(x^{\prime}\right)\right|<\varepsilon^{"} .
$$

Let $A:=\forall \varepsilon \forall x \exists \delta \forall x^{\prime} P\left(x, x^{\prime}, \delta, \varepsilon\right)$ Then $A$ being true is equivalent to saying that the function $f$ is continuous at each point $x \in] a, b\left[.{ }^{27}\right.$

Let $B:=\forall \varepsilon \exists \delta \forall x \forall x^{\prime} P\left(x, x^{\prime}, \delta, \varepsilon\right)$. Then $B$ being true is equivalent to saying that the function $f$ is uniformly continuous in $] a, b\left[.{ }^{28}\right.$

The difference between $A$ and $B$ is that in statement $A$ the variable $\delta$ whose existence is required may depend on both $\varepsilon$ and $x$, i.e., $\delta=\delta(\varepsilon, x$

On the other hand, to satisfy $B$, a $\delta$ must be found which still may depend on $\varepsilon$ but it must be suitable for all $x \in] a, b[$, i.e., $\delta=\delta(\varepsilon$.
Remark 3.16 (Partially quantified statement functions). What was said in remark 3.14 about partial qualification of two-variable predicates generalizes to more than two variables: If $A$ is a statement function with $n$ variables and we use quantifiers for only $m<n$ of those variables then $n-m$ variables in the resulting expression remain unbound and this expression becomes a statement function of those unbound variables.

For example, if $A(w, x, y, z)$ is a four-variable predicate then $B:(x, z) \mapsto[\forall y \neg \exists w A(w, x, y, z)]$ defines a two-variable predicate $B$ which inherits the UoDs for $x$ and $z$ from the original statement function $A$.

### 3.6.4 Quantifiers and Negation (Understand this!)

Negation of statements involving quantifiers is governed by
Theorem 3.6 (De Morgan's laws for quantifiers). Let $A$ be a statement function with $\mathrm{UoD} \mathscr{U}$. Then

$$
\begin{aligned}
& \text { a. } \neg(\forall x A(x)) \Leftrightarrow \exists x \neg A(x) \quad \begin{array}{l}
\text { "It is not true that } A(x) \text { is true for all } x " \Leftrightarrow \text { "There is } \\
\text { some } x \text { for which } A(x) \text { is not true" }
\end{array} \\
& \text { b. } \neg(\exists x A(x)) \Leftrightarrow \forall x \neg A(x) \begin{array}{l}
\text { "There is no } x \text { for which } A(x) \text { is true" } \Leftrightarrow \text { " } A(x) \text { is not } \\
\text { true for all } x \text { " }
\end{array}
\end{aligned}
$$

[^22]PROOF of a: Not given here but you can find it in ch. 2 on logic, subchapter 3.11 (De Morgan's Laws for Quantifiers) of [4] Bryant, Kirby Course Notes for MAD 2104.
PROOF of b: Let $\mathscr{U}_{x}$ be the UoD for $x$.
The truth of $\neg(\exists x A(x))$ means that $\exists x A(x)$ is false, i.e., $A(x)$ is false for all $x \in \mathscr{U}_{x}$. This is equivalent to stating that $\neg A(x)$ is true for all $x \in \mathscr{U}_{x}$ and this is by definition, the truth of $\forall x \neg A(x)$.

You can use the formulas above for negation of statements of more than one variable with more than one quantifier using the following method, demonstrated here by example.

Example 3.31. Negate the statement $\exists x \forall y P(x, y)$, i.e., move the $\neg$ operator of $\neg \exists x \forall y P(x, y)$ to the right past all quantifiers.

The key is to introduce an intermittent predicate $A: x \mapsto A(x):=[\forall y P(x, y)]$. We obtain

$$
\begin{aligned}
{[\neg \exists x \forall y P(x, y)] \Leftrightarrow[\neg \exists x A(x)] \stackrel{(\mathbf{b})}{\Leftrightarrow}[\forall x \neg A(x)] } & \Leftrightarrow[\forall x(\neg \forall y P(x, y))] \\
& \stackrel{\text { a) }}{\Leftrightarrow}[\forall x(\exists y \neg P(x, y))] .
\end{aligned}
$$

Example 3.32. As in example 3.31, negate the statement $\exists x \forall y P(x, y)$ but do so using parentheses instead of explicitly defining an intermittent predicate.

Here is the solution:

$$
\begin{aligned}
{[\neg \exists x \forall y P(x, y)] \Leftrightarrow[\neg \exists x(\forall y P(x, y))] \stackrel{(\mathbf{b})}{\Leftrightarrow}[\forall x \neg(\forall y P(x, y))] } & \Leftrightarrow[\forall x(\neg \forall y P(x, y))] \\
& \stackrel{(a)}{\Leftrightarrow}[\forall x(\exists y \neg P(x, y))] .
\end{aligned}
$$

### 3.7 Proofs (Understand this!)

We have informally discussed proofs in examples 3.18 and 3.19 of chapter 3.3.5 (Arrow and Implication Operators) on p. 37 and seen in two simple cases how a proof can be done by building a single truth table for an if...then statement and showing that it is a tautology. In this chapter we take a deeper look at the concept of "proof".

Many subjects discussed here follow closely ch. 3 (Methods of Proofs) of [4] Bryant, Kirby Course Notes for MAD 2104.

### 3.7.1 Building Blocks of Mathematical Theories

Some of the terminology definitions in notations 3.2 and 3.4 were taken almost literally from ch. 3 (Methods of Proofs), subchapter 1 (Logical Arguments and Formal Proofs) of [4] Bryant, Kirby Course Notes for MAD 2104.

Notations 3.2 (Axioms, rules of inferences and assertions).
a. An axiom is a statement that is true by definition. No justification such as a proof needs to be given.
b. A rule of inference is a logical rule that is used to deduce the truth of a statement from the truth of others.
c. For some statements it is not clear whether they are true for false. Even if a statement is known to be true there might be someone like a student taking a test who is given the task to demonstrate, i.e., prove its truth. In this context we call a statement an assertion and we call it a valid assertion if it can be shown to be true. An assertion which is not known to be true by anyone is often called a conjecture.

Example 3.33. Let $A:=$ "all continuous functions are differentiable" (known to be false ${ }^{29}$ ) and $B:=$ "all differentiable functions are continuous" (known to be true). A homework problem in calculus may ask the students to figure out which of the four statements $A, \neg A, B, \neg B$ are valid assertions and give proofs to that effect.

## Remark 3.17.

a. Goldbach's conjecture states that every even integer greater than 2 can be expressed as the sum of two primes, i.e., integers $p$ greater than 1 which can be divided evenly by no natural number other than $p(p / p=1)$ or $1(p / 1=p)$. Goldbach came up with this in 1742 , more than 250 years ago. No one has been able until now to either prove the validity of this assertion or provide a counterexample to prove its falsehood.
b. Fermat's conjecture was that there are no four numbers $a, b, c, n \in \mathbb{N}$ such that $n>2$ and $a^{n}+b^{n}=c^{n} .{ }^{30}$ This was stated by Pierre de Fermat in 1637 who then claimed that he had a proof. Unfortunately he never got around to write it down. A successful proof was finally published in 1994 by Andrew Wiles. Accordingly, Fermat's conjecture was rechristened Fermat's Last Theorem.

Notations 3.3 (Proofs). A proof is the demonstration that an assertion is valid. This demonstration must be detailed enough so that a person with sufficient expert knowledge can understand that we

[^23]do indeed have a statement which is true for all logically possible combinations of T/F values. To show that the arguments given in this demonstration are valid, available tools are
a. the rules of inference which wil be discussed in section 3.7.2 (Rules of Inference) on p.59
b. logical equivalences for statements (see ch.3.3.6 (Biconditional and Logical Equivalence Operators - Part 2) on p.ch.42).

In almost all cases the assertion in question is of the form "if $P$ then $C$ ". Proving it means showing that the statement $P \rightarrow C$ is a tautology, i.e., it can be replaced by the stronger $P \Rightarrow C$ statement. The proof then consists of the demonstration that the combination $P$ : true, $C$ : false can be ruled out as logically impossible. In other words, assuming $P$ : true, i.e., the truth of the premise, it must be shown that $C$ : true, i.e., the conclusion then also is necessarily true.

Usually a proof is broken down into several "sub-proofs" which can be proved separately and where some or all of those steps again will be broken down into several steps ... You can picture this as a hierarchical upside down tree with a single node at the top. At the most detailed level at the bottom we have the leaf nodes. The proof of the entire statement is represented by that top node.

Notations 3.4 (Theorems, lemmata and corollaries).
a. A theorem is an assertion that can be proved to be true using definitions, axioms, previously proven theorems, and rules of inference.
b. A lemma (plural: lemmata) is a theorem whose main importance is that it can used to prove other theorems.
c. A corollary is a theorem whose truth is a fairly easy consequence of another theorem.

Remark 3.18 (Terminology is different outside logic). The terminology given in the above definitions is specific to the subject of mathematical logic. In other branches of mathematics and hence outside this chapter 3 different meanings are attached to those terms:

Each one of lemma, proposition, theorem, corollary is a theorem as defined above in notations 3.2, i.e., a statement that can be proved to be true. We distinguish those terms by comparing them to propositions:
a. Theorems are considered more important than propositions.
b. The main purpose of a lemma is to serve as a tool to prove other propositions or theorems.
c. A corollary is a fairly easy consequence of some lemma, proposition, theorem or other corollary.

> It was mentioned as a footnote to the definition of a statement (def. 3.3 on $p .26$ ) that what we call a statement, [4] Bryant, Kirby calls a proposition and that we deviate from that approach because mathematics outside logic uses "proposition" to denote a theorem of lesser importance.

Any mathematical theory must start out with a collection of undefined terms and axioms that specify certain properties of those undefined terms.

There is no way to build a theory without undefined terms because the following will happen if you try to define every term: You define $T_{2}$ in terms of $T_{1}$, then you define $T_{3}$ in terms of $T_{2}$, etc. Two possibilities:

1. Each of $T_{1}, T_{2}, T_{3}, \ldots$ are different and you end up with an infinite sequence of definitions.
2. At least one of those terms is repeated and there will be a circular chain of definitions.

Neither case is acceptable if you want to specify the foundations of a mathematical system.
Example 3.34. Here are a few important examples of mathematical systems and their ingredients.
a. In Euclid's geometry of the plane some of the undefined terms are "point", "line segment" and "line". The five Euclidean axioms specify certain properties which relate those undefined terms. You may have heard of the fifth axiom, Euclid's parallel postulate. It has been reproduced here with small alterations from Wikipedia's "Euclidean geometry" entry: ${ }^{31}$
(It is postulated that) "if a line segment falling on two line segments makes the interior angles on the same side less than two right angles, the two line segments, if produced indefinitely, meet on that side on which are the angles less than the two right angles".
b. In the so called Zermelo-Fraenkel set theory which serves as the foundation for most of the math that has been done in the last 100 years, the concept of a "set" and the relation "is an element of" $(\in)$ are undefined terms.
c. Chapters 1 and 2 of [1] Beck/Geoghegan list several axioms which stipulate the existence of a nonempty set called $\mathbb{Z}$ whose elements are called "integers" which you can "add" and "multiply". Certain algebraic properties such as " $a+b=b+a$ " and " $c \cdot(a+b)=(c \cdot a)+(c \cdot a)$ " are given as true and so is the existence of an additive neutral unit " 0 " and a multiplicative neutral unit " 1 ". Besides those algebraic properties the existence of a strict subset $\mathbb{N}$ called "positive integers" is assumed which has, among others, the property that any $z \in \mathbb{Z}$ either satisfies $z \in \mathbb{Z}$ or $-z \in \mathbb{Z}$ or $z=0$. Finally there is the induction axiom which states that if you create the sequence $1,1+1,(1+1)+1$, $\ldots$ then you capture all of $\mathbb{N}$. This axiom is the basis for the principle of proof by mathematical induction (see def.2.12 on p. 19).

Once we have the undefined terms and axioms for a mathematical system, we can begin defining new terms and proving theorems (or lemmas, or corollaries) within the system.

Remark 3.19 (Axioms vs. definitions). You can define anything you want but if you are not careful you may have a logical contradiction and the set of all items that satisfy that definition is empty. In contrast, axioms will postulate the existence of an item or an entire collection of items which satisfy all axioms. If the axioms contradict each other we have a theory which is inconsistent and the only way to deal with it is to discard it and rework its foundations An example for this was set theory in its early stages. Anything that you could phrase as "Let $A$ be the set which contains ..." was fair game to define a set. We saw in remark 2.2 (Russell's Antinomy) on p. 13 that this lead to problems so serious that they caused some of the leading mathematicians of the time to revisit the foundations of mathematics.

Example 3.35. For example you can define an oddandeven integer to be any $z \in \mathbb{Z}$ which satisfies that $z-212$ is an even number and $z+48$ is an odd number and you can prove great things for such $z$. The problem is of course that the set of all oddandeven integers is empty! We have a definition which is useless for all practical purposes, but no mathematical harm is done.

[^24]On the other hand, if you add as an additional axiom for $\mathbb{Z}$ in example 3.34.c that $\mathbb{Z}$ must contain one or more oddandeven integers then you are in a conundrum because you postulated the existence of a set $\mathbb{Z}$ which satisfies all axioms and the existence of such a set is logically impossible!

### 3.7.2 Rules of Inference

Remark 3.20 (Most important rules of inference). In Notations 3.2 on p. 56 we described the term "rule of inference" as "a logical rule that is used to deduce the truth of a statement from the truth of others". The most important rules of inference are those that allow you to draw a conclusion of the form "if $A$ is true then I am allowed to deduce the the truth of $C$." This basically amounts to having is a list of premises $A_{1}, A_{2}, \ldots, A_{n}$ and a conclusion $C$ such that

$$
\begin{equation*}
\text { the compound statement }\left[A_{1} \wedge A_{2} \wedge \cdots \wedge A_{n}\right] \rightarrow C \text { is a tautology. } \tag{3.80}
\end{equation*}
$$

In other words, the column for the conclusion $C$ in the truth table for this statement must have the value true for each combination of truth values which is not logically impossible.

Observe that the order of the premises does not matter because the and connective is commutative.

Theorem 3.7. Let $P_{1}, P_{2}, \ldots, P_{n}$ and $C$ be statements. Then the statement $\left(P_{1} \wedge P_{2} \wedge \cdots \wedge P_{n}\right) \rightarrow C$ is a tautology if and only if the following combination of truth values is logically impossible:

$$
\begin{equation*}
P_{j} \text { is true for each } j=1,2, \ldots, n \text { and } C \text { is false. } \tag{3.81}
\end{equation*}
$$

## PROOF:

Let $P:=\left(P_{1} \wedge P_{2} \wedge \cdots \wedge P_{n}\right)$. Then " $P_{j}$ is true for each $j=1,2, \ldots, n$ " means according to the definition of the $\wedge$ operator the same as the truth of $P$. Hence proving the theorem is equivalent to proving that the statement $P \rightarrow C$ is a tautology if and only if the combination of truth values

$$
\begin{equation*}
P \text { is true and } C \text { is false is logically impossible. } \tag{3.82}
\end{equation*}
$$

In other words, we must prove that $P \rightarrow C$ is a tautology if and only if the row with the combination P:T, C:F, i.e., row 3, is logically impossible and can be ignored. This is is obvious as row 3 is the only one for which $P \rightarrow C$ evaluates to false.

|  | $P$ | $C$ | $P \rightarrow C$ |
| :---: | :---: | :---: | :---: |
| 1. | $F$ | $F$ | $T$ |
| 2. | $F$ | $T$ | $T$ |
| 3. | $T$ | $F$ | false |
| 4. | $T$ | $T$ | $T$ |

Notations 3.5. Rules of inference are commonly written in the following form:


Read ".$\therefore$ " as "therefore". The following, more compact notation can also be found:

| $A_{1}, A_{2}, \ldots, A_{n}$ |
| :--- |
| $\therefore C$ |

Theorem 3.8 (The three most important inference rules). The following lists three inference rules, i.e., those arrow statements are indeed tautoloties:

| Modus Ponens | $A$ |
| :--- | :--- |
| (Law of detachment - the mode that af- <br> firms the antecedent (the premise)) | $A \rightarrow C$ |


| Modus Tollens <br> (The mode that Denies the consequent <br> (the conclusion)) | $\neg C$ |
| :--- | :--- |
|  | $A \rightarrow C$ |


|  | $A \rightarrow B$ |
| :--- | :--- |
| Hypothetical syllogism | $B \rightarrow C$ |
| $\therefore A \rightarrow C$ |  |

Here is the compact notation:

| Modus Ponens | Modus Tollens | Hypothetical syllogism |
| :--- | :--- | :--- |
| $A, A \rightarrow C$ | $\neg C, A \rightarrow C$ | $A \rightarrow B, B \rightarrow C$ |
| $\therefore C$ | $\therefore \neg A$ | $\therefore A \rightarrow C$ |

Note that the proof that the hypothetical syllogism is a tautology was given in thm.3.1 on $p .37$
PROOF:

Example 3.36. Here are five more inference rules.

Disjunction Introduction
A
$\therefore A \vee B$

Conjunction elimination
$A \wedge B$
$\therefore A$

|  | $A \vee B$ |
| :--- | :--- |
| Disjunctive syllogism | $\neg A$ |
|  | $\therefore B$ |

$$
\begin{array}{ll} 
& A \\
\text { Conjunction introduction } & B \\
& \therefore A \wedge B \\
\text { Constructive dilemma } & (A \rightarrow B) \wedge(C \rightarrow D) \\
& A \vee C \\
& \therefore B \vee D \tag{3.90}
\end{array}
$$

Compact notation:

| Disjunction Introduction | Conjunction elimination | Disjunctive syllogism |
| :---: | :--- | :--- |
| $A$ $A \wedge B$ | $A \vee B, \neg A$ <br> $\therefore A \vee B$ | $\therefore A$ |
| Conjunction introduction | Constructive dilemma |  |
| $A, B$ $(A \rightarrow B) \wedge(C \rightarrow D), A \vee C$ |  |  |
| $\therefore A \wedge B$ | - |  |

None of the rules of inference that were given in this chapter involve quantifiers. You can find information about that topic in ch.2, section 1.6 (Rules of Inference for Quantifiers) of [4] Bryant, Kirby Course Notes for MAD 2104.

### 3.7.3 An Example of a Direct Proof

We illustrate in detail a mathematical proof by applying some the tools you have learned so far in this chapter on logic. For an example we will prove the theorem that each polynomial is differentiable. We define a polynomial as a function $f(x)=\sum_{j=0}^{n} c_{j} x^{j}$ for some $n=0,1,2, \ldots$, i.e., for some $n \in \mathbb{Z}_{\geqq 0}$ and we write $\mathscr{D}$ for the set of all differentiable functions. We now can formulate our theorem.

Theorem 3.9. Given the statements
a: $\quad A:="\left(n \in \mathbb{Z}_{\geqq 0}\right) \wedge\left(c_{0} \in \mathbb{R}\right) \wedge\left(c_{1} \in \mathbb{R}\right) \wedge \cdots \wedge\left(c_{n} \in \mathbb{R}\right) \wedge\left(f(x)=\sum_{j=0}^{n} c_{j} x^{j}\right)^{\prime \prime}$,
b: $\quad B:=" f(x) \in \mathscr{D}^{\prime \prime}$,
the following is valid: $A \Rightarrow B .{ }^{32}$

## PROOF:

We first collect the necessary ingredients.
We define the following statements which serve as abbreviations so that the formulas we will build are reasonably compact.

$$
\begin{array}{ll}
\text { a: } & Z_{j}:=" j \in \mathbb{Z}_{\geq 0} ", \\
\text { b: } & C_{j}:=Z_{j} \wedge " c_{j} \in \mathbb{R}^{\prime \prime}, \\
c:{ }^{33} & X_{j}:=Z_{j} \wedge " x^{j} \in \mathscr{D}^{\prime \prime}, \\
\text { d: } & D_{j}:=Z_{j} \wedge " c_{j} x^{j} \in \mathscr{D}^{\prime \prime}, \\
\text { e: } & E:=Z_{n} \wedge " f(x)=\sum_{j=0}^{n} c_{j} x^{j "}, \\
f: & B:=" f(x) \in \mathscr{D}^{\prime \prime} \text { (repeated for convenient reference) }
\end{array}
$$

We now can write our theorem as

$$
\begin{equation*}
\left(Z_{n} \wedge C_{0} \wedge C_{1} \wedge \cdots \wedge C_{n} \wedge E\right) \rightarrow B \tag{3.91}
\end{equation*}
$$

We assume that the following three theorems were proved previously, hence we may use them without giving a proof.

Theorem Thm-1: If $p(x)$ is a power of $x$, i.e., $p(x)=x^{n}$ for some $n=0,1,2, \ldots$, then is $p(x)$ differentiable.

We rewrite Thm-1 as an implication which uses the statements above. Let

$$
A_{1}:=Z_{n} \wedge " p(x)=x^{n \prime \prime}, \quad B_{1}:=X_{n}
$$

Then Thm-1 states that $A_{1} \Rightarrow B_{1} .{ }^{34}$
Theorem Thm-2: The product of a constant (real number) and a differentiable function is differentiable.
We rewrite Thm-2 as an implication. Let

$$
\begin{aligned}
& A_{2}:=" c \in \mathbb{R}^{\prime \prime} \wedge " h(x) \in \mathscr{D}^{\prime \prime} \wedge " g(x)=c \cdot h(x)^{\prime \prime}, \\
& B_{2}:=" h(x) \in \mathscr{D}^{\prime \prime},
\end{aligned}
$$

Then Thm-2 states that $A_{2} \Rightarrow B_{2}$.
Theorem Thm-3: The sum of differentiable functions is differentiable

[^25]We rewrite Thm-3 as an implication. Let

$$
\begin{aligned}
& A_{3}:=" Z_{n} \wedge " h_{1}(x) \in \mathscr{D}^{\prime \prime} \wedge " h_{2}(x) \in \mathscr{D}^{\prime \prime} \wedge \cdots \wedge " h_{n}(x) \in \mathscr{D}^{\prime \prime} \wedge " g(x)=\sum_{j=0}^{n} h_{j}(x) ", \\
& B_{3}:=" g(x) \in \mathscr{D}^{\prime \prime},
\end{aligned}
$$

Then Thm-3 states that $A_{3} \Rightarrow B_{3}$.

| Assertion |  | Reason |
| :--- | :--- | :--- |
| $a:$ | $Z_{0}, Z_{1}, \ldots Z_{n}$ | evident from $\mathbb{Z}_{\geq 0}=\{0,1,2, \ldots\}$ |
| $b:$ | $C_{0}, C_{1}, \ldots C_{n}$ | part of the premise of $A \rightarrow B($ see (3.91)) |
| $c:$ | $Z_{j} \rightarrow X_{j}(j=0,1, \ldots n)$ | Thm- with $n:=j$ |
| $d:$ | $X_{j}(j=0,1, \ldots n)$ | c and modus ponens |
| $e:$ | $\left(Z_{j} \wedge C_{j} \wedge X_{j}\right) \rightarrow D_{j}(j=0,1, \ldots n)$ | Thm-2 with $c:=c_{j}$ and $h(x):=x^{j}$ |
| $f:$ | $D_{j}(j=0,1, \ldots n)$ | e and modus ponens |
| $\boldsymbol{g}:$ | $E$ | part of the premise of $A \rightarrow B$ |
| $h:$ | $\left(Z_{n} \wedge D_{0} \wedge D_{1} \wedge \cdots \wedge D_{n} \wedge E\right) \rightarrow B$ | $g$ and Thm-3 with $h_{j}(x):=c_{j} x^{j}$ and $g(x):=f(x)$ |
| $\boldsymbol{i}:$ | $B$ | $h$ and modus ponens |

We have demonstrated that the truth of the premise $A$ of our theorem implies that of its conclusion $B$ and this proves the theorem.

Remark 3.21. Let us reflect on the steps involved in the proof above.
a: Break down all statements involved - not only those in the theorem you want to prove but also in all theorems, axioms and definitions you reference - into reusable components and name those components with a symbol so that it is easier to understand what assertions you employ and how they lead to the truth of other assertions. Example: $D_{j}$ references the component $Z_{j} \wedge$ " $c_{j} x^{j} \in \mathscr{D}^{\prime \prime}$ (which itself references the component $\left.Z_{j}=" j \in \mathbb{Z}_{\geqq 0}^{\prime \prime}\right)$.
b: Rewrite the theorem to be proved as an implication $A \Rightarrow B$.
c: Do the same for the three other theorems that we assumed as already having been proved.
The following is specific to our example but can be modified to other problems.
d: Start by using the premise $A$ and the definition $\mathbb{Z}_{\geqq 0}:=\{0,1,2, \ldots\}$ to get the first two rows. Show that what you have implies the truth of the premise of Thm-1 and then use the modus ponens inference rule to deduce the truth of its conclusion $X_{j}$. This allows $X_{j}$ to become an additional assertion.
e: Use that new assertion to obtain the truth of the premise of Thm-2 and then use again modus ponens to deduce the truth of its conclusion $D_{j}$. Now $D_{j}$ becomes an additional assertion.
f: Use that new assertion to obtain the truth of the premise of Thm-3 and then use again modus ponens to deduce the truth of $D_{j}$. Now $D_{j}$ becomes an additional assertion.

### 3.7.4 Invalid Proofs Due to Faulty Arguments

Remark 3.22 (Fallacies in logical arguments). People who are not very analytical often commit the following errors in their argumentation:

Affirming the Consequent
(proving the wrong direction)

Denying the Antecedent
(indirect proof in the wrong direction)

$$
\begin{equation*}
P \rightarrow Q \tag{3.92}
\end{equation*}
$$

$Q$
$\therefore P$
$P \rightarrow Q$
$\neg P$
$\therefore \neg Q$

Circular Reasoning
The argument incorporates use of the (not yet proven) conclusion

The reason that the above are fallacies stems from the fact that the above "rules of inferences" are not tautologies.

Example 3.37 (Fallacies in reasoning). a. Affirming the Consequent:
"If you are a great mathematician then you can add $2+2$ ". It is true that you can add $2+2$. You conclude that you are a great mathematician.
b. Denying the Antecedent:
"If this animal is a cat then it can run quickly". This is not a cat. You conclude that this animal cannot run quickly.
c. Circular Reasoning: ${ }^{35}$
"If $x y$ is divisible by 5 then $x$ is divisible by 5 or $y$ is divisible by 5 ".
The following incorrect proof uses the yet to be proven fact that the factors can be divided evenly by 5 .

## PROOF:

If $x y$ is divisible by 5 then $x y=5 k$ for some $k \in \mathbb{Z}$. But then $x=5 m$ or $y=5 n$ for some $m, n \in \mathbb{Z}$ (this is the spot where the conclusion was used). Hence $x$ is divisible by 5 or $y$ is divisible by 5 .

### 3.8 Categorization of Proofs (Understand this!)

There are different methods by which you can attempt to prove an "if ... then" statement $P \Rightarrow Q$. They are:
a. Trivial proof
b. Vacuous proof
c. Direct proof
d. Proof by contrapositive
e. Indirect proof (proof by contradiction)
f. Proof by cases

[^26]
### 3.8.1 Trivial Proofs

The underlying principle of a trivial proof is the following: If we know that the conclusion $Q$ is true then any implication $P \Rightarrow Q$ is valid, regardless of the hypothesis $P$.
Example 3.38 (Trivial proof). Prove that if it rains at least 60 days per year in Miami then $25+35=$ 60.

PROOF: There is nothing to prove as it is known that $25+35=60$. It is irrelevant whether or not in rains (or snows, if you prefer) 60 days per year in Miami.

### 3.8.2 Vacuous Proofs

The underlying principle of a vacuous proof is that a wrong premise allows you to conclude anything you want: Both $P: F, Q: F$ and $P: F, Q: T$ yield true for $P \rightarrow Q$.

For example, it was mentioned in remark 2.3 (Elements of the empty set and their properties) on $p .13$ that you can state anything you like about the elements of the empty set as there are none. The underlying principle of proving this kind of assertion is that of a vacuous proof. We prove here assertion d of that remark.
Theorem 3.10. Let $A$ be any set. Then $\emptyset \subseteq A$.
PROOF:
According to the definition of $\subseteq$ we must prove that if $x \in \emptyset$ then $x \in A$.
So let $x \in \emptyset$. We stop right here: " $x \in \emptyset$ " is a false statement regardless of the nature of $x$ because the empty set, by definition, does not contain any elements. It follows that $x \in A$.
Remark 3.23. You may ask: But is it not equally true that if $x \in \emptyset$ then $x \notin A$ ? The answer to that is YES, it is equally true that $x \in A$ ? and $x \notin A$ ?, but so what? First you'll find me an $x$ that belongs to the empty set and only then am I required to show you that it both does and does not belong to $A$ !

### 3.8.3 Direct Proofs

In a direct proof of $P \Rightarrow Q$ we assume the truth of the hypothesis $P$ and then employ logical equivalences, including the rules of inference, to show the truth of $Q$.

We proved in chapter 3.7.3 (An example of a direct proof) on p. 61 that each polynomial is differentiable (theorem 3.9). That was an example of a direct proof.

### 3.8.4 Proof by Contrapositive

A proof by contrapositive makes use of the logical equivalence $(P \Rightarrow Q) \Leftrightarrow(\neg Q \Rightarrow \neg P)$ (see the contrapositive law (3.42) on $p 46$ ). We give a direct proof of $\neg Q \Rightarrow \neg P$, i.e., we assume the falseness of $Q$ and prove that then $P$ must also be false. Here is an example.
Theorem 3.11. Let $A, B$ be two subsets of some universal set $\Omega$ such that $A \cap B^{\complement}=\emptyset$. Then $A \subseteq B$.
PROOF: We prove the contrapositive instead: If $A \nsubseteq B$ then $A \cap B^{\complement} \neq \emptyset$.
So let us assume $A \nsubseteq B$. This means that not every element of $A$ also belongs to $B$. In other words, there exists some $x \in A$ such that $x \notin B$. But then $x \in A \backslash B=A \cap B^{\complement}$, i.e., $A \cap B^{\complement} \neq \emptyset$.

We have proved from the negated conclusion $A \nsubseteq B$ the negated premise $A \cap B^{\complement} \neq \emptyset$.

### 3.8.5 Proof by Contradiction (Indirect Proof)

A proofs by contradiction are a generalization of proofs by contrapositive. We assume that it is possible for the implication $P \Rightarrow Q$ that the premise $P$ can be true and $Q$ can be false at the same time and construct the assumption of the truth of $P \cap \neg Q$ a statement $R$ such that both $R$ and $\neg R$ must be true. Here is an example.

Theorem 3.12. Let $A \subseteq \mathbb{Z}$ with the following properties:

$$
\begin{align*}
m, n & \in A \Rightarrow m+n \in A  \tag{3.95}\\
m, n & \in A \Rightarrow m n \in A  \tag{3.96}\\
& 0  \tag{3.97}\\
& \notin A  \tag{3.98}\\
\text { if } n & \in \mathbb{Z} \text { then either } n \in A \text { or }-n \in A \text { or } n=0 .
\end{align*}
$$

Then $1 \in A$.
Proof by contradiction: Assume that $A$ is a set of integers with properties (3.95) - (3.98) but that $1 \notin A$. We will show that then $1 \in A$ must be true. This finishes the proof because it is impossible that both $1 \notin A$ and $1 \in A$ are true.
a. It follows from $1 \notin A$ and (3.98) and $1 \neq 0$ that $-1 \in A$.
b. It now follows from (3.96) that $(-1) \cdot(-1) \in A$, i.e., $1 \in A$.

We have reached our contradiction.
Remark 3.24. In this simple proof the statement $R$ for which both $R$ and $\neg R$ were shown to be true happens to be the conclusion $1 \in A$. This generally does not need to be the case.

### 3.8.6 Proof by Cases

Sometimes an assumption $P$ is too messy to take on in its entirety and it is easier to break it down into two or more cases $P_{1}, P_{2}, \ldots, P_{n}$ each of which only covers part of $P$ but such that $P_{1} \vee P_{2} \vee \cdots \vee P_{n}$ covers all of it, i.e., we assume

$$
\begin{equation*}
P_{1} \vee P_{2} \vee \cdots \vee P_{n} \Leftrightarrow P \tag{3.99}
\end{equation*}
$$

Proof by cases then rests on the following theorem:
Theorem 3.13. Let $P, Q, P_{1} \vee P_{2} \vee \cdots \vee P_{n}$ be statements such that (3.99) is true. Then

$$
\begin{equation*}
(P \Rightarrow Q) \Leftrightarrow\left[\left(P_{1} \Rightarrow Q\right) \vee\left(P_{2} \Rightarrow Q\right) \vee \ldots\left(P_{n} \Rightarrow Q\right)\right] \tag{3.100}
\end{equation*}
$$

Proof (outline): You would do the proof by induction. Prove (3.100) first for $n=2$ by expressing $A \rightarrow B$ as $\neg A \vee B$ and then building a truth table that compares $\left(\neg\left(P_{1} \vee P_{2}\right)\right) \vee Q$ with $\neg P_{1} \vee Q \vee \neg P_{2} \vee Q$.

Then do the induction step in which (3.99) becomes $P_{1} \vee P_{2} \vee \cdots \vee P_{n+1} \Leftrightarrow P$ by setting $A:=$ $P_{1} \vee P_{2} \vee \cdots \vee P_{n}$ and this way reducing the proof of $(3.100)$ for $n+1$ to that of 2 components. You make the validity of $(A \Rightarrow Q) \Leftrightarrow\left[\left(P_{1} \Rightarrow Q\right) \vee\left(P_{2} \Rightarrow Q\right) \vee \ldots\left(P_{n} \Rightarrow Q\right)\right]$ the induction assumption.

Theorem 3.14. Prove that for any $x \in \mathbb{R}$ such that $x \neq 5$ we have

$$
\begin{equation*}
\frac{x}{x-5}>0 \Rightarrow[(x<0) \text { or }(x>5)] \tag{3.101}
\end{equation*}
$$

PROOF: There are two cases for which $x /(x-5)>0$ :
either both $x>0$ and $x-5>0$ or both $x<0$ and $x-5<0$. We write
$P:=" x /(x-5)>0$ ", ${ }^{36} P_{1}:=x>0$ and $x-5>0, P_{2}:=x<0$ and $x-5<0$. Then $P=P_{1} \vee P_{2}$.
case 1. $P_{1}$ :
Obviously $x>0$ and $x-5>0$ if and only if $x>5$, so we have proved $P_{1} \Rightarrow(x>5)$.
case 2. $P_{2}$ :
Obviously $x<0$ and $x-5<0$ if and only if $x<0$, so we have proved $P_{2} \Rightarrow(x<0)$.
We now conclude from $P=P_{1} \vee P_{2}$ and theorem 3.13 the validity of (3.101).

## 4 Functions and Relations (Study this!)

### 4.1 Cartesian Products and Relations

Definition 4.1 (Cartesian Product of two sets). The cartesian product of two sets $A$ and $B$ is

$$
A \times B:=\{(a, b): a \in A, b \in B\}
$$

i.e., it consists of all pairs $(a, b)$ with $a \in A$ and $b \in B$.

Two elements $\left(a_{1}, b_{1}\right)$ and $\left(a_{2}, b_{2}\right)$ are called equal if and only if $a_{1}=a_{2}$ and $b_{1}=b_{2}$. In this case we write $\left(a_{1}, b_{1}\right)=\left(a_{2}, b_{2}\right)$.

It follows from this definition of equality that the pairs $(a, b)$ and $(b, a)$ are different unless $a=b$. In other words, the order of $a$ and $b$ is important. We express this by saying that the cartesian product consists of ordered pairs.

As a shorthand, we abbreviate $A^{2}:=A \times A$.
Example 4.1 (Coordinates in the plane). Here is the most important example of a cartesian product of two sets. Let $A=B=\mathbb{R}$. Then $\mathbb{R} \times \mathbb{R}=\mathbb{R}^{2}=\{(x, y): x, y \in \mathbb{R}\}$ is the set of pairs of real numbers. I am sure you are familiar with what those are: They are just points in the plane, expressed by their $x$ - and $y$-coordinates.
Examples of such points are are: $(1,0) \in \mathbb{R}^{2}$ (a point on the $x$-axis), $(0,1) \in \mathbb{R}^{2}$ (a point on the $y$-axis), $(1.234,-\sqrt{2}) \in \mathbb{R}^{2}$.

You should understand why we do not allow two pairs to be equal if we flip the coordinates: Of course $(1,0)$ and $(0,1)$ are different points in the $x y$-plane!

Remark 4.1 (Function graphs as subsets of cartesian products). We gave the preliminary definition of a function in def.3.1, p. 24 of ch.3.1 (Prologue: Notation for Functions). ${ }^{37}$ A function

$$
f: X \rightarrow Y ; \quad y=f(x)
$$

which assigns each $x \in X$ to a unique function value $f(x) \in Y$, e.g., $f(x)=x^{2}$, is characterized by its graph

$$
\Gamma_{f}:=\{(x, f(x)): x \in X\}
$$

[^27]which is a subset of the cartesian product $X \times Y$. For example, if $X=[-2,3]$ and $Y=[0,10]$ then $\Gamma_{f}:=\left\{\left(x, x^{2}\right):-2 \leqq x \leqq 3\right\}$ is a subset of $[-2,3] \times[0,10]$. We will examine the connection between functions and their graphs in detail later in this chapter.

Remark 4.2 (Empty cartesian product). Note that $A \times B=\emptyset$ if and only if $A=\emptyset$ or $B=\emptyset$ or both are empty.
Definition 4.2 (Relation). Let $X$ and $Y$ be two sets and $R \subseteq X \times Y$ a subset of their cartesian product $X \times Y$. We call $R$ a relation on $(X, Y)$. A relation on $(X, X)$ is simply called a relation on $X$. If $(x, y) \in R$ we say that $x$ and $y$ are related and we usually write $x R y$ instead of $(x, y) \in R$.

A relation on $X$ is reflexive if $x R x$ for all $x \in X$. It is symmetric if $x_{1} R x_{2}$ implies $x_{2} R x_{1}$ for all $x_{1}, x_{2} \in X$. It is transitive if $x_{1} R x_{2}$ and $x_{2} R x_{3}$ implies $x_{1} R x_{3}$ for all $x_{1}, x_{2}, x_{3} \in X$. It is antisymmetric if $x_{1} R x_{2}$ and $x_{2} R x_{1}$ implies $x_{1}=x_{2}$ for all $x_{1}, x_{2} \in X$.

Here are some examples of relations.
Example 4.2 (Equality as a relation). Given a set $X$ let $R:=\{(x, x): x \in X\}$, i.e., $x R y$ if and only if $x=y$. This defines a relation on $X$ which is reflexive, symmetric and transitive.
Example 4.3 (Set inclusion as a relation). Given a set $X$ let $R:=\{(A, B): A, B \subseteq X$ and $A \subseteq B\}$, i.e., $A R B$ if and only if $A \subseteq B$. This defines a relation which is reflexive, antisymmetric and transitive.

Example 4.4 (Size of sets as a relation). Let $X$ be a finite set, i.e., a set which only contains finitely many elements. For $A \subseteq X$ let $|A|$ be the number of its elements. Let

$$
R:=\{(A, B): A, B \subseteq X \text { and }|A|=|B|\},
$$

i.e., $A R B$ if and only if $A$ and $B$ possess the same number of elements. This defines a relation on the power set $2^{X}$ of $X$ which is reflexive, symmetric and transitive.

Example 4.5 (Empty relation). Given two sets $X$ and $Y$ let $R:=\emptyset$. This empty relation is the only relation on $(X, Y)$ if $X$ or $Y$ is empty.

Example 4.6. Let $X:=\mathbb{R}^{2}$ be the $x y$-plane. For any point $\vec{x}=\left(x_{1}, x_{2}\right)$ in the plane let $\|\vec{x}\|:=$ $\sqrt{x_{1}^{2}+x_{2}^{2}}$ be its length ${ }^{38}$ and let $R:=\left\{(\vec{x}, \vec{y}) \in \mathbb{R}^{2} \times \mathbb{R}^{2}:\|\vec{x}\|=\|\vec{y}\|\right\}$. In other words, two points in the plane are related when they have the same length: they are located on a circle with radius $r=\|\vec{x}\|=\|\vec{y}\|$. The relation $R$ is reflexive, symmetric and transitive but not antisymmetric.

The relations given in examples 4.2, 4.4, 4.5 and 4.6 are reflexive, symmetric and transitive. Such relations are so important that they deserve a special name:

Definition 4.3 (Equivalence relations and equivalence classes). Let $R$ be a relation on a set $X$ which is reflexive, symmetric and transitive. We call such a relation an equivalence relation on $X$. It is customary to write $x \sim y$ rather than $x R y$ (or $(x, y) \in R$ ) and we say that $x$ and $y$ are equivalent

Given an equivalence relation " $\sim$ " on a set $X$ and $x \in X$ let

$$
\begin{equation*}
[x]_{\sim}:=\{y \in X: y \sim x\}=\{\text { all items equivalent to } x\} . \tag{4.1}
\end{equation*}
$$

We call $[x]_{\sim}$ the equivalence class of $x$. If it is clear from the context what equivalence relation is referred to then we simply write $[x]$ instead of $[x]_{\sim}$.

[^28]Relations which are reflexive, antisymmetric and transitive like the relation of example 4.3 (set inclusion) allow to compare items for "bigger" and "smaller" or "before" and "after". They also deserve a special name:
Definition 4.4 (Partial Order Relation). Let $R$ be a relation on a set $X$ which is reflexive, antisymmetric and transitive. We call such a relation a partial ordering of $X$. or a partial order relation on $X$. ${ }^{39}$ It is customary to write " $x \preceq y$ " or " $y \succeq x$ " rather than " $x R y$ " for a partial ordering $R$. We say that " $x$ before $y$ " or " $y$ after $x$ ".

If " $x \preceq y$ " defines a partial ordering on $X$ then $(X, \preceq)$ is called a partially ordered set set or a POset.
Remark 4.3. The properties of a partial ordering can now be phrased as follows:

$$
\begin{array}{lll}
x \preceq x \text { for all } x \in X \quad \text { reflexivity, } \\
x \preceq y \text { and } y \preceq x \Rightarrow y=x & \text { antisymmetry, } \\
x \preceq y \text { and } y \preceq z \Rightarrow & x \preceq z & \text { transitivity. } \tag{4.4}
\end{array}
$$

Remark 4.4 (Partial orderings and reflexivity). Note the following:
A. According to the above definition, the following are partial orderings of $X$ :

1. $X=\mathbb{R}$ and $x \preceq y$ if and only if $x \leqq y$.
2. $X=2^{\Omega}$ for some set $\Omega$ and $A \preceq B$ if and only if $A \subseteq B$ (example 4.3).
3. $X=\mathbb{R}$ and $x \preceq y$ if and only if $x \geqq y$.
B. The following relations are not partial orderings of $X$ because none of them is reflexive.
4. $X=\mathbb{R}$ and $x \preceq y$ if and only if $x<y$.
5. $X=2^{\Omega}$ for some set $\Omega$ and $A \preceq B$ if and only if $A \subset B$ (i.e., $A \subseteq B$ but $A \neq B$ ).
6. $X=\mathbb{R}$ and $x \preceq y$ if and only if $x>y$.

Note that each one of those three relations is antisymmetric. For example, let us look at $x<y$. It is indeed true that the premise $[x<y$ and $y<x]$ allows us to conclude that $y=x$ as there are no such numbers $x$ and $y$ and a premise that is known never to be true allows us to conclude anything we want!
C. An equivalence relation $\sim$ is a never a partial ordering of $X$ except in the very uninteresting case where you have $x \sim y$ if and only if $x=y$.
D. A partial ordering of $X$, as any relation on $X$ in general, is inherited by any subset $A \subseteq X$ as follows: Let $\preceq$ be a partial ordering on a set $X$ and let $A \subseteq X$. We define a relation $\preceq_{A}$ on A as follows: Let $x, y \in A$. Then $x \preceq_{A} y$ if and only if $x \preceq y$.

Definition 4.5 (Inverse Relation). Let $X$ and $Y$ be two sets and $R \subseteq X \times Y$ a relation on $(X, Y)$. Let

$$
R^{-1}:=\{(y, x):(x, y) \in R\} .
$$

Clearly $R^{-1}$ is a subset of $Y \times X$ and hence a relation on $(Y, X)$. We call $R^{-1}$ the inverse relation to the relation $R$.
Example 4.7. Let $R:=\left\{\left(x, x^{3}\right): x \in \mathbb{R}\right\}$. Then $R$ is the relation on $\mathbb{R}$ which represents the function $y=f(x)=x^{3}$. We obtain

$$
R^{-1}=\left\{\left(x^{3}, x\right): x \in \mathbb{R}\right\}=\left\{\left(y, y^{1 / 3}\right): y \in \mathbb{R}\right\} .
$$

In other words, $R^{-1}$ represents the inverse function $x=f^{-1}(y)=y^{1 / 3}$.

[^29]
### 4.2 Functions (Mappings) and Families

### 4.2.1 Some Preliminary Observations about Functions

Remark 4.5 (A layman's definition of a function). We look at the set $\mathbb{R}$ of all real numbers ${ }^{40}$ and the function $y=f(x)=\sqrt{4-x^{2}}$ which associates with certain real numbers $x$ (the "argument" or "independent variable") another real number $y=\sqrt{4-x^{2}}$ (the "function value" or "dependent variable"):
$f(0)=\sqrt{(4-0}=2, f(2)=f(-2)=\sqrt{(4-4}=0, f(2 / 3)=f(-2 / 3)=\sqrt{(36-4) / 9}=\sqrt{30} / 3, \ldots$
You can think of this function as a rule or law which specifies what item $y$ is obtained as the output or result if the item $x$ is provided as input.

Let us take a closer look at the function $y=f(x)=\sqrt{4-x^{2}}$ and its properties:
a. For some real numbers $x$ there is no function value: For example, if $x=10$ then $4-x^{2}=$ -96 is negative and the square root cannot be taken.
b. For some other $x$, e.g., $x=0$ or $x=2 / 3$, there is a function value $f(x)$. A moment's reflection shows that the biggest possible set of potential arguments for our function, called by some authors the natural domain of the function (e.g., [2] Brewster/Geoghegan), is the interval $[-2,2]$. It is customary to write $D_{f}$ for the natural domain of a function $y=f(x)$.
c. For a given $x$ there is never more than one function value $f(x)$. This property allows us to think of a function as an assignment rule: It assigns to certain arguments $x$ a unique function value $f(x)$. We observed in $\mathbf{b}$ that $f(x)$ exists if and only if $x \in[-2,2]$.
d. Not every $y \in \mathbb{R}$ is suitable as a function value: A square root cannot be negative, hence no $x$ exists such that $f(x)=-1$ or $f(x)=-\pi$.
e. On the other hand, there are numbers $y$ such as $y=0$, which are "hit" more than once by the function: $f(2)=f(-2)=0 .{ }^{41}$
f. Graphs as drawings: We are used to look at the graphs of functions. Here is a picture of the graph of $f(x)=\sqrt{4-x^{2}}$.


[^30]g. Graphs as sets: Drawings as the one above have limited precision (the software should have drawn a perfect half circle with radius 2 about the origin but there seem to be wedges at $x \approx \pm 1.8$ ). Also, how would you draw a picture of a function which assigns a 3-dimensional vector ${ }^{42}(x, y, z)$ to its distance $w=F(x, y, z)=\sqrt{x^{2}+y^{2}+z^{2}}$ from the zero vector $(0,0,0)$ ? You would need four dimensions, one each for $x, y, z, w$, to draw the graph!
To express the graph of a function without a picture, let us look at a verbal description: The graph of a function $f(x)$ is the collection of the pairs $(x, f(x))$ for all points $x$ which belong to the set $[-2,2]$ of potential arguments (see a). In mathematical parlance: The graph of the function $f(x)$ is the set
$$
\Gamma_{f}:=\left\{(x, f(x)): x \in D_{f}\right\}
$$
(see remark 4.1 on p. 67).

We now make adjustments to some of those properties which will get us closer to the definition of a function as it is used in abstract mathematics.

Remark 4.6 (A better definition of a function). We make the following alterations to remark 4.5.

- We require an upfront specification of the set $A$ of items that will be allowed as input (arguments) for the function and we require that $y=f(x)$ makes sense for each $x \in A$. Given the function $y=f(x)=\sqrt{4-x^{2}}$ from above this means that $A$ must be a subset of $[-2,2]$.
- We require an upfront specification of the set $B$ of items that will be allowed as output (function values) for the function. This set must be so big that each $x \in A$ has a function value $y \in B$. We do not mind if $B$ contains redundant $y$ values. For $y=f(x)=\sqrt{4-x^{2}}$ any superset of the closed interval $[0,2]$ will do. We may choose, e.g., $B:=[0,2]$ or $B:=[-2,2 \pi]$ or $B:=[0,4]$ or $B:=\mathbb{R} \cup\{$ all inhabitants of Chicago $\}$.
Doing so gives us the following: A function consists of three items: a set $A$ of inputs, a set $B$ of outputs and an assignment rule $x \mapsto f(x)$ with the following properties:

1. For all inputs $x \in A$ there is a function value $f(x) \in B$.
2. For any input $x \in A$ there is never more than one function value $f(x) \in B$. It follows from property 1 that each $x \in A$ uniquely determines its function value $y=f(x)$. This property is what allows us to think of a function as an assignment rule: It assigns to each $x \in A$ a unique function value $f(x) \in B$.
3. Not every $y \in B$ needs to be a function value $f(x)$ for some $x \in A$, i.e., the set $\{x \in A$ : $f(x)=y\}$ can be empty.
4. On the other hand there may be numbers $y$ which are "hit" more than once by $f$.

Example: Let $A:=\mathbb{N}, B:=\mathbb{R}, f(x):=(-1)^{x}$. Then both -1 and 1 are mapped to infinitely often by $f$.

[^31]5. The graph $\Gamma_{f}$ of a function $f(x)$ is the collection of the pairs $(x, f(x))$ for all points $x$ which belong to the set $A$, i.e.,
\[

$$
\begin{equation*}
\Gamma_{f}:=\{(x, f(x)): x \in A\} . \tag{4.5}
\end{equation*}
$$

\]

$\Gamma_{f}$ has the following properties:
5a. $\quad \Gamma_{f} \subseteq A \times B$, i.e., $\Gamma_{f}$ is a relation on $(A, B)$ (see def.4.2 on p.68).
5b. For each $x \in A$ there exists a unique $y \in B$ such that $(x, y) \in \Gamma_{f}$
5c. If $x \mapsto g(x)$ is another function with inputs $A$ and outputs $B$ which is different from $x \mapsto f(x)$ (i.e., there is at least one $a \in A$ such that $f(a) \neq g(a)$ ) then the graphs $\Gamma_{f}$ and $\Gamma_{g}$ do not coincide
6. Conversely, if $A$ and $B$ are two nonempty sets, then any relation $\Gamma$ on $(A, B)$ which satifies $\mathbf{5 a}$ and $\mathbf{5 b}$ uniquely determines a function $x \mapsto f(x)$ with inputs $A$ and outputs $B$ as follows: For $a \in A$ we define $f(a)$ to be the element $b \in B$ for which $(a, b) \in \Gamma$. We know from $\mathbf{5 b}$ that such $b$ exists and is uniquely determined.

Here is a complicated way of looking at the example above: Let $X=[-2,2]$ and $Y=\mathbb{R}$. Then $y=$ $f(x)=\sqrt{4-x^{2}}$ is a rule which "maps" each element $x \in X$ to a uniquely determined number $y \in Y$ which depends on $x$ as follows: Subtract the square of $x$ from 4, then take the square root of that difference.

Mathematicians are very lazy as far as writing is concerned and they figured out long ago that writing "depends on xyz" all the time not only takes too long, but also is aesthetically very unpleasing and makes statements and their proofs hard to understand. They decided to write " $(x y z)$ " instead of "depends on $x y z$ " and the modern notion of a function or mapping $y=f(x)$ was born.

Here is another example: if you say $f(x)=x^{2}-\sqrt{2}$, it's just a short for "I have a rule which maps a number $x$ to a value $f(x)$ which depends on $x$ in the following way: compute $x^{2}-\sqrt{2}$." It is crucial to understand from which set $X$ you are allowed to pick the "arguments" $x$ and it is often helpful to state what kinds of objects $f(x)$ the $x$-arguments are associated with, i.e., what set $Y$ they will belong to.

We now are ready to give the precise definition of a function.

### 4.2.2 Definition of a Function and Some Basic Properties

We have seen in remark 4.6 on $p .71$ that a function can be thought of equivalently as an assignment rule $x \mapsto f(x)$ or as a graph. Mathematicians prefer the latter because "assignment rule" is a rather vague term (an undefined term in the sense of ch. 3.7.1 (Building blocks of mathematical theories) on p.56) whereas " $g$ raph" is entirely defined in the language of sets. This chapter starts with the official definition of a function. It then deals with the following concepts: composition of functions, injective, surjective, bijective and inverse functions, restriction and extension of functions.

Definition 4.6 (Mappings (functions)). Given are two arbitrary nonempty sets $X$ and $Y$ and a relation $\Gamma$ on $(X, Y)$ (see 4.2 on p.68) which satisfies the following:

$$
\begin{equation*}
\text { for each } x \in X \text { there exists exactly one } y \in Y \text { such that }(x, y) \in \Gamma \text {. } \tag{4.6}
\end{equation*}
$$

We call the triplet $f(\cdot):=(X, Y, \Gamma)$ a function or mapping from $X$ to $Y$. The set $X$ is called the domain or source and $Y$ is called the codomain or target of the mapping $f(\cdot)$. We will mostly use the words "domain" and "codomain" in this document.

Usually mathematicians simply write $f$ instead of $f(\cdot)$ We mostly follow that convention, but sometimes include the " $(\cdot)$ " part to emphasize that a function rather than an "ordinary"' element of a set is involved.

Let $x \in X$. We write $f(x)$ for the uniquely determined $y \in Y$ such that $(x, y) \in \Gamma$. We write $\Gamma_{f}$ or $\Gamma(f)$ if we want to stress that $\Gamma$ is the relation associated with the function $f=(X, Y, \Gamma)$. We call $\Gamma$ the graph of the function $f$. Clearly

$$
\begin{equation*}
\Gamma=\Gamma_{f}=\Gamma(f)=\{(x, f(x)): x \in X\} \tag{4.7}
\end{equation*}
$$

It is customary to write

$$
\begin{equation*}
f: X \rightarrow Y, \quad x \mapsto f(x) \tag{4.8}
\end{equation*}
$$

instead of $f=(X, Y, \Gamma)$ and we henceforth follow that convention. We abbreviate that to $f: X \rightarrow Y$ if it is clear or irrelevant how to compute $f(x)$ from $x$. We read " $a \mapsto b$ " as " $a$ is assigned to $b$ " or " $a$ maps to $b^{\prime \prime}$ and refer to $\mapsto$ as the maps to operator or assignment operator.

Domain elements $x \in X$ are called independent variables or arguments and $f(x) \in Y$ is called the function value of $x$. The subset

$$
\begin{equation*}
f(X):=\{y \in Y: y=f(x) \text { for some } x \in X\}=\{f(x): x \in X\} \tag{4.9}
\end{equation*}
$$

of $Y$ is called the range or image of the function $f(\cdot) .{ }^{43}$
We say " $f$ maps $X$ into $Y$ " and " $f$ maps the domain value $x$ to the function value $f(x)$ ".
Figure 4.1 on $p .74$ illustrates the graph of a function as a subset of $X \times Y$.
Remark 4.7 (Mappings vs. functions). Mathematicians do not always agree $100 \%$ on their definitions. The issue of what is called a function and what is called a mapping is subject to debate. Some mathematicians call a mapping a function only if its codomain is a subset of the real numbers ${ }^{44}$ but the majority does what I'll try to adhere to in this document: I use "mapping" and "function" interchangeably and I'll talk about real valued functions rather than just functions if the codomain is part of $\mathbb{R}$ (see (4.12) on p .84 ).

Remark 4.8. The symbol $x$ chosen for the argument of the function is a dummy variable in the sense that it does not matter what symbol you use.

The following each define the same function with domain $[0, \infty[$ and codomain $\mathbb{R}$ which assigns to any non-negative real number its (positive) square root:

$$
\begin{array}{ll}
f:[0, \infty[\rightarrow \mathbb{R}, & x \mapsto \sqrt{x}, \\
f:[0, \infty[\rightarrow \mathbb{R}, & y \mapsto \sqrt{y}, \\
f:[0, \infty[\rightarrow \mathbb{R}, & f(\gamma)=\sqrt{\gamma}
\end{array}
$$

Matter of fact, not even the symbol you choose for the function matters as long as the operation (here: assign a number to its square root) is unchanged. In other words, the following still describe

[^32]

Figure 4.1: Graph of a function.
the same function as above:

$$
\begin{array}{ll}
\varphi:[0, \infty[\rightarrow \mathbb{R}, & t \mapsto \sqrt{t}, \\
A:[0, \infty[\rightarrow \mathbb{R}, & x \mapsto \sqrt{x}, \\
g:[0, \infty[\rightarrow \mathbb{R}, & g(A)=\sqrt{A}
\end{array}
$$

In contrast, the following three functions all are different from each other and none of them equals $f$ because domain and /or codomain do not match:

$$
\begin{aligned}
& \psi:] 0, \infty[\rightarrow \mathbb{R}, \quad x \mapsto \sqrt{x} \quad \text { (different domain). } \\
& B:[0, \infty[\rightarrow] 0, \infty[, \quad x \mapsto \sqrt{x} \quad \text { (different codomain), } \\
& h:[0,1[\rightarrow[0,1[, \quad x \mapsto \sqrt{x} \quad \text { (different domain and codomain). }
\end{aligned}
$$

Definition 4.7 (Function composition). Given are three nonempty sets $X, Y$ and $Z$ and two functions $f: X \rightarrow Y$ and $g: Y \rightarrow Z$. Given $x \in X$ we know the meaning of the expression $g(f(x))$ :
$y:=f(x)$ is the function value of $x$ for the function $f$, i.e., the unique $y \in Y$ such that $(x, y) \in \Gamma_{f}$.
$z:=g(y)=g(f(x))$ is the function value of $f(x)$ for the function $g$, i.e., the unique $z \in Z$ such that $(f(x), z)=(f(x), g(f(x))) \in \Gamma_{g}$.

The set $\Gamma:=\{(x, g(f(x)): x \in X)\}$ is a relation on $(X, Z)$ such that
(4.10) for each $x \in X$ there exists exactly one $z \in Z$, namely, $z=g(f(x))$, such that $(x, z) \in \Gamma$.

It follows that $\Gamma$ is the graph of a function $h=(X, Z, \Gamma)$ with function values $h(x)=g(f(x))$ for each $x \in X$. We call $h$ the composition of $f$ and $g$ and we write $h=g \circ f$ (" $g$ after $f$ ").

As far as notation is concerned it is OK to write either of $g \circ f(x)$ or $(g \circ f)(x)$. The additional parentheses may give a clearer presentation if $f$ and/or $g$ are defined by fairly complex formulas.

The following shows how you diagram the composition of two functions. The left picture shows the domains and codomains for each mapping and the right one shows the element assignments.

## Function composition



We have a special name for the "do nothing function" which assigns each argument to itself:
Definition 4.8 (identity mapping). Given any nonempty set $X$, we use the symbol $i d_{X}$ for the identity mapping defined as

$$
i d_{X}: X \rightarrow X, \quad x \mapsto x .
$$

We drop the subscript if it is clear what set is referred to.

### 4.2.3 Examples of Functions

We now give some examples of functions. You might find some of them rather difficult to understand at first reading.
Example 4.8. Let $\Gamma:=\left\{\left(x, x^{3}\right): x \in \mathbb{R}\right\} \subseteq \mathbb{R} \times \mathbb{R}$. Then $f=(\mathbb{R}, \mathbb{R}, \Gamma)$ is the function

$$
f: \mathbb{R} \rightarrow \mathbb{R}, \quad x \mapsto x^{3} .
$$

Example 4.9. Let $\Gamma:=\left\{\left(x, x^{2}+1\right): x \in \mathbb{R}\right\}$. Then $g=(\mathbb{R}, \mathbb{R}, \Gamma)$ is the function

$$
g: \mathbb{R} \rightarrow \mathbb{R}, \quad x \mapsto x^{2}+1
$$

Example 4.10. Let $\Gamma:=\{(a, \ln (a)): a \in] 0, \infty[ \}$. Here $\ln (a)$ denotes the natural logarithm of $a$. Then $h=(] 0, \infty[, \mathbb{R}, \Gamma)$ is the function

$$
h:] 0, \infty[\rightarrow \mathbb{R}, \quad x \mapsto \ln (x) .
$$

Example 4.11. Let $\Gamma:=\{(x, \sqrt{x}): x \in[0, \infty[ \}$. Then $\varphi=([0, \infty[, \mathbb{R}, \Gamma)$ is the function

$$
\varphi:[0, \infty[\rightarrow \mathbb{R}, \quad x \mapsto \sqrt{x}
$$

Example 4.12. Let $\Gamma:=\{(x, \sqrt{x}): x \in[0, \infty[ \}$. We can consider $\Gamma$ as a subset of $[0, \infty[\times \mathbb{R}$ but also as a subset of $[0, \infty[\times[0, \infty[$. In the first case we obtain a function $\varphi=([0, \infty[, \mathbb{R}, \Gamma)$, i.e., the function

$$
\varphi:[0, \infty[\rightarrow \mathbb{R}, \quad x \mapsto \sqrt{x}
$$

In the second case we obtain a different(!) function $\psi=([0, \infty[,[0, \infty[, \Gamma)$, i.e., the function

$$
\psi:[0, \infty[\rightarrow[0, \infty[, \quad x \mapsto \sqrt{x} .
$$

If you have taken multivariable calculus or linear algebra then you know that functions need not necessarily map numbers to numbers but they can also map vectors to numbers, numbers to vectors (curves) or vectors to vectors.

Example 4.13. We define a function which maps two-dimensional vectors to numbers. Let $A:=$ $\left\{\left((x, y) \in \mathbb{R}^{2}: x^{2}+y^{2} \leqq 1\right\}\right.$. Let $\Gamma:=\left\{\left((x, y), \sqrt{1-x^{2}-y^{2}}\right):(x, y) \in A\right\}$. Then $F=(A, \mathbb{R}, \Gamma)$ is the function

$$
F: A \rightarrow \mathbb{R}, \quad(x, y) \mapsto \sqrt{1-x^{2}-y^{2}} .
$$

Note that the domain is not a set of real numbers but of points in the plane and that the graph of $F$ is a set of points $(x, y, z)$ in 3 -dimensional space.
Example 4.14. We define a function which maps numbers to two-dimensional vectors (a curve in the plane). Let $\Gamma:=\{(t,(\sin t, \cos t)): t \in \mathbb{R}\}$. Then $G=\left(\mathbb{R}, \mathbb{R}^{2}, \Gamma\right)$ is the function

$$
G: \mathbb{R} \rightarrow \mathbb{R}^{2}, \quad t \mapsto(\sin t, \cos t)
$$

whose image $G\left(\mathbb{R}\right.$ is the unit circle $\left\{(x, y) \in \mathbb{R}^{2}: x^{2}+y^{2}=1\right\}$ Note that the codomain is not a set of real numbers but the Euclidean plane.
Example 4.15. Let $\Gamma:=\{((x, y),(2 x-y / 3, x / 6+4 y)): x, y \in \mathbb{R}\}$. Then $H=\left(\mathbb{R}^{2}, \mathbb{R}^{2}, \Gamma\right)$ is the function

$$
H: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}, \quad(x, y) \mapsto(2 x-y / 3, x / 6+4 y)
$$

Note that both domain and codomain are the Euclidean plane.
We now reformulate the last example in the framework of linear algebra. Skip this next example if you do not know about matrix multiplication.
Example 4.16. As is customary in linear algebra we now think of $\mathbb{R}^{2}$ as the collection of column vectors $\left\{\binom{x}{y}: x, y \in \mathbb{R}\right\}$ rather than the cartesian product $\mathbb{R} \times \mathbb{R}$ which is the collection of row vectors $\{(x, y): x, y \in \mathbb{R}\}$.

Let $A$ be the $2 \times 2$ matrix

$$
A:=\left(\begin{array}{cc}
2 & -1 / 3 \\
1 / 6 & 4
\end{array}\right)
$$

We then obtain for any pair of numbers $\vec{x}=(x, y)^{T 45}$ that

$$
\begin{aligned}
& A \vec{x}=\left(\begin{array}{cc}
2 & -1 / 3 \\
1 / 6 & 4
\end{array}\right)\binom{x}{y}=\binom{2 x-y / 3}{x / 6+4 y} \\
& \text { Let } \Gamma:=\left\{\left(\binom{x}{y},\binom{2 x-y / 3}{x / 6+4 y}\right): x, y \in \mathbb{R}\right\} \text {. Then } H=\left(\mathbb{R}^{2}, \mathbb{R}^{2}, \Gamma\right) \text { is the function } \\
& H: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}, \quad\binom{x}{y} \mapsto A\binom{x}{y} .
\end{aligned}
$$

Note that both domain and codomain are the Euclidean plane.
If you want to construct a counterexample to a mathematical statement concerning functions it often is best to construct functions with small domain and codomain so that you can draw a picture that completely describes the assignments. The next example will illustrate this.
${ }^{45}$ Here $(x, y)^{T}=\binom{x}{y}$ is the transpose of $(x . y)$, i.e., the operation that switches rows and columns of any matrix. In particular it transforms a row vector into a column vector and vice versa.

## Example 4.17.

Let $X:=\{a, b, c, d\}, Y:=\{x, y, z\}, \Gamma:=\{(a, y),(b, y),(c, z),(d, y)\}$. Then $I=(X, Y, \Gamma)$ is the function which maps the elements of $X$ to $Y$ according to the diagram on the right. Note that nothing was said about the nature of the elements of $X$ and $Y$. One need not know about it to make observations like the following: Examine items 3 and 4 of remark 4.6 (A better definition of a function) on p.71. Convince yourself that $x \in Y$ is an example for 3 : Not every element

of $Y$ needs to be a function value and that $y \in Y$ is an example for 4: There may be elements of $Y$ which are "hit" more than once by the function.

Example 4.18. This example represents a mathematical model for computing probabilities of the outcomes of rolling a fair die and demonstrates that probability can be thought of as a function that maps sets to numbers.

If we roll a die then the outcome will be an integer between 1 and 6 , i.e., the "state space" for this random action will be $X:=\{1,2,3,4,5,6\}$. For $A \subseteq X$ let $\operatorname{Prob}(A)$ denote the probability that rolling the die results in an outcome $x \in A$.

For example $\operatorname{Prob}$ ( an even number occurs $)=\operatorname{Prob}(\{2,4,6\})=50 \%=1 / 2$. Clearly we have for singletons consisting of a single outcome that

$$
\operatorname{Prob}(\{1\})=\operatorname{Prob}(\{2\})=\cdots=\operatorname{Prob}(\{6\})=1 / 6=16 . \overline{6} \%
$$

Your everyday experience tells you that if $A=\left\{x_{1}, x_{2}, \ldots, x_{k}\right\}$ where $x_{j} \in X$ for each index $j$ (and hence $k \leqq 6$ because a set does not contain duplicates) then

$$
\operatorname{Prob}(A)=\operatorname{Prob}\left(\left\{x_{1}\right\}\right)+\operatorname{Prob}\left(\left\{x_{2}\right\}\right)+\cdots+\operatorname{Prob}\left(\left\{x_{k}\right\}\right)=\sum_{j=1}^{k} \operatorname{Prob}\left(\left\{x_{j}\right\}\right)
$$

What if $A$ is the event that the roll of the die does not result in any outcome, i.e., $A=\emptyset$ ? We do not worry about the die getting stuck in mid-air or the dog snatching it before we get a chance to see the outcome and consider this event impossible, i.e., $\operatorname{Prob}(\emptyset)=0$.

We now have a probability associated with every $A \subseteq X$, i.e., with every $A \in 2^{X}$ and can finally write this probability as a function. Let $\Gamma:=\{(A, \operatorname{Prob}(A)): A \subseteq X\}$. Then $P=\left(2^{X},[0,1], \Gamma\right)$ is the function

$$
P: 2^{X} \rightarrow[0,1], \quad A \mapsto \operatorname{Prob}(A)
$$

Why do we use $[0,1]$ and not $\mathbb{R}$ as the codomain? The answer is that we could have done so but no event has a probablity that exceeds $100 \%$ or is negative, so $[0,1]$ is big enough and by choosing this set as the codomain we do not deviate from standard presentation of mathematical probability theory.

To understand the next example you need to be familiar with the concepts of continuity, differentiability and antiderivatives (integrals) of functions of a single variable. Just skip the parts where you lack the background.

Example 4.19. Let $a \in \mathbb{R} \cup\{-\infty\}$ and $b \in \mathbb{R} \cup\{\infty\}$ and let $X:=] a, b$ [ be the open (end points $a, b$ are excluded) interval of all real numbers between $a$ and $b$. Let $\left.x_{0} \in\right] a, b[$ be "fixed but arbitrary".

The following is known from calculus (see [11] Stewart, J: Single Variable Calculus): Let $f$ : $] a, b[\rightarrow \mathbb{R}$ be a function which is continuous on $] a, b[$. Then
a. $\quad f$ is integrable, i.e., for any $\alpha, \beta \in \mathbb{R}$ such that $a<\alpha<\beta<b$ the definite integral $\int_{\alpha}^{\beta} f(u) d u$ exists. For a definition of integrability see section 4.1 of [11] Stewart, J: Single
Variable Calculus.
b. Integration is "linear", i.e., it is additive: $\int_{\alpha}^{\beta}(f(u)+g(u)) d u=\int_{\alpha}^{\beta} f(u) d u+\int_{\alpha}^{\beta} g(u) d u$, and you also can "pull out" constant $\lambda \in \mathbb{R}: \int_{\alpha}^{\beta} \lambda f(u) d u=\lambda \int_{\alpha}^{\beta} f(u) d u$.
c. Integration is "monotone":

If $f(x) \leqq g(x)$ for all $\alpha \leqq x \leqq \beta$ then $\int_{\alpha}^{\beta}(f(u)) d u \leqq \int_{\alpha}^{\beta} g(u) d u$.
d. $\quad f$ has an antiderivative : There exists a function $F:] a, b\left[\rightarrow \mathbb{R}\right.$ whose derivative $F^{\prime}(\cdot)$ exists on all of $] a, b\left[\right.$ and coincides with $f$, i.e., $F^{\prime}(x)=f(x)$ for all $\left.x \in\right] a, b[$.
e. This antiderivative satisfies $\int_{\alpha}^{\beta} f(u) d u=F(\beta)-F(\alpha)$ for all $a<\alpha<\beta<b$ and it is not uniquely defined: If $C \in \mathbb{R}$ then $F(\cdot)+C$ is also an antiderivative of $f$.
On the other hand, if both $F_{1}$ and $F_{2}$ are antiderivatives for $f$ then their difference $G(\cdot):=$ $F_{2}(\cdot)-F_{1}(\cdot)$ has the derivative $G^{\prime}(\cdot)=f(\cdot)-f(\cdot)$ which is constant zero on $] a, b[$. It follows that any two antiderivatives only differ by a constant.
To summarize the above: If we have one antiderivative $F$ of $f$ then any other antiderivative $\tilde{F}$ is of the form $\tilde{F}(\cdot)=F(\cdot)+C$ for some real number $C$.
This fact is commonly expressed by writing $\int f(u) d u=F(x)+C$ for the indefinite integral (an integral without integration bounds).
f. It follows from $\mathbf{e}$ that if some $c_{0} \in \mathbb{R}$ is given then there is only one antiderivative $F$ such that $F\left(x_{0}\right)=c_{0}$.
Here is a quick proof: Let $G$ be another antiderivative of $f$ such that $G\left(x_{0}\right)=c_{0}$. Because $G-F$ is constant we have for all $x \in] a, b[$ that

$$
G(x)-F(x)=\text { const }=G\left(x_{0}\right)-F\left(x_{0}\right)=0
$$

i.e., $G=F$.

After those reminders about integration we are ready to define a function $I(\cdot)$ for which both domain and codomain are sets of functions.

Given are $a, b$ and $x_{0}$ as above and $c_{0} \in \mathbb{R}$. Let

$$
\begin{aligned}
\mathscr{F} & :=\{f:] a, b[\rightarrow \mathbb{R} \text { such that } f \text { is continuous on }] a, b[ \}, \\
\mathscr{G} & :=\{g:] a, b[\rightarrow \mathbb{R} \text { such that } g \text { is differentiable on }] a, b\left[\text { and } g\left(x_{0}\right)=c_{0}\right\} .
\end{aligned}
$$

It follows from the preparatory remarks that for each $f \in \mathscr{F}$ there exists a unique $F \in \mathscr{G}$ which is an antiderivative for $f$. We now define a function $I: \mathscr{F} \rightarrow \mathscr{G}$ by specifying its graph as the set

$$
\Gamma:=\left\{(f(\cdot), g(\cdot)): f \in \mathscr{F}, g \in \mathscr{G}, g \text { is an antiderivative of } f \text { and } g\left(x_{0}=c_{0}\right\} .\right.
$$

It follows from $\mathbf{e}$ that the function value $I(f)(x)$ for any $a<x<b$ is

$$
I(f)(x)=\int_{x_{0}}^{x} f(u) d u-I(f)\left(x_{0}\right)=\int_{x_{0}}^{x} f(u) d u-c_{0}
$$

### 4.2.4 Injective, Surjective and Bijective functions

Definition 4.9 (Surjective, injective, bijective). Let $f: X \rightarrow Y$. As usual the graph of $f$ is denoted $\Gamma_{f}$.
a. Surjectivity: In general it is not true that $f(X)=\{f(x): x \in X\}$ equals the entire codomain $Y$, i.e., that
for each $y \in Y$ there exists at least one $x \in X$ such that $(x, y) \in \Gamma_{f}$.
But if $f(X)=Y$, i.e., if (4.12) holds, we call $f$ surjective and we say that $f$ maps $X$ onto $Y$.
b. Injectivity: In general it is not true that if $y \in f(X)$ then $y=f(x)$ for a unique $x$, i.e., that if there is another $x_{1} \in X$ such that also $y=f\left(x_{1}\right)$ then it follows that $x_{1}=x$. But if this is the case, i.e., if
for each $y \in Y$ there exists at most one $x \in X$ such that $(x, y) \in \Gamma_{f}$.
then we call $f$ injective .
We can express (4.13) also as follows: If $x, x_{1} \in X$ and $y \in Y$ are such that $(x, y) \in \Gamma_{f}$ and $\left(x_{1}, y\right) \in \Gamma_{f}$ then it follows that $x_{1}=x$.
c. Bijectivity: Let $f: X \rightarrow Y$ be both injective and surjective. Such a function is called bijective. We often write $f: X \xrightarrow{\sim} Y$ for a bijective function $f$.

It follows from (4.12) and (4.13) that $f$ is bijective if and only if

$$
\begin{equation*}
\text { for each } y \in Y \text { there exists exactly one } x \in X \text { such that }(x, y) \in \Gamma_{f} \text {. } \tag{4.14}
\end{equation*}
$$

We rewrite (4.14) by employing $\Gamma_{f}$ 's inverse relation $\Gamma_{f}^{-1}=\{(y, x):(x . y) \in \Gamma\}$ (see def. 4.5 on p.69) and obtain

$$
\begin{equation*}
\text { for each } y \in Y \text { there exists exactly one } x \in X \text { such that }(y, x) \in \Gamma_{f}^{-1} \text {. } \tag{4.15}
\end{equation*}
$$

But this implies, according to (4.6), that $\Gamma_{f}^{-1}$ is the graph of a function $g:=\left(Y, X, \Gamma_{f}^{-1}\right)$ with domain $Y$ and codomain $X$ where, for a given $y \in Y, g(y)$ stands for the uniquely determined $x \in X$ such that $(y, x) \in \Gamma_{f}^{-1}$. Note that

$$
\begin{equation*}
\Gamma_{f}^{-1}=\Gamma_{g} . \tag{4.16}
\end{equation*}
$$

We call $g$ the inverse mapping or inverse function of $f$ and write $f^{-1}$ instead of $g$.

## Remark 4.9.

a. It follows from (4.16) that

$$
\begin{equation*}
\Gamma_{f}^{-1}=\Gamma_{f^{-1}} . \tag{4.17}
\end{equation*}
$$

b. Each $x \in X$ is mapped to $y=f(x)$ which is the only element of $Y$ such that $f^{-1}(y)=x$, c. Each $y \in Y$ is mapped to $x=f^{-1}(y)$ which is the only element of $X$ such that $f(x)=y$.
d. It follows from $\mathbf{b}$ and $\mathbf{c}$ that

$$
\begin{equation*}
\text { if } x \in X, y \in Y \text { then } f(x)=y \Leftrightarrow x=f^{-1}(y) \tag{4.18}
\end{equation*}
$$

e. It also follows from $\mathbf{b}$ and $\mathbf{c}$ that $f^{-1}(f(x))=x$ for all $x \in X$ and $f\left(f^{-1}(y)\right)=y$ for all $y \in Y$. In other words, $f^{-1} \circ f=i d_{X}$ and $f \circ f^{-1}=i d_{Y}$. Here is the picture:

Inverse function


Theorem 4.1 (Characterization of inverse functions). Let $X$ and $Y$ be nonempty sets and $f: X \rightarrow Y$. Then the following are equivalent:
a. fis bijective.
b. There exists $g: Y \rightarrow X$ such that both $g \circ f=i d_{X}$ and $f \circ g=i d_{Y}$.

PROOF of $\boldsymbol{a} \Rightarrow \boldsymbol{b}$ : We have seen in part $\boldsymbol{e}$ of remark 4.9 that $g:=f^{-1}$ satisfies $\boldsymbol{b}$.
PROOF of $\boldsymbol{b} \Rightarrow \boldsymbol{a}$ : We must show that $f$ is both surjective and injective. First we show that $f$ is surjective. Let $y \in Y$. we must find some $x \in X$ such that $f(x)=y$. Let $x:=g(y)$. Then

$$
f(x)=f(g(y))=f \circ g(y)=i d_{Y}(y)=y
$$

We have $f(x)=y$ and this proves surjectivity. Now we show that $f$ is injective. Let $x_{1}, x_{2} \in X$ and $y \in Y$ such that $f\left(x_{1}\right)=f\left(x_{2}\right)=y$. We are done if we can prove that $x_{1}=x_{2}$. We have

$$
x_{1}=i d_{X}\left(x_{1}\right)=g \circ f\left(x_{1}\right)=g\left(f\left(x_{1}\right)\right)=g(y)=g\left(f\left(x_{2}\right)\right)=g \circ f\left(x_{2}\right)=i d_{X}\left(x_{2}\right)=x_{2},
$$

i.e., $x_{1}=x_{2}$. This proves injectivity of $f$.

Remark 4.10. [Horizontal and vertical line tests] Let $X$ and $Y$ be nonempty sets and $f: X \rightarrow Y$. The following needs to be taken with a grain of salt because $X$ and $Y$ need not be sets of real numbers.

Let $R \subseteq X \times Y$.
a. (4.6) on p. 72 states that $R$ is the graph of a function with domain $X$ and codomain $Y$ if and only if it passes the "vertical line test": Any "vertical line", i.e., any subset of $X \times Y$ of the form $V\left(x_{0}\right):=\left\{\left(x_{0}, y\right): y \in Y\right\}$ for a fixed $x_{0} \in X$ intersects $R$ in exactly one point.
b. (4.12) on p. 79 states that $R$ is the graph of a surjective function with domain $X$ and codomain $Y$ if and only if it passes in addition to the vertical line test the following "horizontal line test": any "horizontal line", i.e., any subset of $X \times Y$ of the form $H\left(y_{0}\right):=\left\{\left(x, y_{0}\right): x \in X\right\}$ for a fixed $y_{0} \in Y$ intersects $R$ in at least one point.
c. (4.13) on p. 79 states that $R$ is the graph of an injective function with domain $X$ and codomain $Y$ if and only if it passes in addition to the vertical line test the following horizontal line test: any "horizontal line", i.e., any subset of $X \times Y$ of the form $H\left(y_{0}\right):=\left\{\left(x, y_{0}\right): x \in X\right\}$ for a fixed $y_{0} \in Y$ intersects $R$ in at most one point.
d. It follows from (4.14) on p. 79 but also from the above that that $R$ is the graph of a bijective function with domain $X$ and codomain $Y$ if and only if it passes in addition to the vertical line test the following horizontal line test: any "horizontal line", i.e., any subset of $X \times Y$ of the form $H\left(y_{0}\right):=\left\{\left(x, y_{0}\right): x \in X\right\}$ for a fixed $y_{0} \in Y$ intersects $R$ in exactly one "point". Note the symmetry between this test and the one for vertical lines. The above is another indication that the inverse graph $R^{-1}$ of a bijective function is a graph of a function (the inverse function $f^{-1}$ ).

Remark 4.11. Abstract math is about proving theorems and propositions. Functions are very important tools for many proofs. I many instances it is very important to know or to show that a certain function is injective or surjective or both. But these properties depend on the choice of domain and codomain as this simple example shows.

Let $f: A \rightarrow B$ be the function $f(x):=x^{2}$.

$$
\begin{array}{ll}
A=\mathbb{R}, B=\mathbb{R}: & f \text { is neither injective nor surjective } \\
A=]-2,3[, B=[0,9[: & f \text { is surjective but not injective } \\
A=] 0,3[, B=[0,9]: & f \text { is injective but not surjective } \\
A=] 0,3[, B=] 0,9[: & f \text { is bijective }
\end{array}
$$

The above demonstrates why domain and codomain are important for the specification of a function.

Proposition 4.1. Let $X, Y, Z \neq \emptyset$. Let $f: X \rightarrow Y$ and $g: Y \rightarrow Z$.
a. If both $f, g$ are injective then $g \circ f$ is injective.
b. If both $f, g$ are surjective then $g \circ f$ is surjective.
c. If both $f, g$ are bijective then $g \circ f$ is bijective.

The proof of $\boldsymbol{a}$ and $\boldsymbol{b}$ is left as exercise 4.6.
PROOF of $\boldsymbol{c}$ : Follows from $\boldsymbol{a}$ and $\boldsymbol{b}$ because bijective $=$ injective + surjective.
Corollary 4.1. Let $X, Y, Z \neq \emptyset$. Let $f: X \rightarrow Y$ and $g: Y \rightarrow Z$.
a. If $f$ is bijective and $g$ is injective then both $g \circ f$ and $f \circ g$ are injective.
b. If $f$ is bijective and $g$ is surjective then both $g \circ f$ and $f \circ g$ are surjective.
c. If $f$ is bijective and $g$ is bijective then both $g \circ f$ and $f \circ g$ are bijective.

## PROOF:

a follows from prop.4.1.a because bijective functions are injective.
$\boldsymbol{b}$ follows from prop.4.1.b because bijective functions are surjective.
c follows from prop.4.1.c.
We now examine conditions under which there are functions $f: X \rightarrow Y$ and $g: Y \rightarrow X$ such that $g \circ f=i d_{X}$, i.e.,

$$
\begin{equation*}
g(f(x))=x \text { for all } x \in X: \tag{4.20}
\end{equation*}
$$



Proposition 4.2. Let $X, Y \neq \emptyset$. Let $f: X \rightarrow Y$ and $g: Y \rightarrow X$ such that $g \circ f=i d_{X}$. Then
a. $f$ is injective,
b. $g$ is surjective.

PROOF of $a$ : Let $x_{1}, x_{2} \in X$. If $f\left(x_{1}\right)=f\left(x_{2}\right)$ then

$$
x_{1}=i d_{X}\left(x_{1}\right)=g\left(f\left(x_{1}\right)\right)=g\left(f\left(x_{2}\right)\right)=i d_{X}\left(x_{2}\right)=x_{2}
$$

This proves injectivity of $f$.
PROOF of $\boldsymbol{b}$ : Let $x_{0} \in X$. Let $y:=f\left(x_{0}\right)$. Then $g(y)=g\left(f\left(x_{0}\right)\right)=g \circ f\left(x_{0}\right)=x_{0}$. We found for an arbitrary $x_{0}$ in the codomain of $g$ some $y$ which maps to $x_{0}$ and we have proved surjectivity of $g$.

Proposition 4.3. Let $X, Y \neq \emptyset$.
a. Let $f: X \rightarrow Y$. If $f$ is injective then there exists $g: Y \rightarrow X$ such that $g \circ f=i d_{X}$ and any such function $g$ is necessarily surjective.
b. Let $g: Y \rightarrow X$. If $g$ is surjective then there exists $f: X \rightarrow Y$ such that $g \circ f=i d_{X}$ and any such function $f$ is necessarily injective.

PROOF of $a$ : Let $Y^{\prime}:=f(X)$ and

$$
f^{\prime}: X \rightarrow Y^{\prime}, \quad x \mapsto f(x)
$$

i.e., $f(x)=f^{\prime}(x)$ for all $x \in X$. The only difference between $f$ and $f^{\prime}$ is that we shrunk the codomain from $Y$ to $f(X)$, thus making $f^{\prime}$ not only injective but also surjective, hence bijective. It follows that the inverse $\left(f^{\prime}\right)^{-1}: Y^{\prime} \rightarrow X$ exists.

Let $x_{0}$ be an arbitrary, but fixed, element of $X$. We define $g: Y \rightarrow X$ as follows.

$$
g(y):= \begin{cases}\left(f^{\prime}\right)^{-1}(y) & \text { if } y \in Y^{\prime} \\ x_{0} & \text { if } y \notin Y^{\prime}\end{cases}
$$

Let $x \in X$. Then $f(x) \in Y^{\prime}$, hence $g \circ f(x)=g \circ f^{\prime}(x)=\left(f^{\prime}\right)^{-1}\left(f^{\prime}(x)\right)=x$. This proves $g \circ f=i d_{X}$. We observe that $g$ is surjective according to prop.4.2b.
PROOF of $\boldsymbol{b}$ : For $x \in X$ let $Y_{x}:=\{y \in Y: g(y)=x\}$. It follows from the surjectivity of $g$ that $Y_{x} \neq \emptyset$ and we may select for each $x \in X$ some $y_{x} \in Y_{x} .{ }^{46}$

Let $f: X \rightarrow Y$ be the function $x \mapsto y_{x}$. Let $x \in X$. Then

$$
g \circ f(x)=g\left(y_{x}\right)=x
$$

The first equality follows from the definition of $f$ and the second one is true because $y_{x} \in Y_{x}$. It follows from prop.4.2a that $f$ is injective.

There are special names for functions $f$ and $g$ which are related by (4.20).
Definition 4.10 (Left inverses and right inverses). Let $X, Y \neq \emptyset$. Let $f: X \rightarrow Y$ and $g: Y \rightarrow X$ such that $g \circ f=i d_{X}$. We call $g$ a left inverse of $f$ and we call $f$ a right inverse of $g$.

We combine that last definitions with the preceding two proposition and obtain

[^33]Theorem 4.2. Let $X, Y \neq \emptyset$.
a. Let $f: X \rightarrow Y$. Then $f$ is injective $\Leftrightarrow f$ has a left inverse (which is necessarily surjective).
b. Let $g: Y \rightarrow X$. Then $g$ is surjective $\Leftrightarrow g$ has a right inverse (which is necessarily injective).
c. An injection $X \rightarrow Y$ exists $\Leftrightarrow$ a surjection $Y \rightarrow X$ exists.

PROOF of $\boldsymbol{a} \Rightarrow)$ : prop.4.3a.
PROOF of $\boldsymbol{a} \Leftarrow)$ : prop.4.2a.
PROOF of $\boldsymbol{b} \Rightarrow)$ : prop.4.3b.
PROOF of $\boldsymbol{b} \Leftarrow)$ : prop.4.2b.
PROOF of $\boldsymbol{c} \Rightarrow)$ : Let $f: X \rightarrow Y$ be injective. According to part a there exists a left inverse $g: Y \rightarrow X$ and this function is surjective
PROOF of $\boldsymbol{c} \Leftarrow)$ : Let $g: Y \rightarrow X$ be surjective. According to part $\boldsymbol{b}$ there exists a right inverse $f: X \rightarrow Y$ and this function is injective

Remark 4.12. Let $X$ and $Y$ be two nonempty sets. No assumptions are made concerning how $X$ and $Y$ might be related.
a. Let $a \in Y$. Then the function

$$
\begin{equation*}
f: X \xrightarrow{\sim}\{a\} \times X ; \quad x \mapsto(a, x) \tag{4.21}
\end{equation*}
$$

is bijective because $f$ has the function $(a, x) \mapsto x$ as an inverse.
b. The sets $\{0\} \times X$ and $\{1\} \times Y$ are disjoint.
c. An injection/surjection/bijection $X \rightarrow Y$ exists if and only if an injection/surjection/bijection $\{0\} \times X \rightarrow\{1\} \times Y$ exists.
Definition 4.11 (Restriction/Extension of a function). Given are three nonempty sets $A, X$ and $Y$ such that $A \subseteq X$. Let $f: X \rightarrow Y$ a function with domain $X$. We define the restriction of $f$ to $A$ as the function

$$
\begin{equation*}
\left.f\right|_{A}: A \rightarrow Y \quad \text { defined as }\left.\quad f\right|_{A}(x):=f(x) \text { for all } x \in A \tag{4.22}
\end{equation*}
$$

Conversely let $f: A \rightarrow Y$ and $\varphi: X \rightarrow Y$ be functions such that $f=\left.\varphi\right|_{A}$. We then call $\varphi$ an extension of $f$ to $X$.
Example 4.20. For an example let $X:=\mathbb{R}, A:=[0,1]$ and $f(x):=3 x^{2}(x \in[0,1])$. For any $\alpha \in \mathbb{R}$ the function $\varphi_{\alpha}: \mathbb{R} \rightarrow \mathbb{R}$ defined as $\varphi_{\alpha}(x):=3 x^{2}$ if $0 \leqq x \leqq 1$ and $\alpha x$ otherwise defines a different extension of $f$ to $\mathbb{R}$.
Notations 4.1. As the only difference between $f$ and $\left.f\right|_{A}$ is the domain, it is customary to write $f$ instead of $\left.f\right|_{A}$ to make formulas look simpler if doing so does not give rise to confusions.
Remark 4.13. The restriction $\left.f\right|_{A}$ is always uniquely determined by $f$. Such is not the case for extensions if $A$ is a strict subset of $X$ unless some conditions are imposed on the nature of the extension.

For example, if we had asked for continuity ${ }^{47}$ of the extension $\varphi_{\alpha}$ of $f$ in example 4.20 above, only $\varphi_{1}(x)=3 x^{2}$ if $0 \leqq x \leqq 1$ and $x$ otherwise would qualify.

[^34]
### 4.2.5 Operations on Real Valued Functions

Definition 4.12 (real functions). Let $X$ be an arbitrary, nonempty set. If the codomain $Y$ of a mapping

$$
f: X \rightarrow Y ; \quad x \mapsto f(x)
$$

is a subset of $\mathbb{R}$, then we call $f(\cdot)$ a real function or real valued function.
Remember that this definition does not exclude the case $Y=\mathbb{R}$ because $Y \subseteq \mathbb{R}$ is in particular true if both sets are equal.

Real functions are a pleasure to work with because, given any fixed argument $x_{0}$, the object $f\left(x_{0}\right)$ is just an ordinary number. In particular you can add, subtract, multiply and divide real valued functions. Of course, division by zero is not allowed:

Definition 4.13 (Operations on real valued functions). Let $X$ be an arbitrary nonempty set.
Given are two real valued functions $f(\cdot), g(\cdot): X \rightarrow(R)$ and a real number $\alpha$. The sum $f+g$, difference $f-g$, product $f g$ or $f \cdot g$, quotient $f / g$, and scalar product $\alpha f$ are defined by doing the operation in question with the numbers $f(x)$ and $g(x)$ for each $x \in X$. In other words these items are defined by the following equations:

$$
\begin{align*}
(f+g)(x) & :=f(x)+g(x), \\
(f-g)(x) & :=f(x)-g(x), \\
(f g)(x) & :=f(x) g(x),  \tag{4.23}\\
(f / g)(x) & :=f(x) / g(x) \quad \text { for all } x \in X \text { where } g(x) \neq 0, \\
(\alpha f)(x) & :=\alpha \cdot g(x) .
\end{align*}
$$

Before we list some basic properties of addition and scalar multiplication of functions (the operations that interest us the most), let us have a quick look at constant functions.
Definition 4.14 (Constant functions). Let $Y$ be a nonempty set and $y_{0} \in Y$. You can think of $y_{0}$ as a function from any nonempty set $X$ to $Y$ as follows:

$$
y_{0}(\cdot): X \rightarrow Y ; \quad x \mapsto y_{0} .
$$

In other words, the function $y_{0}(\cdot)$ assigns to each $x \in X$ one and the same value $y_{0}$. We call such a function which only takes a single value a constant function.

The most important constant real-valued function is the zero function $0(\cdot)$ which maps any $x \in X$ to the number zero. We usually just write 0 for this function unless doing so would confuse the reader. Note that scalar multiplication $(\alpha f)(x)=\alpha \cdot g(x)$ is a special case of multiplying two functions $(g f)(x)=g(x) f(x)$ : Let $g(x)=\alpha$ for all $x \in X$ (constant function $\alpha$ ).

One last definition before we get so see some examples:
Definition 4.15 (Negative function). Let $X$ be an arbitrary, nonempty set and let $f: X \rightarrow \mathbb{R}$. The function

$$
-f(\cdot): X \rightarrow \mathbb{R} ; \quad x \mapsto-f(x) .
$$

is called negative $f$ or minus $f$. We usually write $-f$ for $-f(\cdot)$.

All those last definitions about sums, products, scalar products, . . of real valued functions are very easy to understand if you remember that, for any fixed $x \in X$, you just deal with ordinary numbers!

Example 4.21 (Arithmetic operations on real valued functions).
For simplicity, let $X:=\mathbb{R}_{+}=\{x \in \mathbb{R}: x \geqq 0\}$. Let

$$
\begin{array}{ll}
f: \mathbb{R}_{\geq 0} \rightarrow \mathbb{R} ; & x \mapsto(x-1)(x+1) \\
g: \mathbb{R}_{\geq 0} \rightarrow \mathbb{R} ; & x \mapsto(x-1) \\
h: \mathbb{R}_{\geqq 0} \rightarrow \mathbb{R} ; & x \mapsto(x+1)
\end{array}
$$

Then

$$
\begin{aligned}
(f+h)(x) & =(x-1)(x+1)+x+1=x^{2}-1+x+1=x(x+1) \forall x \in \mathbb{R}_{\geqq 0}, \\
(f-g)(x) & =(x-1)(x+1)-(x-1)=x^{2}-1-x+1=x(x-1) \forall x \in \mathbb{R}_{\geqq 0}, \\
(g h)(x) & =(x-1)(x+1)=f(x) \forall x \in \mathbb{R}_{\geqq 0}, \\
(f / h)(x) & =(x-1)(x+1) /(x+1)=x-1=g(x) \forall x \in \mathbb{R}_{\geqq 0}, \\
(f / g)(x) & =(x-1)(x+1) /(x-1)=x+1=h(x) \forall x \in \mathbb{R}_{\geqq 0} \backslash\{1\}
\end{aligned}
$$

It is really, really important for you to understand that $f / g$ and $h$ are not the same functions. Here is the reason. $f / g$ is not defined for $x=1$ because $\frac{(1-1)(1+1)}{1-1}=" 0 / 0$ ". The domain of $f / g$ is $\mathbb{R}_{\geqq 0} \backslash\{1\}$. It is different from $\mathbb{R}_{\geq 0}$, the domain of $h$. It follows that both functions are different.

Many more properties of functions will be discussed later. Now we look at families, sequences and some additional properties of sets.

### 4.2.6 Sequences, Families and Functions as Families

Definition 4.16 (Indices). Given is an expression of the form

$$
X_{\imath} .
$$

We say that $X$ is indexed by or subscripted by or tagged by $\imath$. We call $\imath$ the index or subsript of $X$ and we call $X_{i}$ an indexed item .

Remark 4.14. Both $X$ and $\imath$ can occur in many different ways. Here is a collection of indexed items:

$$
x_{7}, A_{\alpha}, k_{T}, \mathfrak{H}_{2 / 9}, f_{x}, x_{t}, h_{\mathscr{A}},, i_{\mathbb{R}}, H_{2 \pi}
$$

Some of the indices in this collection are highly unusual: Not only are some of them negative but they are fractions (e.g., 2/9) or irrational (e.g., $2 \pi$ ) Others don't even look like numbers (e.g., $\alpha, T, x$, $t, \mathscr{A}$ and $\mathbb{R}$ ). It is not clear from the information available to us whether those indices are names of variables which represent numbers or whether they represent functions, sets or other mathematical objects. There is one exception: The index $\mathbb{R}$ of $i_{\mathbb{R}}$ denotes the set of all real numbers.

We can turn any set into a "family" by tagging each of its members with an index. As an example, look at the following two indexed versions of the set $S_{2}$ from example 2.1 on $p .12$ :

$$
\begin{aligned}
& F=\left(a_{1}, e_{1}, e_{2}, i_{1}, i_{2}, i_{3}, o_{1}, o_{2}, o_{3}, o_{4}, u_{3}, u_{5}, u_{9}, u_{11}, u_{99}\right) \\
& G=\left(a_{k}, e_{-\sqrt{2}}, e_{1}, i_{-6}, i_{\mathscr{B}}, i_{\mathbb{R}}, o_{7}, o_{2 / 3}, o_{-8}, o_{3}, u_{A}, u_{B}, u_{C}, u_{D}, u_{E}\right)
\end{aligned}
$$

We note several things:
a. F has the kind of indices that we are familiar with: all of them are positive integers.
b. Some of the indices in $F$ occur multiple times. For example, 3 occurs as an index for $i_{3}, o_{3}, u_{3}$.
c. All of the indices in $G$ are unique.
d. As in remark 4.14, some of the indices are very unusual.

The last point is not much of a problem as mathematicians are used to very unusual notation but point (b), the non-uniqueness of indices, is something that we want to avoid. From now on we ask for the following: The indices of an indexed collection must belong to some set $J$ and each index $\imath \in J$ must be used exactly once. Remember that this automatically takes care of the duplicate indices problem as a set never contains duplicate values (see def.2.1 on $p .12$ ). We also demand that there is a set $X$ such that each indexed item $x_{\imath}$ belongs to $X$.

We now are ready to give the definition of a family:
Definition 4.17 (Indexed families). Let $J$ and $X$ be nonempty set and assume that
for each $\imath \in J$ there exists exactly one indexed item $x_{\imath} \in X$.
Let $R:=\left\{\left(\imath, x_{\imath}\right): \imath \in J\right\}$. Then $R$ is a relation on $(J, X)$ which satisfies (4.6) of the definition of a function

$$
F: J \rightarrow X, \quad \imath \mapsto F(\imath):=x_{\imath}
$$

(see def.4.6 on p.72) whose graph $\Gamma_{F}$ equals $R$.
We write $\left(x_{\imath}\right)_{\imath \in J}$ for this function if we want to emphasize that we are interested in the collection of indexed elements $x_{\imath}$ rather than the function $F(\cdot)$ or the relation $R$. Reasons for this will be given in rem.4.16 on p.87.
$\left(x_{\imath}\right)_{\imath \in J}$ is called an indexed family or simply a family in $X$ and $J$ is called the index set of the family. For each $\jmath \in J, x_{\jmath}$ is called a member of the family $\left(x_{\imath}\right)_{\imath \in J}$.
$\imath$ is a dummy variable: $\left(x_{\imath}\right)_{\imath \in J}$ and $\left(x_{k}\right)_{k \in J}$ describe the same family as long as $\imath \mapsto x_{\imath}$ and $k \mapsto x_{k}$ describe the same function $F: J \rightarrow X$. This should not come as a surprise to you if you recall remark 4.8 on p.73.

Remark 4.15. The codomain $X$ does not occur in the notation $\left(x_{\imath}\right)_{\imath \in J}$. This is not a problem because we do not care about surjectivity or injectivity of families. The only thing that matters about the set $X$ is that it is big enough to contain each indexed item. Here are two possible choices for a codomain.
a. If there is a universal set $X$ which contains all tagged items of the family then $X$ is acceptable as the codomain.
b. If there is no universal set then you can think of

$$
X=\bigcup\left[x_{\imath}: \imath \in J\right]:=\left\{x: x=x_{\imath_{0}} \text { for some } \imath_{0} \in I\right\}
$$

as the codomain. ${ }^{48}$

[^35]Definition 4.18 (Equality of families). Two families $\left(x_{i}\right)_{i \in I}$ and $\left(y_{j}\right)_{j \in J}$ are equal if
a. $\quad I=J$,
b. $\quad x_{i}=y_{i}$ for all $i \in I$.

Note 4.1 (Simplified notation for families).
If there is no confusion about the index set then it can be dropped from the specification of a family and we simply write $\left(x_{\imath}\right)_{\imath}$ instead of $\left(x_{\imath}\right)_{\imath \in J}$. Even the index outside the right parentheses may be omitted and we write $\left(x_{\imath}\right)$ if it is clear that we are talking about families.

For example, a proposition may start as follows: Let $\left(A_{\alpha}\right)$ and $\left(B_{\alpha}\right)$ be two families of subsets of $\Omega$ indexed by the same set. Then .....

It is clear from the formulation that we deal in fact with two families $\left(A_{\alpha}\right)_{\alpha \in J}$ and $\left(B_{\alpha}\right)_{\alpha \in J}$. Nothing is said about $J$, probably because the proposition is valid for any index set or because this set was fixed once and for all earlier on for the entire section.

Example 4.22. Here is an example of a family of subsets of $\mathbb{R}$ which are indexed by real numbers: Let $J=[0,1]$ and $X:=2^{\mathbb{R}}$. For $0 \leqq x \leqq 1$ let $A_{x}:=[x, 2 x]$ be the set of all real numbers between $x$ and $2 x$. Then $\left(A_{x}\right)_{x \in[0,1]}$ is such a family.

Remark 4.16. If a family is just some kind of function, why bother with yet another definition? The answer to this is that there are many occasions in which writing something as a collection of indexed items rather than as a function makes things easier to understand. For example, take a peek at theorem 5.1 (De Morgan's Law) on p.97. One of the formulas there states that for any indexed family $\left(A_{\alpha}\right)_{\alpha \in I}$ of subsets of a universal set $\Omega$ it is true that

$$
\left(\bigcup_{\alpha} A_{\alpha}\right)^{\complement}=\bigcap_{\alpha} A_{\alpha}^{\complement}
$$

Without the notion of a family you might have to say something like this: Let $A: I \rightarrow 2^{\Omega}$ be a function which assigns its arguments to subsets of $\Omega$. Then

$$
\left(\bigcup_{\alpha} A(\alpha)\right)^{\complement}=\bigcap_{\alpha} A(\alpha)^{\complement} .
$$

The additional parentheses around the index $\alpha$ just add complexity to the formula.
Example 4.23 (Sequences as families). You have worked with special families before: those where $J=\mathbb{N}$ or $J=\mathbb{Z}_{\geq 0}$ and $X$ is a subset of the real numbers. Example: $x_{n}:=1 / n$. Here

$$
\left(x_{n}\right)_{n \in \mathbb{N}} \text { corresponds to the indexed collection } 1, \frac{1}{2}, \frac{1}{3}, \ldots, \frac{1}{n}, \ldots
$$

## Such families are called sequences:

Definition 4.19 (Sequences and subsequences). Let $n_{\star}$ be an integer and let $J:=\left\{n_{\star}, n_{\star}+1, n_{\star}+\right.$ $2, \ldots\}=\left\{k \in \mathbb{Z}: k \geq n_{\star}\right\}$. Let $X$ be an arbitrary nonempty set. An indexed family in $X$ with index set $J$ is called a sequence in $X$ with start index $n_{\star}$.

Let $I$ be an infinite subset of $J$. Because any two different integers $m$ and $n$ can be compared, i.e., either $m<n$ or $n<m$, ${ }^{49}$ there are integers

$$
\begin{equation*}
n_{1}<n_{2}<n_{3}<\ldots \tag{4.24}
\end{equation*}
$$

such that $I=\left\{n_{j}: j \in \mathbb{N}\right\}$. An indexed family in $X$ with index set $I$ is called a subsequence of the original sequence.

As for families, the name of the index variable of a sequence does not matter as long as it is applied consistently. It does not matter whether you write $\left(x_{j}\right)_{j \in J}$ or $\left(x_{n}\right)_{n \in J}$ or $\left(x_{\beta}\right)_{\beta \in J}$.

Note 4.2 (Simplified notation for sequences).
a. It is customaryx to choose either of $i, j, k, l, m, n$ as the symbol of the index variable of a sequence and to stay away from other symbols whenever possible.
b. If $J$ is defined as above then is not unusual to see " $\left(x_{n}\right)_{n \geq n_{\star}}$ " instead of " $\left(x_{n}\right)_{n \in J}$ ". By default the index set for a sequence is $J=\mathbb{N}=\{1,2,3,4, \ldots\}$ and we can then write $\left(x_{n}\right)_{n}$ or just $\left(x_{n}\right)$
c. Customary notation for subsequences is either of $\left(x_{n_{j}}\right)_{j \in \mathbb{N}},\left(x_{n_{j}}\right)_{j \geqq 1},\left(x_{n_{j}}\right)_{j}$ or simply ( $x_{n_{j}}$ ).

Compare this to note 4.1 about simplified notation for families.
Here are two more examples of sequences.
Example 4.24 (Oscillating sequence). If $x_{j}:=(-1)^{j}\left(j \in \mathbb{Z}_{\geqq 0}\right)$, then

$$
x_{0}=1, \quad x_{2}=-1, \quad x_{2}=1, \quad x_{3}=-1, \quad x_{4}=1, \quad x_{5}=-1, \ldots
$$

Example 4.25 (Series (summation sequence) ). Let $s_{k}:=1+2^{-1}+2^{-2}+\ldots+2^{-k}(k=1,2,3, \ldots)$ :

$$
\begin{aligned}
& s_{1}=1, \quad s_{2}=1+1 / 2=2-1 / 2, \quad s_{3}=1+1 / 2+1 / 4=2-1 / 4, \quad \ldots, \\
& s_{k}=1+1 / 2+\ldots+2^{k-1}=2-2^{k-1} ; \quad s=1+1 / 2+1 / 4+1 / 8+\ldots \quad \text { "infinite sum". }
\end{aligned}
$$

You obtain $s_{k+1}$ from $s_{k}=2-2^{k-1}$ by cutting the difference $2^{k-1}$ to the number 2 in half (that would be $2^{k}$ ) and adding that to $s_{k}$. It is intuitively obvious from $s_{k}=2-2^{k-1}$ that the infinite sum $s$ adds up to 2 . Such an infinite sum is called a series. ${ }^{50}$

Remark 4.17. Having defined the family $\left(x_{\imath}\right)_{\imath \in J}$ as the function which maps $\imath \in J$ to $x_{\imath}$ means that a family distinguishes any two of its members $x_{\imath}$ and $x_{\jmath}$ by remembering what their indices are, even if they represent one and the same element of $X$ : Think of " $\left(x_{\imath}\right)_{\imath \in J \text { " as an abreviation for }}$

$$
\begin{equation*}
"\left(\left(\imath, x_{\imath}\right)\right)_{\imath \in J} " \text {. } \tag{4.25}
\end{equation*}
$$

Doing so should also make it much easier to see the equivalence of functions and families: (4.25) looks at its core very much like the graph $\left\{\left(\imath, x_{\imath}\right): \imath \in J\right\}$ of the function $\imath \mapsto x_{\imath}$.

[^36]Remark 4.18 (Families and sequences can contain duplicates). One of the important properties of sets is that they do not contain any duplicates (see def.2.1 (sets) on p.12). On the other hand, remark 4.17 casually mentions that families, and hence sequences as special kinds of families, can contain duplicates. Let us look at this more closely.

The two sets $A:=\{31,20,20,20,31\}$ and $B:=\{20,31\}$ are equal. On the other hand let $J:=\{\alpha, \beta, \pi, \star, Q\}$ and define the family $\left(w_{\imath}\right)_{\imath \in J}$ in $B$ by its associated graph as follows:

$$
\Gamma:=\{(\alpha, 31),(\beta, 20),(\pi, 20),(\star, 20),(Q, 31)\} \text {, i.e. }, \text {, } w_{\alpha}=31, w_{\beta}=20,, w_{\pi}=20, w_{\star}=20, w_{Q}=31 .
$$

The three occurrences of 20 cannot be distinguished as elements of the set $A$. In contrast to this the items $(\beta, 20),(\pi, 20),(\star, 20)$ as elements of $\Gamma \subseteq J \times A=J \times B{ }^{51}$ are different from each other because two pairs $(a, b)$ and $(x, y)$ are equal only if $x=a$ and $y=b$.

## In contrast to sets, families and sequences allow us to incorporate duplicates.

A family $\left(x_{\imath}\right)_{\imath \in J}$ in $X$ is specified by the function $F: J \rightarrow X$ which maps $\imath \in J$ to $F(z)=x_{z}$. Conversely, let $X, Y$ be nonempty sets and let $f: X \rightarrow Y$ be a function with domain $X$ and codomain $Y$. For $x \in X$ let $f_{x}:=f(x)$. Then $f$ can be written as $\left(f_{x}\right)_{x \in X}$, i.e., as a family in $Y$ with index set $X$. Thus we have

Proposition 4.4. The following ways of specifying a function $f: X \rightarrow Y, x \mapsto f(x)$ are equivalent:
a. $f=(X, Y, \Gamma)$ is defined by its graph $\Gamma:=\{(x, f(x)): x \in X\}$.
b. $f$ is defined by the following family in $Y:(f(x))_{x \in X}$

Note that this is one case where we had to explicitly mention the codomain $Y$ in the specification of the family.

There will be a lot more on sequences and series (sequences of sums) in later chapters, but we need to develop more concepts, such as convergence, to continue with this subject. Now let's get back to sets.

### 4.3 Exercises for Ch. 4

### 4.3.1 Exercises for Functions and Relations

## Exercise 4.1.

a. Which of the following is an equivalence relation? a partial ordering? on $\mathbb{R}$ ?
a1. $x R y \Leftrightarrow x<y$,
a2. $x R y \Leftrightarrow x \leqq y$,
a3. $x R y \Leftrightarrow x=y$,
a4. $x R y \Leftrightarrow x \neq y$.
b. Define $x R y \Leftrightarrow x y>0$. Is this an equivalence relation on $\mathbb{R}$ ? on $\mathbb{R}_{\neq 0}$ ? on $\mathbb{R}_{>0}$ ? on $\mathbb{R}_{<0}$ ?

Exercise 4.2. Injectivity and Surjectivity

- Let $f: \mathbb{R} \rightarrow\left[0, \infty\left[; \quad x \mapsto x^{2}\right.\right.$.
- Let $g:\left[0, \infty\left[\rightarrow\left[0, \infty\left[; \quad x \mapsto x^{2}\right.\right.\right.\right.$.

In other words, $g$ is same function as $f$ as far as assigning function values is concerned, but its domain is downsized to $[0, \infty[$.
Answer the following with true or false.

[^37]a. $\quad f$ is surjective
b. $f$ is injective
c. $g$ is surjective
d. $g$ is injective

If your answer is false then give a specific counterexample.
Exercise 4.3 (Excercise 4.2 continued). Let $A \subseteq \mathbb{R}$.
Part 1.

- Let $F_{1}: A \rightarrow\left[-2,20\left[; \quad x \mapsto x^{2}\right.\right.$.
- Let $F_{2}: A \rightarrow\left[2,20\left[; \quad x \mapsto x^{2}\right.\right.$.

What choice of $A$ makes
a. $F_{1}$ surjective?
b. $\quad F_{1}$ injective?
c. $F_{2}$ surjective?
d. $F_{2}$ injective?

Part 2.

- Let $G_{1}: A \rightarrow[-2,20[; \quad x \mapsto \sqrt{x}$.
- Let $G_{2}: A \rightarrow[2,20[; \quad x \mapsto \sqrt{x}$.

What choice of $A$ makes
e. $G_{1}$ surjective?
f. $G_{1}$ injective?
g. $\quad G_{2}$ surjective?
h. $\quad G_{2}$ injective?

Part 3.

- Let $G_{3}: A \rightarrow[-20,2[; \quad x \mapsto \sqrt{x}$.
- Let $G_{4}: A \rightarrow[-20,-2[; \quad x \mapsto \sqrt{x}$.

What choice of $A$ makes
i. $G_{3}$ surjective? k. $\quad G_{3}$ surjective?
j. $G_{4}$ injective? 1. $G_{4}$ injective?

For the questions above

- Write impossible if no choice of $A \subseteq \mathbb{R}$ exists.
- Write NAF for any of $F_{1}, F_{2}, G_{1}, G_{2}, G_{3}, G_{4}$ which does not define a function.

Exercise 4.4. Find $f: X \rightarrow Y$ and $A \subseteq X$ such that $f\left(A^{\complement}\right) \neq f(A)^{\complement}$. Hint: use $f(x)=x^{2}$ and choose $Y$ as a one element only set (which does not leave you a whole lot of choices for $X$ ). See example 4.17 on p.77.

## Exercise 4.5.

a. Prove prop.4.1a: The composition of two injective functions is injective.
b. Prove prop.4.1b: The composition of two surjective functions is surjective.

Exercise 4.6. You proved in the previous exercise that
injective $\circ$ injective $=$ injective,
surjective $\circ$ surjective $=$ surjective.
This exercise illustrates that the reverse is not necessarily true.
Find functions $f:\{a\} \rightarrow\left\{b_{1}, b_{2}\right\}$ and $g:\left\{b_{1}, b_{2}\right\} \rightarrow\{a\}$ such that $h:=g \circ f:\{a\}$ is bijective but such that it is not true that both $f, g$ are injective and it is also not true that both $f, g$ are surjective.

Hint: There are not a whole lot of possibilities. Draw possible candidates for $f$ and $g$ in arrow notation as on p.118. You should easily be able to figure out some examples. Think simple and look at example 4.17 on p.77.

Exercise 4.7. Prove $\mathbf{b}$ of remark 4.12 on p .83 : Let $X$ and $Y$ be two nonempty sets. Then the sets $\{0\} \times X$ and $\{1\} \times Y$ are disjoint.

Exercise 4.8. Prove c of remark 4.12 on p.83: Let $X$ and $Y$ be two nonempty sets.
Then an injection/surjection/bijection $X \rightarrow Y$ exists if and only if an injection/surjection/bijection $\{0\} \times X \rightarrow\{1\} \times Y$ exists.

Exercise 4.9. B/G Project 6.9.:
On $\mathbb{Z} \times(\mathbb{Z} \backslash\{0\})$ we define the relation $\sim$ as follows.

$$
\begin{equation*}
\left(m_{1}, n_{1}\right) \sim\left(m_{2}, n_{2}\right) \Leftrightarrow m_{1} \cdot n_{2}=n_{1} \cdot m_{2} . \tag{4.26}
\end{equation*}
$$

a. Prove that $\sim$ defines an equivalence relation on $\mathbb{Z} \times(\mathbb{Z} \backslash\{0\})$.

## Let

$$
\begin{equation*}
\mathfrak{Q}:=\{[(m, n)]: m, n \in \mathbb{Z} \text { and } n \neq 0\} \tag{4.27}
\end{equation*}
$$

be the set of all equivalence classes of $\sim$. We define two binary operations $\oplus$ and $\otimes$ on $\mathfrak{Q}$ as follows;

$$
\begin{align*}
& {\left[\left(m_{1}, n_{1}\right)\right] \oplus\left[\left(m_{2}, n_{2}\right)\right]:=\left[\left(m_{1} n_{2}+m_{2} n_{1}, n_{1} n_{2}\right)\right],}  \tag{4.28}\\
& {\left[\left(m_{1}, n_{1}\right)\right] \otimes\left[\left(m_{2}, n_{2}\right)\right]:=\left[\left(m_{1} m_{2}, n_{1} n_{2}\right)\right]} \tag{4.29}
\end{align*}
$$

b. Prove that these binary operations are defined consistently: the right hand sides of (4.28) and (4.29) do not depend on the particular choice of elements picked from the sets $\left[\left(m_{1}, n_{1}\right)\right]$ and $\left[\left(m_{2}, n_{2}\right)\right]$. In other words, prove the following:

Let $\left(p_{1}, q_{1}\right) \sim\left(m_{1}, n_{1}\right)$ and $\left(p_{2}, q_{2}\right) \sim\left(m_{2}, n_{2}\right)$. Then

$$
\begin{align*}
{\left[\left(m_{1} n_{2}+m_{2} n_{1}, n_{1} n_{2}\right)\right] } & =\left[\left(p_{1} q_{2}+p_{2} q_{1}, q_{1} q_{2}\right)\right],  \tag{4.30}\\
{\left[\left(m_{1} m_{2}, n_{1} n_{2}\right)\right] } & =\left[\left(p_{1} p_{2}, q_{1} q_{2}\right)\right] . \tag{4.31}
\end{align*}
$$

or, equivalently, then

$$
\begin{align*}
\left(m_{1} n_{2}+m_{2} n_{1}, n_{1} n_{2}\right) & \sim\left(p_{1} q_{2}+p_{2} q_{1}, q_{1} q_{2}\right),  \tag{4.32}\\
\left(m_{1} m_{2}, n_{1} n_{2}\right) & \sim\left(p_{1} p_{2}, q_{1} q_{2}\right) . \tag{4.33}
\end{align*}
$$

### 4.4 Addenda to Ch. 4

This material belongs after def.4.3 (Equivalence relations and equivalence classes) on p. 68 .

Proposition 4.5 (see [1] B/G prop.6.4). Let " $\sim$ " be an equivalence relation on a nonempty set $X$ and $x, y \in X$ Then
a. $x \in[x]$.
b. $x \sim y \Leftrightarrow[x]=[y]$.

PROOF of a: This follows from the reflexivity of " $\sim$ ".
PROOF of b: Let $z \in[x]$. It follows from the definition of $x$ that $z \sim x$, hence $z \sim y$ (transitivity of " $\sim^{\prime \prime}$ ), hence $z \in[y]$. This proves $[x] \subseteq[y]$.

We switch the roles of $x$ and $y$ and repeat the above to obtain $[y] \subseteq[x]$.

Proposition 4.6 (see [1] B/G prop.6.5). Let " $\sim$ " be an equivalence relation on a nonempty set $X$ and $x, y \in X$ Then either $[x]=[y]$ or $[x] \cap[y]=\emptyset$.

PROOF: It suffices to show that if $[x] \cap[y] \neq \emptyset$ then $[x]=[y]$. So let $z \in[x] \cap[y]$. It follows from $z \in[x]$ that $z \sim x$ and from $z \in[y]$ that $z \sim y$. Symmetry and transitivity of " $\sim$ " yield $x \sim y$. It follows from prop.4.5 that $[x]=[y]$.

## Proposition 4.7 (see [1] B/G prop.6.6 for parts a and b).

a. Let " $\sim$ " be an equivalence relation on a nonempty set $X$ and let $\mathscr{P}_{\sim}:=\{[x]: x \in X\}$ be the set of all its equivalence classes. Then $\mathscr{P}_{\sim}$ is a partition of $X$.
b. Conversely, let $\mathscr{P}$ be a partition of $X$ and define a relation " $\sim_{\mathscr{P}}$ " on $X$ as follows: $x \sim \mathscr{P} y \Leftrightarrow$ there is $P \in \mathscr{P}$ such that $x, y \in P$. Then $\sim \mathscr{P}$ is an equivalence relation on $X$.
c. Let " $\sim$ " be an equivalence relation on $X$. Let $\mathscr{P}_{\sim}$ be the associated partition of its equivalence classes. Let " $\sim_{\mathscr{P}_{\sim}}$ " be the equivalence relation associated with the partition $\mathscr{P}_{\sim}$. Then " $\sim_{\mathscr{P}_{\sim}}$ " $=$ " $\sim$ " (i.e., both equivalence relations are equal as subsets of $X \times X$.
d. Let $\mathscr{P}$ be a partition of $X$. Let $\sim_{\mathscr{P}}$ be the associated equivalence relation defined in part b.. Let $\mathscr{P}_{\sim_{\mathscr{P}}}$ be the associated partition of its equivalence classes. Then $\mathscr{P}_{\sim_{\mathscr{P}}}=\mathscr{P}$.

## PROOF of $\boldsymbol{a}$ :

We observe that a set contains no duplicates: If $[x]=[y]$, we do not count $[x]$ and $[y]$ as separate members of $\mathscr{P}_{\sim}$ ! It follows from prop.4.6 that those members are mutually disjoint.

It remains to prove that their union is $X$. Let $x \in X$. Then $x \in[x]$ (reflexivity of " $\sim$ ") and $[x] \in \mathscr{P}_{\sim}$. It follows that $x \in \bigcup\left[P: P \in \mathscr{P}_{\sim}\right]$ and we obtain that $\bigcup\left[P: P \in \mathscr{P}_{\sim}\right]=X$. We have proved that $\mathscr{P}_{\sim}$ is a partition of $X$.
PROOF of $\boldsymbol{b}$ : For the following assume that $x, y, z \in X$.
It follows from $\biguplus[P: P \in \mathscr{P}]=X$ that for each $x \in X$ there exists $P \in \mathscr{P}$ such that $x \in P$. It further follows from the mutual disjointness of the elements of $\mathscr{P}$ that there exists exactly one such $P$ and we are justfied to write $P_{x}$ for this uniquely defined set $P$. In other words, the assignment $x \mapsto P_{x}$ defines a function $P(\cdot): X \rightarrow \mathscr{P}$.

Reflexivity: $x \in P_{x}$ implies $x \sim_{P} x$.
Symmetry: Let $x \sim_{\mathscr{P}} y$. This implies $P_{x}=P_{y}$, hence $y \in P_{x}$, hence $y \sim_{\mathscr{P}} x$.
Transitivity: Let $x \sim_{\mathscr{P}} y$ and Let $y \sim_{\mathscr{P}} z . x \sim_{\mathscr{P}} y$ implies $P_{x}=P_{y}$ and $y \sim_{\mathscr{P}} z$ implies $P_{y}=P_{z}$. It follows that $P_{z}=P_{x}$, i.e., $x \sim \mathscr{P} z$.
PROOF of $c$ :
Let $x, y \in X$. Then

$$
x \sim y \Leftrightarrow[x]=[y] \Leftrightarrow \text { both } x, y \text { belong to the same element of } \mathscr{P}_{\sim} \Leftrightarrow x \sim_{\mathscr{P}_{\sim}} y .
$$

## PROOF of $d$ :

Let $P \in \mathscr{P}$ and $x, y \in X$. Let $[x]$ and $[y]$ be the equivalence classes of $x$ and $y$ for " $\sim_{\mathscr{P}}$ ". If $x \in P$ then

$$
y \in P \Leftrightarrow x \sim_{\mathscr{P}_{\sim}} y \Leftrightarrow[x]=[y] .
$$

It follows that $P=[x]$, hence $P \in \mathscr{P}_{\sim_{\mathscr{P}}}$ This true for any $P \in \mathscr{P}$ and it follows that $\mathscr{P} \subseteq \mathscr{P}_{\sim_{\mathscr{P}}}$.
Now let $x \in X$. If $y \in X$ then

$$
y \in[x] \Leftrightarrow y \sim_{\mathscr{P}} x \Leftrightarrow y \in P_{x} .
$$

It follows that any equivalence class $[x]$ with respect to " $\sim_{\mathscr{P}}$ " is an element of $\mathscr{P}$, hence $\mathscr{P}_{\sim_{\mathscr{P}}} \subseteq \mathscr{P}$. We have shown that $\mathscr{P}_{\sim_{\mathscr{P}}}=\mathscr{P}$ and this finishes the proof of Proof of $d$.

This material belongs at the begining of ch.5.2 (Cartesian Products of more than Two Sets) on p. 99 .

For the following definition see also [1] B/G (Beck/Geoghegan), ch.4.3 (Fishing in a Finite Pool).
Definition 4.20 (Vectors and finite lists). Let $n_{\star}, n^{\star} \in \mathbb{Z}$ such that $n_{\star} \leqq n^{\star}$ and let $J:=\left\{k \in \mathbb{Z}: n_{\star} \leqq k \leqq n^{\star}\right\}$. A finite list $\chi$, is a family with (finite) index set $J$. It is customary to write one of the following:

$$
\begin{align*}
& \chi=x_{n_{\star}}, x_{n_{\star+1}}, x_{n_{\star}+2}, \ldots, x_{n^{\star-1}}, x_{n^{\star}},  \tag{4.34}\\
& \chi=\left(x_{n_{\star}}, x_{n_{\star+1}}, x_{n_{\star}+2}, \ldots, x_{n^{\star-1}}, x_{n^{\star}}\right),  \tag{4.35}\\
& \chi=\left(x_{k}\right)_{n_{\star}^{\star}}^{n^{\star}} . \tag{4.36}
\end{align*}
$$

If all members $x_{j}$ of the list are (real) numbers, we call $\chi$ a vector ${ }^{52}$ of dimension $n^{\star}$. In this case we always surround the list with parentheses, and often use "arrow notation" $\vec{x}$ rather than $\chi$, i.e., we write

$$
\begin{equation*}
\vec{x}=\left(x_{1}, x_{2}, x_{3}, \ldots, x_{n-1}, x_{n}\right) . \tag{4.37}
\end{equation*}
$$

## Example 4.26.

a. $(3.5,-97 \pi, 4, \sqrt{8})$ is both a list and a vector of dimension 4 . We have $x_{1}=3.5, x_{2}-97 \pi, x_{3}=4, x_{4}=\sqrt{8}$.
2. $(3 k+2)_{k=-2}^{3}=-4,-1,2,5,8,11=(-4,-1,2,5,8,11)$ is both a list and a vector of size/dimension 6 .
3. Joe, $5,-6.8$, Dolores is a list of size 4 . This is not a vector because not all of its members are numeric.

This material belongs at the end of ch.4.2.4 (Injective, Surjective and Bijective functions) on p.79.

Definition 4.21 (Binary operations). Let $X$ be a nonempty set. A binary operation, on $X$ is a function

$$
\begin{equation*}
\diamond: X \times X \longrightarrow X ; \quad(x, y) \mapsto x \diamond y:=\diamond(x, y) \tag{4.38}
\end{equation*}
$$

A unary operation, on $X$ is a function

$$
\begin{equation*}
\bullet: X \longrightarrow X ; \quad x \mapsto x^{\bullet}:=\bullet(x) \tag{4.39}
\end{equation*}
$$

[^38]Example 4.27. Some examples of binary operations are.
a. Addition on $X=\mathbb{N}$ or $X=\mathbb{Z}$ or $X=\mathbb{Q}$ or $X=\mathbb{R}$ is a binary operation

$$
\begin{equation*}
+: X \times X \longrightarrow X ; \quad(x, y) \mapsto x+y \tag{4.40}
\end{equation*}
$$

b. Multiplication on $X=\mathbb{N}$ or $X=\mathbb{Z}$ or $X=\mathbb{Q}$ or $X=\mathbb{R}$ is a binary operation

$$
\begin{equation*}
\cdot: X \times X \longrightarrow X ; \quad(x, y) \mapsto x \cdot y . \tag{4.41}
\end{equation*}
$$

c. Let $X$ be a onempty set and $\mathscr{F}:=\{$ functions $f: X \rightarrow X\}$. Function composition

$$
\begin{equation*}
\circ: \mathscr{F} \times \mathscr{F} \longrightarrow \mathscr{F} ; \quad(f, g) \mapsto g \circ f \tag{4.42}
\end{equation*}
$$

where $g \circ f: X \rightarrow X$ is the function defined by $x \mapsto g \circ f(x):=g(f(x))(x \in X)$.
Some examples of unnary operations are.
a. Negative number: Let $X=\mathbb{N}$ or $X=\mathbb{Z}$ or $X=\mathbb{Q}$ or $X=\mathbb{R}$. Then

$$
\begin{equation*}
-: X \rightarrow X ; \quad x \mapsto-x . \tag{4.43}
\end{equation*}
$$

is a unary operation.
b. Reciprocal $X=\mathbb{Q}_{\neq 0}$ or $X=\mathbb{R}_{\neq 0}$

$$
\begin{equation*}
{ }^{-1}: X \rightarrow X ; \quad x \mapsto 1 / x . \tag{4.44}
\end{equation*}
$$

is a unnary operation.
c. Let $X$ be a onempty set and $\mathscr{B}:=\{$ bijective functions $f: X \rightarrow X\}$. Let

$$
\begin{equation*}
{ }^{-1}: \mathscr{B} \rightarrow \mathscr{B} ; f(\cdot) \mapsto f^{-1}(\cdot) \tag{4.45}
\end{equation*}
$$

be the function which assigns to the function $x \mapsto f(x)$ its (uniquely determined) inverse function $\left.y \mapsto f^{-1}(y)\right)$. Then this assignment is a unary operation on $\mathscr{B}$.

## 5 More on Sets (Understand this!)

### 5.1 More on Set Operations (Study this!)

The material in this chapter thematically belongs to ch. 2.1 on $p .12$ but it had to be deferred to this chapter as much of it deals with families of sets, i.e., families $\left(A_{i}\right)_{i}$

Definition 5.1 (Arbitrary unions and intersections). Let $J$ be a nonempty set and let $\left(A_{i}\right)_{i \in J}$ be a family of sets. We define

$$
\begin{align*}
& \bigcup_{i \in I} A_{i}:=\bigcup\left[A_{i}: i \in I\right]:=\left\{x: x \in A_{i_{0}} \text { for some } i_{0} \in I\right\},  \tag{5.1}\\
& \bigcap_{i \in I} A_{i}:=\bigcap\left[A_{i}: i \in I\right]:=\left\{x: x \in A_{i_{0}} \text { for each } i_{0} \in I\right\} . \tag{5.2}
\end{align*}
$$

We call $\bigcup_{i \in I} A_{i}$ the union and $\bigcap_{i \in I} A_{i}$ the intersection of the family $\left(A_{i}\right)_{i \in J}$
It is convenient to allow unions and intersections for the empty index set $J=\emptyset$. For intersections this requires the existence of a universal set $\Omega$. We define

$$
\begin{equation*}
\bigcup_{i \in \emptyset} A_{i}:=\emptyset, \quad \bigcap_{i \in \emptyset} A_{i}:=\Omega \tag{5.3}
\end{equation*}
$$

Note that any statement concerning arbitrary families of sets such as the definition above covers finite lists $A_{1}, A_{2}, \ldots, A_{n}$ of sets ( $J=\{1,2, \ldots, n\}$ ) and also sequences $A_{1}, A_{2}, \ldots$, of sets ( $J=\mathbb{N}$ ).

We give some examples of non-finite unions and intersections. For the first one (prop.5.1) We need to review the behavior of the sequence $x_{n}=1 / n$.

Remark 5.1. The following definition of the limit of a sequence is known from calculus (see [11] Stewart, J: Single Variable Calculus, 7th edition, ch.11.1) and it will be repeated in ch.8.2 (Convergence and Continuity in $\mathbb{R}$ ) of this document (see def.8.10 on p.133).

We say that a sequence $\left(x_{n}\right)$ of real numbers converges to $a \in \mathbb{R}$ as $n \rightarrow \infty$ if for any $\varepsilon \in \mathbb{R}$ (no matter how small) there exists a corresponding (possibly extremely large) $N \in \mathbb{N}$ such that

$$
\begin{equation*}
\left|x_{j}-a\right|<\varepsilon \text { for all } j \geqq N . \tag{5.4}
\end{equation*}
$$

We write

$$
\begin{equation*}
a=\lim _{n \rightarrow \infty} x_{n} \quad \text { or } \quad x_{n} \rightarrow a \tag{5.5}
\end{equation*}
$$

and we call $a$ the limit of the sequence $\left(x_{n}\right)$.
Remark 5.2. It is known from calculus that $\lim _{n \rightarrow \infty} \frac{1}{n}=0$. According to (5.4) this implies that for any $\varepsilon>0$ there exists $N \in \mathbb{N}$ such that

$$
\begin{equation*}
\frac{1}{n}=\left|\frac{1}{n}-0\right|<\varepsilon \text { for all } n \geqq N \tag{5.6}
\end{equation*}
$$

Proposition 5.1. For the following note that $[u, v]=\emptyset$ for $u>v$ and $] u, v[=\emptyset$ for $u \geqq v$ (see (2.10) on $p .18)$. Let $a, b \in \mathbb{R}$. Then

$$
\begin{align*}
& \left.[a, b]=\bigcap_{n \in \mathbb{N}}\right] a-\frac{1}{n}, b+\frac{1}{n}[.  \tag{5.7}\\
& ] a, b\left[=\bigcup_{n \in \mathbb{N}}\left[a+\frac{1}{n}, b-\frac{1}{n}\right],\right. \tag{5.8}
\end{align*}
$$

## A. PROOF of (5.7)

Case 1: We assume that $a>b$. Then $[a, b]=\emptyset$, hence (5.7) is valid if we can show that there exists $N \in \mathbb{N}$ such that $] a-\frac{1}{N}, b+\frac{1}{N}\left[\right.$ is empty. We do this as follows. Let $\varepsilon:=\frac{a-b}{2}$. Then $\varepsilon>0$. According to (5.6) there exists $N \in \mathbb{N}$ such that $\frac{1}{n}<\varepsilon$ for all $n \geqq N$; in particular, $\frac{1}{N}<\varepsilon=\frac{a-b}{2}$. Of course the choice of $N$ depends on $x$. It follows that

$$
\left.\frac{2}{N}<a-b, \text { hence } b+\frac{1}{N}<a-\frac{1}{N}, \text { hence }\right] a-\frac{1}{N}, b+\frac{1}{N}[=\emptyset
$$

Case 2: We assume that $a=b$. Then $[a, b]=\{a\}$. Clearly $a \in] a-\frac{1}{n}, b+\frac{1}{n}[$ for all $n$, hence $\{a\} \subseteq$ $\left.\bigcap_{n}\right] a-\frac{1}{n}, b+\frac{1}{n}[$. The proof for case 2 is done if we can show that if $x \neq a$ then there exists $N \in \mathbb{N}$ such that $x \notin] a-\frac{1}{N}, b+\frac{1}{N}[$.

If $x<a$, let $\varepsilon:=a-x>0$. According to (5.6) there exists $N \in \mathbb{N}$ such that $\frac{1}{N}<\varepsilon=a-x$, i.e., $x<a-\frac{1}{N}$, hence $\left.x \notin\right] a-\frac{1}{N}, b+\frac{1}{N}[$.

If $x>a$, we can similarly find some $N \in \mathbb{N}$ such that $\frac{1}{N}\left\langle x-a\right.$, i.e., $x>a+\frac{1}{N}$, hence $\left.x \notin\right] a-\frac{1}{N}, b+\frac{1}{N}[$. Case 3: We assume that $a<b$. If $n \in \mathbb{N}$ then $] a-\frac{1}{n}, b+\frac{1}{n}\left[\supseteq[a, b]\right.$, hence $\left.[a, b] \subseteq \bigcap_{n \in \mathbb{N}}\right] a-\frac{1}{n}, b+\frac{1}{n}[$.

We finally show " $\supseteq$ ". If $\left.x \in \bigcap_{n \in \mathbb{N}]}\right] a-\frac{1}{n}, b+\frac{1}{n}[$ then

$$
\begin{equation*}
a-\frac{1}{n}<x<b+\frac{1}{n} \text { for all } n \in \mathbb{N} . \tag{5.9}
\end{equation*}
$$

The proof is done if we can show that this implies both $x \geqq a$ and $x \leqq b$.
Assume to the contrary that $x<a$. Let $\varepsilon:=a-x$. Then $\varepsilon>0$. According to (5.6) there exists $N \in \mathbb{N}$ such that $\frac{1}{n}<\varepsilon$ for all $n \geqq N$; in particular, $\frac{1}{N}<\varepsilon$. Of course the choice of $N$ depends on $x$.

We obtain from $a-x=\varepsilon>\frac{1}{N}$ that $x<a-\frac{1}{N}$. This contradicts (5.9) and we have proved that $x \geqq a$. Demonstrating that $x \leqq b$ is similar. We have proved " $\supseteq$ ".
B. PROOF of (5.8): This proof is left as exercise 5.1.

Example 5.1. For any set $A$ we have $A=\bigcup_{a \in A}\{a\}$. According to (5.3) this also is true if $A=\emptyset$.
The following trivial lemma (a lemma is a "proof subroutine" which is not remarkable on its own but very useful as a reference for other proofs) is useful if you need to prove statements of the form $A \subseteq B$ or $A=B$ for sets $A$ and $B$. It is a means to simplity the proofs of [1] B/G (Beck/Geoghegan), project 5.12. You must reference this lemma as the "inclusion lemma" when you use it in your homework or exams. Be sure to understand what it means if you choose $J=\{1,2\}$ (draw one or two Venn diagrams).

Lemma 5.1 (Inclusion lemma). Let $J$ be an arbitrary, nonempty index set. Let $U, X_{j}, Y, Z_{j}, W(j \in J)$ be sets such that $U \subseteq X_{j} \subseteq Y \subseteq Z_{j} \subseteq W$ for all $j \in J$. Then

$$
\begin{equation*}
U \subseteq \bigcap_{j \in J} X_{j} \subseteq Y \subseteq \bigcup_{j \in J} Z_{j} \subseteq W \tag{5.10}
\end{equation*}
$$

PROOF: Note that we need at various places in this proof the existence of some $j_{0} \in J$, i.e. the assumption that $J \neq \emptyset$ is essential.
a. Let $x \in U$. Then $x \in X_{j}$ for all $j \in J$, hence $x \in \bigcap_{j \in J} X_{j}$. This proves the first inclusion.
b. Now let $x \in \bigcap_{j \in J} X_{j}$ and $j_{0} \in J$. Then $x \in X_{j}$ for all $j \in J$; in particular, $x \in X_{j_{0}}$. It follows from $X_{j_{0}} \subseteq Y$ that $x \in Y$ and we have shown the second inclusion.
c. Let $x \in Y$ and $j_{0} \in J$. It follows from $Y \subseteq Z_{j_{0}}$ that $x \in Z_{j_{0}}$. But then $x \in\{z: z \in$ $Z_{j}$ for some $\left.j \in J\right\}$, i.e., $x \in \bigcup_{j \in J} Z_{j}$. This proves the third inclusion.
d. Finally, assume $x \in \bigcup_{j \in J} Z_{j}$ It follows from the definitions of unions that there exists $j_{0} \in J$ such that $x \in Z_{j_{0}}$. But then $x \in W$ as $W$ contains $Z_{j_{0}}$. It follows that $\bigcup_{j \in J} Z_{j} \subseteq W$. This finishes the proof of the rightmost inclusion.

Definition 5.2 (Disjoint families). Let $J$ be a nonempty set. We call a family of sets $\left(A_{i}\right)_{i \in J}$ a mutually disjoint family if for any two different indices $i, j \in J$ it is true that $A_{i} \cap A_{j}=\emptyset$, i.e., if any two sets in that family with different indices are mutually disjoint.

Definition 5.3 (Partition). Let $\mathfrak{A} \subseteq 2^{\Omega}$. We call $\mathfrak{A}$ a partition or a partitioning of $\Omega$ if
a. $A \cap B=\emptyset$ for any two $A, B \in \mathfrak{A}$,
b. $\Omega=\biguplus[A: A \in \mathfrak{A}]$.

We extend this definition to arbitrary families and hence finite collections and sequences of subsets of $\Omega$ : Let $J$ be an arbitrary nonempty set, let $\left(A_{j}\right)_{j \in J}$ be a family of subsets of $\Omega$. We call $\left(A_{j}\right)_{j \in J}$ a partition of $\Omega$ if it is a mutually disjoint family which satisfies $\Omega=\biguplus\left[A_{j}: j \in J\right]$, i.e., if $\mathfrak{A}:=\left\{A_{j}: j \in J\right\}$ is a partition of $\Omega$.
Note that duplicate nonempty sets cannot occur in a disjoint family of sets because otherwise the condition of mutual disjointness does not hold.

Example 5.2. Here are some examples of partitions.
a. For any set $\Omega$ the collection $\{\{\omega\}: \omega \in \Omega\}$ is a partition of $\Omega$.
b. The empty set is a partition of the empty set and it is its only partition. Do you see that this is a special case of a?
c. The set of half open intervals $] k, k+1]: k \in \mathbb{Z}\}$ is a partitioning of $\mathbb{R}$.
d. Given is a strictly increasing sequence $n_{0}=0<n_{1}<n_{2}<\ldots$ of non-negative integers. For $k \in \mathbb{N}$ let $A_{k}:=\left\{j \in \mathbb{N}: n_{k-1}<j \leqq n_{k}\right\}$. Then the set $\left\{A_{k}: k \in \mathbb{N}\right\}$ is a partition of $\mathbb{N}\left(\right.$ not of $\left.\mathbb{Z}_{\geq 0}!\right)$

Theorem 5.1 (De Morgan's Law). Let there be a universal set $\Omega$ (see (2.6) on p.15). Then the following "duality principle" holds for any indexed family $\left(A_{\alpha}\right)_{\alpha \in I}$ of sets:

$$
\begin{equation*}
\text { a. }\left(\bigcup_{\alpha} A_{\alpha}\right)^{\complement}=\bigcap_{\alpha} A_{\alpha}^{\complement} \quad \text { b. }\left(\bigcap_{\alpha} A_{\alpha}\right)^{\complement}=\bigcup_{\alpha} A_{\alpha}^{\complement} \tag{5.11}
\end{equation*}
$$

To put this in words, the complement of an arbitrary union is the intersection of the complements and the complement of an arbitrary intersection is the union of the complements.

Generally speaking the formulas are a consequence of the duality principle for set operations which states that any true statement involving a family of subsets of a universal sets can be converted into its "dual" true statement by replacing all subsets with their complements, all unions with intersections and all intersections with unions.

PROOF of De Morgan's law, formula a:

1) First we prove that $\left(\bigcup_{\alpha} A_{\alpha}\right)^{\complement} \subseteq \bigcap_{\alpha} A_{\alpha}^{\complement}$ :

Assume that $x \in\left(\bigcup_{\alpha} A_{\alpha}\right)^{\complement}$. Then $x \notin \bigcup_{\alpha} A_{\alpha}$ which is the same as saying that $x$ does not belong to any of the $A_{\alpha}$. That means that x belongs to each $A_{\alpha}^{\complement}$ and hence also to the intersection $\bigcap_{\alpha} A_{\alpha}^{\complement}$.
2) Now we prove that $\left(\bigcup_{\alpha} A_{\alpha}\right)^{\complement} \supseteq \bigcap_{\alpha} A_{\alpha}^{\complement}$ : Let $x \in \bigcap A_{\alpha}^{\complement}$. Then $x$ belongs to each of the $A_{\alpha}^{\complement}$ and hence to none of the $A_{\alpha}$. Then it also does not belong to the union of all the $A_{\alpha}$ and must therefore belong to the complement $\left(\bigcup_{\alpha} A_{\alpha}\right)^{\complement}$. This completes the proof of formula $a$. The proof of formula $b$ is very similar and given as exercise 5.4 on p.101.

You should draw the Venn diagrams involving just two sets $A_{1}$ and $A_{2}$ for both formulas a and $b$ so that you understand the visual representation of De Morgan's law.

Proposition 5.2 (Distributivity of unions and intersections). Let $\left(A_{i}\right)_{i \in I}$ be an arbitrary family of sets and let $B$ be a set. Then

$$
\begin{align*}
& \bigcup_{i \in I}\left(B \cap A_{i}\right)=B \cap \bigcup_{i \in I} A_{i},  \tag{5.12}\\
& \bigcap_{i \in I}\left(B \cup A_{i}\right)=B \cup \bigcap_{i \in I} A_{i} . \tag{5.13}
\end{align*}
$$

PROOF: We only prove (5.12). The proof of (5.13) is left as exercise 6.3.
PROOF of " $\subseteq$ ": It follows from $B \cap A_{i} \subseteq A_{i}$ for all $i$ that $\bigcup_{i}\left(B \cap A_{i}\right) \subseteq \bigcup_{i} A_{i}$. Moreover, $B \cap A_{i} \subseteq B$ for all i implies $\bigcup_{i}\left(B \cap A_{i}\right) \subseteq \bigcup_{i} B$ which equals $B$. It follows that $\bigcup_{i}\left(B \cap A_{i}\right)$ is contained in the intersection $\left(\bigcup_{i} A_{i}\right) \cap B$.
PROOF of " $\supseteq$ ": Let $x \in B \cap \bigcup_{i} A_{i}$. Then $x \in B$ and $x \in A_{i^{\star}}$ for some $i^{\star} \in I$, hence $x \in B \cap A_{i^{\star}}$, hence $x \in \bigcup_{i}\left(B \cap A_{i}\right)$.

Proposition 5.3 (Rewrite unions as disjoint unions). Let $\left(A_{j}\right)_{j \in \mathbb{N}}$ be a sequence of sets which all are contained within the universal set $\Omega$. For $n \in \mathbb{N}$ let $B_{n}:=\bigcup_{j=1}^{n} A_{j}=A_{1} \cup A_{2} \cup \cdots \cup A_{n}$

Further, let $C_{1}:=A_{1}=B_{1}$ and $C_{n+1}:=A_{n+1} \backslash B_{n}(n \in \mathbb{N})$. Then
a. The sequence $\left(B_{n}\right)_{j}$ is increasing: $m<n \Rightarrow B_{m} \subseteq B_{n}$,
b. For each $n \in \mathbb{N}, \bigcup_{j=1}^{n} A_{j}=\bigcup_{j=1}^{n} B_{j}$,
c. The sets $C_{j}$ are mutually disjoint and $\bigcup_{j=1}^{n} A_{j}=\bigcup_{j=1}^{n} B_{j}=\biguplus_{j=1}^{n} C_{j}$.

PROOF of $\boldsymbol{a}$ and of $\boldsymbol{b}$ : Left as exercise 5.2 (p.101).

PROOF of $c$ : Let $1 \leqq j \leqq n$. We note that $C_{j} \subseteq A_{j} \subseteq B_{j} \subseteq B_{n}$ and obtain

$$
C_{j} \cap C_{n+1} \subseteq B_{n} \cap C_{n+1}=B_{n} \cap\left(A_{n+1} \backslash B_{n}\right)=B_{n} \cap\left(A_{n+1} \cap B_{n}^{\complement}\right)=A_{n+1} \cap\left(B_{n} \cap B_{n}^{\complement}\right)=\emptyset .
$$

We have proved that for any $j, k \in \mathbb{N}$ such that $j<k$ the sets $C_{j}$ and $C_{k}$ have empty intersection (we replaced $n+1$ with $k$ ) and it follows that the entire sequence of sets $C_{j}$ is disjoint.

We finally prove that $\bigcup_{j=1}^{n} A_{j}=\bigcup_{j=1}^{n} B_{j}=\biguplus_{j=1}^{n} C_{j}$. The first equation follows from $\boldsymbol{b}$. To prove the second equation we first show that $\biguplus_{j=1}^{n} C_{j} \subseteq \bigcup_{j=1}^{n} A_{j}$. This is immediate from $C_{n} \subseteq A_{n}$ for all $n \in \mathbb{N}$.

We finally prove that $\bigcup_{j=1}^{n} A_{j} \subseteq \biguplus_{j=1}^{n} C_{j}$. Let $x \in \bigcup_{j=1}^{n} A_{j}$. Then $x \in A_{j}$ for at least one $1 \leqq j \leqq n$. Let $j_{0}$ be the smallest such $j$. If $j_{0}=1$ then $x \in C_{1}$ because $C_{1}=A_{1}$, hence $x \in \biguplus_{j=1}^{n} C_{j}$ and we are done. Otherwise $x \notin A_{j}$ for all $1 \leqq j<j_{0}$, hence $x \notin \bigcup_{j=1}^{j_{0}-1} A_{j}=B_{j_{0}-1}$, hence $x \in A_{j_{0}} \backslash B_{j_{0}-1}$, i.e., $x \in C_{j_{0}}$. It follows that $x \in \biguplus_{j=1}^{n} C_{j}$.

### 5.2 Cartesian Products of more than Two Sets

Remark 5.3 (Associativity of cartesian products). Assume we have three sets $A, B$ and $C$. We can then look at

$$
\begin{aligned}
& (A \times B) \times C=\{((a, b), c): a \in A, b \in B, c \in C\} \\
& A \times(B \times C)=\{(a,(b, c)): a \in A, b \in B, c \in C\}
\end{aligned}
$$

The mapping

$$
F:(A \times B) \times C \rightarrow A \times(B \times C), \quad((a, b), c) \mapsto(a,(b, c))
$$

is bijective because it has the mapping

$$
G: A \times(B \times C) \rightarrow(A \times B) \times C, \quad(a,(b, c)) \mapsto((a, b), c)
$$

as an inverse. For both $(A \times B) \times C$ and $A \times(B \times C)$ there are bijections to the set $\{(a, b, c): a \in$ $A, b \in B, c \in C\}$ of all triplets $(a, b, c)$ : the obvious bijections would be $(a, b, c) \mapsto((a, b), c)$ and $(a, b, c) \mapsto((a, b), c)$.

This remark leads us to the following definition:
Definition 5.4 (Cartesian Product of three or more sets). The cartesian product of three sets $A, B$ and $C$ is defined as

$$
A \times B \times C:=\{(a, b, c): a \in A, b \in B, c \in C\}
$$

i.e., it consists of all pairs $(a, b, c)$ with $a \in A, b \in B$ and $c \in C$.

More generally, for $N$ sets $X_{1}, X_{2}, X_{3}, \ldots, X_{N}(N \in \mathbb{N})$, we define the cartesian product as ${ }^{53}$

$$
X_{1} \times X_{2} \times X_{3} \times \ldots \times X_{N}:=\left\{\left(x_{1}, x_{2}, \ldots, x_{N}\right): x_{j} \in X_{j} \text { for all } 1 \leqq j \leqq N\right\}
$$

[^39]Note that the elements of this set are finite lists in the sense of def.4.20 (Vectors and finite lists) on p. 93 .

Two elements $\left(x_{1}, x_{2}, \ldots, x_{N}\right)$ and $\left(y_{1}, y_{2}, \ldots, y_{N}\right)$ of $X_{1} \times X_{2} \times X_{3} \times \ldots \times X_{N}$ are called equal if and only if $x_{j}=y_{j}$ for all $j$ such that $1 \leqq j \leqq N$. In this case we write $\left(x_{1}, x_{2}, \ldots, x_{N}\right)=\left(y_{1}, y_{2}, \ldots, y_{N}\right)$.

As a shorthand, we abbreviate $X^{N}:=\underbrace{X \times X \times+\cdots \times X}_{N \text { times }}$.
Example 5.3 ( $N$-dimensional coordinates). Here is the most important example of a cartesian product of $N$ sets. Let $X_{1}=X_{2}=\ldots=X_{N}=\mathbb{R}$. Then

$$
\mathbb{R}^{N}=\left\{\left(x_{1}, x_{2}, \ldots, x_{N}\right): x_{j} \in \mathbb{R} \text { for } 1 \leqq j \leqq N\right\}
$$

is the set of points in $N$-dimensional space. You may not be familiar with what those are unless $N=2$ (see example 4.1 above) or $N=3$.

In the 3 -dimensional case it is customary to write $(x, y, z)$ rather than $\left(x_{1}, x_{2}, x_{3}\right)$. Each such triplet of real numbers represents a point in (ordinary 3-dimensional) space and we speak of its $x$-coordinate, $y$-coordinate and $z$-coordinate.

For the sake of completeness: If $N=1$, the item $(x) \in \mathbb{R}^{1}$ (where $x \in \mathbb{R}$; observe the parentheses around $x$ ) is considered the same as the real number $x$. In other words, we "identify" $\mathbb{R}^{1}$ with $\mathbb{R}$. Such a "one-dimensional point" is simply a point on the $x$-axis.

A short note on vectors and coordinates: For $N \leqq 3$ you can visualize the following: Given a point $x$ on the $x$-axis or in the plane or in 3-dimensional space, there is a unique arrow that starts at the point whose coordinates are all zero (the origin) and ends at the location marked by the point $x$. Such an arrow is customarily called a vector.

Because it makes sense in dimensions $1,2,3$, an $N$-tuple ( $x_{1}, x_{2}, \ldots, x_{N}$ ) of numbers is called a vector of dimension $N .{ }^{54}$ You will read more about this in ch. 9 about vectors and vector spaces on page 155.

This is worth while repeating: We can uniquely identify each $x \in \mathbb{R}^{N}$ with the corresponding vector: an arrow that starts in $(\underbrace{0,0, \ldots, 0}_{N \text { times }})$ and ends in $x$.

More will be said about $n$-dimensional space in section 9, p. 155 on vectors and vector spaces.

Example 5.4 (Parallelepipeds). Let $a_{1}<b_{1}, a_{2}<b_{2}, a_{3}<b_{3}$ be real numbers. Then

$$
\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \times\left[a_{3}, b_{3}\right]=\left\{(x, y, z): a_{1} \leqq x \leqq b_{1}, a_{2} \leqq y \leqq b_{2}, a_{3} \leqq z \leqq b_{3}\right\}
$$

is the parallelepiped (box or quad parallel to the coordinate axes) with sides $\left[a_{1}, b_{1}\right],\left[a_{2}, b_{2}\right]$ and $\left[a_{3}, b_{3}\right]$. This generalizes in an obvious manner to $N$ dimensions:

Let $N \in \mathbb{N}$ and $a_{j}<b_{j}\left(j \in \mathbb{N}, j \leqq N, a_{j}, b_{j} \in \mathbb{R}\right)$. Then

$$
\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \times \cdots \times\left[a_{N}, b_{N}\right]=\left\{\left(x_{1}, x_{2}, \ldots, x_{N}\right): a_{j} \leqq x_{j} \leqq b_{j}, j \in \mathbb{N}, j \leqq N\right\}
$$

is the parallelepiped with sides $\left[a_{1}, b_{1}\right], \ldots,\left[a_{N}, b_{N}\right]$.

[^40]We now introduce cartesian products of an entire family of sets $\left(X_{i}\right)_{i \in I}$.
Definition 5.5 (Cartesian Product of a family of sets). Let $I$ be an arbitrary, nonempty set (the index set). Let $\left(X_{i}\right)_{i \in I}$ be a family of nonempty sets $X_{i}$. The cartesian product of the family $\left(X_{i}\right)_{i \in I}$ is the set

$$
\prod_{i \in I} X_{i}:=\left(\prod X_{i}\right)_{i \in I}:=\left\{\left(x_{i}\right)_{i \in I}: x_{k} \in X_{k} \forall k \in I\right\}
$$

of all familes $\left(x_{i}\right)_{i \in I}$ each of whose members $x_{j}$ belongs to the corresponding set $X_{j}$. The " $\Pi$ " is the greek "upper case" letter "Pi" (whose lower case incarnation " $\pi$ " you are probably more familiar with).

Two elements $\left(x_{i}\right)_{i \in I}$ and $\left(y_{k}\right)_{k \in I}$ of $\prod_{i \in I} X_{i}$ are called equal if and only if $x_{j}=y_{j}$ for all $j \in I$. In this case we write $\left(x_{i}\right)_{i \in I}=\left(y_{k}\right)_{k \in I} .{ }^{55}$

If all sets $X_{i}$ are equal to one and the same set $X$, we also write $X^{I}:=\prod_{i \in I} X:=\prod_{i \in I} X_{i}$.
Remark 5.4. We note that each element $\left(y_{x}\right)_{x \in X}$ of the cartesian product $Y^{X}$ is the function

$$
y(\cdot): X \rightarrow Y, \quad x \mapsto y_{x}
$$

(see def.4.17 (indexed families) and the subsequent remarks concerining the equivalence of functions and families). In other words,

$$
\begin{equation*}
Y^{X}=\{f: f \text { is a function with domain } X \text { and codomain } Y\} . \tag{5.14}
\end{equation*}
$$

### 5.3 Exercises for Ch. 5

Exercise 5.1. Prove (5.8) of prop.5.1 on p.96: Let $a, b \in \mathbb{R}$. Then $] a, b\left[=\bigcup_{n \in \mathbb{N}}[a+1 / n, b-1 / n]\right.$. Adapt the proof of (5.7) but note that this one is simpler. There are only two cases to worry about: $a \geqq b$ (very easy!) vs $a<b$.

Exercise 5.2. Prove $\mathbf{a}$ and $\mathbf{b}$ of prop. 5.3 (Rewrite unions as disjoint unions) on p.98:
Let $\left(A_{j}\right)_{j \in \mathbb{N}}$ such that $A_{j} \subseteq \Omega$ for all $j \in \mathbb{N}$. For $n \in \mathbb{N}$ let $B_{n}:=\bigcup_{j=1}^{n} A_{j}=A_{1} \cup A_{2} \cup \cdots \cup A_{n}$
Further, let $C_{1}:=A_{1}=B_{1}$ and $C_{n+1}:=A_{n+1} \backslash B_{n}(n \in \mathbb{N})$. Then
a. The sequence ${ }_{n}^{\left(B_{j}\right)_{j}}$ is $\underset{n}{\text { increasing: }} m<n \Rightarrow B_{m} \subseteq B_{n}$,
b. For each $n \in \mathbb{N}, \bigcup_{j=1}^{n} A_{j}=\bigcup_{j=1}^{n} B_{j}$.

Exercise 5.3. (See example 4.22 on p .87 ). Let $X:=[0,2]$. For $0 \leqq x \leqq 2$ let $A_{x}:=[x, 2 x]$.
a. What is $\bigcap\left[A_{x}: x \in X\right]$ ?
b. What is $\bigcup\left[A_{x}: x \in X\right]$ ?

Exercise 5.4. Prove $\mathbf{b}$ of thm.5.1 (De Morgan's Law):
Let $\left(A_{\alpha}\right)_{\alpha \in I}$ be a family of subsets of a universal set $\Omega$. Then $\left(\bigcap_{\alpha} A_{\alpha}\right)^{\complement}=\bigcup_{\alpha} A_{\alpha}^{\complement}$.

[^41]Exercise 5.5. Prove the second formula of prop. 5.2 (Distributivity of unions and intersections): Let $\left(A_{i}\right)_{i \in I}$ be an arbitrary family of sets and let $B$ be a set. Then

$$
\bigcap_{i \in I}\left(B \cup A_{i}\right)=B \cup \bigcap_{i \in I} A_{i} .
$$

## 6 Sets and Functions, Direct and Indirect Images (Study this!)

### 6.1 Direct Images and Indirect Images (Preimages) of a Function

Definition 6.1. Let $X, Y$ be two nonempty sets and $f: X \rightarrow Y$ be an arbitrary function with domain $X$ and codomain $Y$. Let $A \subseteq X$ and $B \subseteq Y$. Let

$$
\begin{align*}
f(A) & :=\{f(x): x \in A\},  \tag{6.1}\\
f^{-1}(B) & :=\{x \in X: f(x) \in B\} . \tag{6.2}
\end{align*}
$$

We call $f(A)$ the direct image ${ }^{56}$ of $A$ under $f$ and we call We call $f^{-1}(B)$ the indirect image or preimage of $B$ under $f$.

## Notational conveniences I:

$$
\text { If we have a set that is written as }\{\ldots\} \text { then we may write } f\{\ldots\} \text { instead of } f(\{\ldots\}) \text { and } f^{-1}\{\ldots\}
$$ instead of $f^{-1}(\{\ldots\})$. Specifically for singletons $\{x\} \subseteq X$ and $\{y\} \subseteq Y$ we obtain $f\{x\}$ and $f^{-1}\{y\}$.

Many mathematicians will write $f^{-1}(y)$ instead of $f^{-1}\{y\}$ but this author sees no advantages doing so whatsover. There seemingly are no savings with respect to time or space for writing that alternate form but we are confounding two entirely separate items: a subset $f^{-1}\{y\}$ of $X$ v.s. the function value $f^{-1}(y)$ of $y \in Y$ which is an element of $X$. We are allowed to talk about the latter only in case that the inverse function $f^{-1}$ of $f$ exists.

## Notational conveniences II:

In measure theory and probability theory the following notation is also very common:
$\{f \in B\}:=f^{-1}(B),\{f=y\}:=f^{-1}\{y\}$.
Let $a<b \in \mathbb{R}$. We write $\{a \leqq f \leqq b\}:=f^{-1}([a, b]),\{a<f<b\}:=f^{-1}(] a, b[)$,
$\left.\left.\{a \leqq f<b\}:=f^{-1}\left(\left[a, b[),\{a<f \leqq b\}:=f^{-1}(] a, b\right]\right),\{f \leqq b\}:=f^{-1}(]-\infty, b\right]\right)$, etc.
Example 6.1 (Direct images). Let $f: \mathbb{R} \rightarrow \mathbb{R} ; \quad f(x)=x^{2}$.
a. $\left.f(]-4,-2[)=\left\{x^{2}: x \in\right]-4,-2[ \}=\left\{x^{2}:-4<x<-2\right\}=\right] 4,16[$.
b. $f([1,2])=\left\{x^{2}: x \in[1,2]\right\}=\left\{x^{2}: 1 \leqq x \leqq 2\right\}=[1,4]$.
c. $\quad f([5,6])=\left\{x^{2}: x \in[5,6]\right\}=\left\{x^{2}: 5 \leqq x \leqq 6\right\}=[25,36]$.
d. $f(]-4,-2[\cup[1,2] \cup[5,6])=\left\{x^{2}: x \in\right]-4,-2[$ or $x \in[1,2]$ or $x \in[5,6]\}$ $=] 4,16[\cup[1,4] \cup[25,36]=[1,16[\cup[25,36]$.

Example 6.2 (Direct images). Let $f: \mathbb{R} \rightarrow \mathbb{R} ; \quad f(x)=x^{2}$.
a. $\left.f(]-4,2[)=\left\{x^{2}: x \in\right]-4,2[ \}=\left\{x^{2}:-4<x<2\right\}=\right] 4,16[$.
b. $f([1,3])=\left\{x^{2}: x \in[1,3]\right\}=\left\{x^{2}: 1 \leqq x \leqq 3\right\}=[1,9]$.
c. $\quad f(]-4,2[\cap[1,3])=\left\{x^{2}: x \in\right]-4,2[$ and $x \in[1,3]\}=\left\{x^{2}: 1 \leqq x<2\right\}=[1,4[$.

And here are the results for the preimages of the same sets with respect to the same function $x \mapsto x^{2}$.
Example 6.3 (Preimages). Let $f: \mathbb{R} \rightarrow \mathbb{R} ; \quad f(x)=x^{2}$.

[^42]a. $\quad f^{-1}(]-4,-2[)=\left\{x \in \mathbb{R}: x^{2} \in\right]-4,-2[ \}=\{-4<f<-2\}=\emptyset$.
b. $f^{-1}([1,2])=\left\{x \in \mathbb{R}: x^{2} \in[1,2]\right\}=\{1 \leqq f \leqq 2\}=[-\sqrt{2},-1] \cup[1, \sqrt{2}]$.
c. $f^{-1}([5,6])=\left\{x \in \mathbb{R}: x^{2} \in[5,6]\right\}=\{5 \leqq f \leqq 6\}=[-\sqrt{6},-\sqrt{5}] \cup[\sqrt{5}, \sqrt{6}]$.
d. $f^{-1}(]-4,-2[\cup[1,2] \cup[5,6])=\left\{x \in \mathbb{R}: x^{2} \in\right]-4,-2\left[\right.$ or $x^{2} \in[1,2]$ or $\left.x^{2} \in[5,6]\right\}$
$=[-\sqrt{2},-1] \cup[1, \sqrt{2}] \cup[-\sqrt{6},-\sqrt{5}] \cup[\sqrt{5}, \sqrt{6}]$.
Example 6.4 (Preimages). Let $f: \mathbb{R} \rightarrow \mathbb{R} ; \quad f(x)=x^{2}$.
a. $\left.f^{-1}(]-4,2[)=\left\{x \in \mathbb{R}: x^{2} \in\right]-4,2[ \}=\left\{x \in \mathbb{R}:-4<x^{2}<2\right\}=\right]-2,2[$.
b. $\quad f^{-1}([1,3])=\left\{x \in \mathbb{R}: x^{2} \in[1,3]\right\}=\left\{x \in \mathbb{R}: 1 \leqq x^{2} \leqq 3\right\}=[-\sqrt{3}, 1] \cup[1, \sqrt{3}]$.
c. $f^{-1}(]-4,2[\cap[1,3])=\left\{x \in \mathbb{R}: x^{2} \in\right]-4,2\left[\right.$ and $\left.x^{2} \in[1,3]\right\}$
$\left.\left.=\left\{x \in \mathbb{R}: 1 \leqq x^{2}<2\right\}=\right]-\sqrt{2},-1\right] \cup[1, \sqrt{2}[$.

Unless stated otherwise, $X, Y$ and $f$ are as defined above for the remainder of this chapter: $f: X \rightarrow$ $Y$ is a function with domain $X$ and codomain $Y$.

Proposition 6.1. Some simple properties:

$$
\begin{align*}
f(\emptyset) & =f^{-1}(\emptyset)=\emptyset  \tag{6.3}\\
A_{1} & \subseteq A_{2} \subseteq X \Rightarrow f\left(A_{1}\right) \subseteq f\left(A_{2}\right) \quad \text { (monotonicity of } f\{\ldots\} \text { ) }  \tag{6.4}\\
B_{1} & \left.\subseteq B_{2} \subseteq Y \Rightarrow f^{-1}\left(B_{1}\right) \subseteq f^{-1}\left(B_{2}\right) \quad \text { (monotonicity of } f^{-1}\{\ldots\}\right)  \tag{6.5}\\
x & \in X \Rightarrow f(\{x\})=\{f(x)\}  \tag{6.6}\\
f(X) & =Y \Leftrightarrow f \text { is surjective }  \tag{6.7}\\
f^{-1}(Y) & =X \quad \text { always! } \tag{6.8}
\end{align*}
$$

PROOF: Left as an exercise.
Proposition 6.2 ( $f^{-1}$ is compatible with all basic set ops). In the following we assume that $J$ is an arbitrary index set, and that $B \subseteq Y, B_{j} \subseteq Y$ for all $j$. Then

$$
\begin{align*}
f^{-1}\left(\bigcap_{j \in J} B_{j}\right) & =\bigcap_{j \in J} f^{-1}\left(B_{j}\right)  \tag{6.9}\\
f^{-1}\left(\bigcup_{j \in J} B_{j}\right) & =\bigcup_{j \in J} f^{-1}\left(B_{j}\right)  \tag{6.10}\\
f^{-1}\left(B^{\complement}\right) & =f^{-1}(B)^{\complement}  \tag{6.11}\\
f^{-1}\left(B_{1} \backslash B_{2}\right) & =f^{-1}\left(B_{1}\right) \backslash f^{-1}\left(B_{2}\right)  \tag{6.12}\\
f^{-1}\left(B_{1} \Delta B_{2}\right) & =f^{-1}\left(B_{1}\right) \Delta f^{-1}\left(B_{2}\right) \tag{6.13}
\end{align*}
$$

PROOF of (6.9): Let $x \in X$. Then

$$
\begin{align*}
x \in f^{-1}\left(\bigcap_{j \in J} B_{j}\right) & \Leftrightarrow f(x) \in \bigcap_{j \in J} B_{j} \quad \text { (def preimage) } \\
& \Leftrightarrow \forall j f(x) \in B_{j} \quad(\text { def } \cap) \\
& \Leftrightarrow \forall j x \in f^{-1}\left(B_{j}\right) \quad \text { (def preimage) }  \tag{6.14}\\
& \Leftrightarrow x \in \bigcap_{j \in J} f^{-1}\left(B_{j}\right) \quad(\text { def } \cap)
\end{align*}
$$

PROOF of (6.10): Let $x \in X$. Then

$$
\begin{align*}
x \in f^{-1}\left(\bigcup_{j \in J} B_{j}\right) & \Leftrightarrow f(x) \in \bigcup_{j \in J} B_{j} \quad(\text { def preimage }) \\
& \Leftrightarrow \exists j_{0}: f(x) \in B_{j_{0}} \quad(\text { def } \cup)  \tag{6.15}\\
& \left.\Leftrightarrow \exists j_{0}: x \in f^{-1}\left(B_{j_{0}}\right) \quad \text { (def preimage }\right) \\
& \Leftrightarrow x \in \bigcup_{j \in J} f^{-1}\left(B_{j}\right) \quad(\text { def } \cup)
\end{align*}
$$

PROOF of (6.11): Let $x \in X$. Then

$$
\begin{align*}
x \in f^{-1}\left(B^{\complement}\right) & \Leftrightarrow f(x) \in B^{\complement} \quad(\text { def preimage }) \\
& \Leftrightarrow f(x) \notin B \quad(\text { def }(\cdot) \complement)  \tag{6.16}\\
& \Leftrightarrow x \notin f^{-1}(B) \quad(\text { def preimage }) \\
& \left.\Leftrightarrow x \in f^{-1}(B)^{\complement} \quad(\cdot) \complement\right)
\end{align*}
$$

PROOF of (6.12): Let $x \in X$. Then

$$
\begin{align*}
x \in f^{-1}\left(B_{1} \backslash B_{2}\right) & \Leftrightarrow x \in f^{-1}\left(B_{1} \cap B_{2}^{\complement}\right) \quad(\operatorname{def} \backslash) \\
& \Leftrightarrow x \in f^{-1}\left(B_{1}\right) \cap f^{-1}\left(B_{2}^{\complement}\right) \quad(\text { see (6.9) }  \tag{6.17}\\
& \Leftrightarrow x \in f^{-1}\left(B_{1}\right) \cap f^{-1}\left(B_{2}\right)^{\complement} \quad(\text { see (6.11) } \\
& \Leftrightarrow x \in f^{-1}\left(B_{1}\right) \backslash f^{-1}\left(B_{2}\right) \quad(\operatorname{def} \backslash)
\end{align*}
$$

PROOF of (6.13): This follows from $B_{1} \Delta B_{2}=\left(B_{1} \backslash B_{2}\right) \cup\left(B_{2} \backslash B_{1}\right)$ and (6.10) and (6.12).
Proposition 6.3 (Properties of the direct image). In the following we assume that $J$ is an arbitrary index set, and that $A \subseteq X, A_{j} \subseteq X$ for all $j$. Then

$$
\begin{align*}
f\left(\bigcap_{j \in J} A_{j}\right) & \subseteq \bigcap_{j \in J} f\left(A_{j}\right)  \tag{6.18}\\
f\left(\bigcup_{j \in J} A_{j}\right) & =\bigcup_{j \in J} f\left(A_{j}\right) \tag{6.19}
\end{align*}
$$

PROOF of (6.18): This follows from the monotonicity of the direct image (see 6.4):

$$
\begin{aligned}
\bigcap_{j \in J} A_{j} \subseteq A_{i} \forall i \in J & \Rightarrow f\left(\bigcap_{j \in J} A_{j}\right) \subseteq f\left(A_{i}\right) \forall i \in J \\
& \Rightarrow f\left(\bigcap_{j \in J} A_{j}\right) \subseteq \bigcap_{i \in J} f\left(A_{i}\right) \quad(\operatorname{def} \cap)
\end{aligned}
$$

First proof of (6.19)) - "Expert proof":

$$
\begin{align*}
y \in f\left(\bigcup_{j \in J} A_{j}\right) & \Leftrightarrow \exists x \in X: f(x)=y \text { and } x \in \bigcup_{j \in J} A_{j} \quad(\operatorname{def} f(A))  \tag{6.20}\\
& \Leftrightarrow \exists x \in X \text { and } j_{0} \in J: f(x)=y \text { and } x \in A_{j_{0}} \quad(\operatorname{def} \cup)  \tag{6.21}\\
& \Leftrightarrow \exists x \in X \text { and } j_{0} \in J: f(x)=y \text { and } f(x) \in f\left(A_{j_{0}}\right) \quad(\operatorname{def} f(A))  \tag{6.22}\\
& \Leftrightarrow \exists j_{0} \in J: y \in f\left(A_{j_{0}}\right) \quad(\operatorname{def} f(A))  \tag{6.23}\\
& \Leftrightarrow y \in \bigcup_{j \in J} f\left(A_{j}\right) \quad(\operatorname{def} \cup) \tag{6.24}
\end{align*}
$$

Alternate proof of (6.19) - Proving each inclusion separately. Unless you have a lot of practice reading and writing proofs whose subject is the equality of two sets you should write your proof the following way:
A. Proof of " $\subseteq$ ":

$$
\begin{align*}
y \in f\left(\bigcup_{j \in J} A_{j}\right) & \Rightarrow \exists x \in X: f(x)=y \text { and } x \in \bigcup_{j \in J} A_{j} \quad(\operatorname{def} f(A))  \tag{6.25}\\
& \Rightarrow \exists j_{0} \in J: f(x)=y \text { and } x \in A_{j_{0}} \quad(\operatorname{def} \cup)  \tag{6.26}\\
& \Rightarrow y=f(x) \in f\left(A_{j_{0}}\right)(\operatorname{def} f(A))  \tag{6.27}\\
& \Rightarrow y \in \bigcup_{j \in J} f\left(A_{j}\right) \quad(\operatorname{def} \cup) \tag{6.28}
\end{align*}
$$

B. Proof of " $\supseteq$ ":

This follows from the monotonicity of $A \mapsto f(A)$ (see 6.4):

$$
\begin{align*}
A_{i} \subseteq \bigcup_{j \in J} A_{j} \forall i \in J & \Rightarrow f\left(A_{i}\right) \subseteq f\left(\bigcup_{j \in J} A_{j}\right) \forall i \in J  \tag{6.29}\\
& \Rightarrow \bigcup_{i \in J} f\left(A_{i}\right) \subseteq f\left(\bigcup_{j \in J} A_{j}\right) \forall i \in J \quad(\operatorname{def} \cup) \tag{6.30}
\end{align*}
$$

The "elementary" proof is barely longer than the first one, but it is so much easier to understand!
Remark 6.1. In general you will not have equality in (6.18). Counterexample: $f(x)=x^{2}$ with domain $\mathbb{R}$ : Let $\left.\left.A_{1}:=\right]-\infty, 0\right]$ and $A_{2} ;=\left[0, \infty\left[\right.\right.$. Then $A_{1} \cap A_{2}=\{0\}$, hence $f\left(A_{1} \cap A_{2}\right)=$ $f(\{0\})=\{0\}$. On the other hand, $f\left(A_{1}\right)=f\left(A_{2}\right)=[0, \infty]$, hence $f\left(A_{1}\right) \cap f\left(A_{2}\right)=[0, \infty]$. Clearly, $\{0\} \subsetneq[0, \infty]$.

Proposition 6.4 (Direct images and preimages of function composition). Let $X, Y, Z$ be an arbitrary, nonempty sets. Let $f: X \rightarrow Y$ and $g: Y \rightarrow Z$ and let $W \subseteq Z$. Then

$$
\begin{align*}
(g \circ f)(U) & =g(f(U)) \text { for all } U \subseteq X  \tag{6.31}\\
(g \circ f)^{-1} & =f^{-1} \circ g^{-1}, \text { i.e., } \quad(g \circ f)^{-1}(W)=f^{-1}\left(g^{-1}(W)\right) \text { for all } W \subseteq Z . \tag{6.32}
\end{align*}
$$

PROOF of (6.31): Left as exercise ??.
PROOF of (6.32):
a. " $\subseteq$ ": Let $W \subseteq Z$ and $x \in(g \circ f)^{-1}(W)$. Then $g(f(x))=(g \circ f)(x) \in W$, hence $f(x) \in g^{-1}(W)$. But then $x \in f^{-1}\left(g^{-1}(W)\right)$. This proves " $\subseteq$ ".
b. " $\supseteq$ ": Let $W \subseteq Z, h:=g \circ f$, and $x \in f^{-1}\left(g^{-1}(W)\right)$. Then $f(x) \in g^{-1}(W)$, hence $h(x)=$ $g(f(x)) \in W$, hence $x \in h^{-1}(W)=(g \circ f)^{-1}(W)$. This proves " $\supseteq$ ".

Proposition 6.5 (Indirect image and fibers of $f$ ). Let $X, Y$ be nonempty sets and let $f: X \rightarrow Y$ be a function. We define on the domain $X$ a relation " $\sim$ " as follows:

$$
\begin{equation*}
x_{1} \sim x_{2} \Leftrightarrow f\left(x_{1}\right)=f\left(x_{2}\right) . \tag{6.33}
\end{equation*}
$$

a. " $\sim$ " is an equivalence relation. Its equivalence classes, which we denote by $[x]_{f},{ }^{57}$ are

$$
\begin{equation*}
[x]_{f}=\{a \in X: f(a)=f(x)\}=f^{-1}\{f(x)\} . \quad(x \in X) \tag{6.34}
\end{equation*}
$$

b. If $A \subseteq X$ then

$$
\begin{equation*}
f^{-1}(f(A))=\bigcup_{a \in A}[a]_{f} \tag{6.35}
\end{equation*}
$$

The proof that " $\sim$ " is an equivalence relation is left as exercise 6.2.
PROOF of (6.34): The equation on the left is nothing but the definition of the equivalence classes generated by an equivalence relation. The equation on the right follows from $f(a)=f(x) \Leftrightarrow a \in f^{-1}\{f(x)\}$, which is true according to the definition of preimages.
PROOF of (6.35):
As $f(A)=f\left(\bigcup_{a \in A}\{x\}\right)=\bigcup_{a \in A}\{f(x)\}$ (see 6.19), it follows that

$$
\begin{align*}
f^{-1}(f(A)) & =f^{-1}\left(\bigcup_{a \in A}\{f(a)\}\right)  \tag{6.36}\\
& =\bigcup_{a \in A} f^{-1}\{f(a)\} \quad(\text { see 6.10) }  \tag{6.37}\\
& =\bigcup_{a \in A}[a]_{f} \quad \text { (see 6.34) } \tag{6.38}
\end{align*}
$$

## Corollary 6.1.

$$
\begin{equation*}
\text { If } A \subseteq X \text { then } f^{-1}(f(A)) \supseteq A \text {. } \tag{6.39}
\end{equation*}
$$

PROOF: It follows from $x \sim x$ for all $x \in X$ that $x \in[x]_{f}$, i.e., $\{x\} \subseteq[x]_{f}$ for all $x \in X$. But then

$$
\begin{equation*}
A=\bigcup_{x \in A}\{x\} \subseteq \bigcup_{x \in A}[x]_{f}=f^{-1}(f(A)) \tag{6.40}
\end{equation*}
$$

where the last equation holds because of (6.35).

## Proposition 6.6.

(6.41)

$$
\text { If } B \subseteq Y \text { then } f\left(f^{-1}(B)\right)=B \cap f(X) \text {. }
$$

[^43]PROOF of " $\subseteq$ ":
Let $y \in f\left(f^{-1}(B)\right)$. There exists $x_{0} \in f^{-1}(B)$ such that $f\left(x_{0}\right)=y$ (def direct image). We have
a. $\quad x_{0} \in f^{-1}(B) \Rightarrow f\left(x_{0}\right) \in B$ (def. of preimage)
b. Of course $x_{0} \in X$. Hence $f\left(x_{0}\right) \in f(X)$.
$\boldsymbol{a}$ and $\boldsymbol{b}$ together imply that $y=f\left(x_{0}\right) \in B \cap f(X)$.
PROOF of " $\supseteq$ ":
Let $y \in f(X) \cap B$. We must prove that $y \in f\left(f^{-1}(B)\right)$.
It follows from $y \in f(X)$ that there exists $x_{0} \in X$ such that $y=f\left(x_{0}\right)$.
Because $f\left(x_{0}\right) \in B$, we conclude that $x_{0} \in f^{-1}(B)$ (def preimage), hence $y=f\left(x_{0}\right) \in f\left(f^{-1}(B)\right)$.
We have shown that if $y \in f(X)$ and $y \in B$ then $y \in f\left(f^{-1}(B)\right)$. The proof of " $\supseteq$ " is completed.
Remark 6.2. Be sure to understand how the assumption $y \in f(X)$ was used.

## Corollary 6.2.

$$
\begin{equation*}
\text { If } B \subseteq Y \text { then } f\left(f^{-1}(B)\right) \subseteq B \tag{6.42}
\end{equation*}
$$

Trivial as $f\left(f^{-1}(B)\right)=B \cap f(X) \subseteq B$.

## Proposition 6.7.

a. Let $A \subseteq X$. If $f: X \rightarrow Y$ is injective then $f^{-1}(f(A))=A$.
b. Let $B \subseteq Y$. If $f: X \rightarrow Y$ is surjective then $f\left(f^{-1}(B)\right)=B$.
c. Let $A \subseteq X$ and $B \subseteq Y$. If $f: X \rightarrow Y$ is injective and if $B=f(A)$ then $f^{-1}(B)=A$.
d. Let $A \subseteq X$ and $B \subseteq Y$. If $f: X \rightarrow Y$ is surective and if $f^{-1}(B)=A$ then $B=f(A)$.
e. Let $A \subseteq X$ and $B \subseteq Y$. If $f: X \rightarrow Y$ is bijective then $B=f(A) \Leftrightarrow f^{-1}(B)=A$.

PROOF: Left as exercise 6.4.
Proposition 6.8. Let $J$ be an arbitrary nonempty index set and let $A \subseteq X, A_{j} \subseteq X$ for all $j$.
Let $f: X \rightarrow Y$ be bijective. Then the following all are true:

$$
\begin{align*}
f\left(\bigcap_{j \in J} A_{j}\right) & =\bigcap_{j \in J} f\left(A_{j}\right)  \tag{6.43}\\
f\left(\bigcup_{j \in J} A_{j}\right) & =\bigcup_{j \in J} f\left(A_{j}\right)  \tag{6.44}\\
f\left(A^{\complement}\right) & =f(A)^{\complement}  \tag{6.45}\\
f\left(A_{1} \backslash A_{2}\right) & =f\left(A_{1}\right) \backslash f\left(A_{2}\right)  \tag{6.46}\\
f\left(A_{1} \Delta A_{2}\right) & =f\left(A_{1}\right) \Delta f\left(A_{2}\right) \tag{6.47}
\end{align*}
$$

PROOF: Left as exercise 6.5.
Proposition 6.9. Let $f: X \rightarrow Y$ be bijective. Let $J$ be an arbitrary nonempty index set and let $\left(A_{j}\right)_{j \in J}$ be a partition of $X$, i.e., if $i \neq j$ then $A_{i} \cap A_{j}=\emptyset$ and $X=\biguplus_{j} A_{j}$. Assume further that none of the $A_{j}$ are enpty. For $j \in J$ let $B_{j}:=f\left(A_{j}\right)$. Then
a. $\quad\left(B_{j}\right)_{j \in J}$ is a partition of $Y$.
b. For $j \in J$ we look at the restriction $\left.f\right|_{A_{j}}: A_{j} \rightarrow Y$ to $A_{j}$. Then $\left.f\right|_{A_{j}}\left(A_{j}\right)=B_{j}$ and the function

$$
f_{j}: A_{j} \rightarrow B_{j}, \quad x \mapsto f_{j}(x):=\left.f\right|_{A_{j}}(x)=f(x)
$$

is a bijection.
PROOF: Left as exercise 6.6.
Corollary 6.3. Let $f: X \rightarrow Y$ be bijective. Let $A \subset X, A \neq \emptyset$ (strict inclusion, so $A^{\complement} \neq \emptyset$ ). Then both

$$
f_{A}: A \rightarrow f(A), \quad x \rightarrow f(x) \quad \text { and } \quad f_{A^{\complement}}: A^{\complement} \rightarrow f\left(A^{\complement}\right), \quad x \rightarrow f(x)
$$

are bijections.
PROOF: This follows from prop.6.9, applied to $J=\{1,2\}, A_{1}=A, A_{2}=A^{\complement}$.
Corollary 6.4. Let $f: X \rightarrow Y$ be bijective. Let $a \subset X$ and assume that $X \neq\{a\}$. Then

$$
\tilde{f}: X \backslash\{a\} \rightarrow Y \backslash\{f(a)\}, \quad x \rightarrow f(x)
$$

also is bijective. ${ }^{58}$
PROOF: This follows from 6.4 applied to $A=\{a\}$ and the fact that $f(\{a\})=\{f(a)\}$.
The following two propositions allow you to replace bijective and surjective functions with more suitable ones that inherit bijectivity or surjectivity. This will come in handy in ch.7.1 on p. 115 when we prove propositions concerning cardinality.

The first proposition shows how to preserve bijectivity if two function values need to be switched around.
Proposition 6.10. Let $X, Y \neq \emptyset$, let $f: X \rightarrow Y$ be bijective and let $x_{1}, x_{2} \in X$. Let

$$
g(x):= \begin{cases}f\left(x_{2}\right) & \text { if } x=x_{1},  \tag{6.48}\\ f\left(x_{1}\right) & \text { if } x=x_{2}, \\ f(x) & \text { if } x \neq x_{1}, x_{2} .\end{cases}
$$

(In other words, we swap two function arguments). Then $g: X \rightarrow Y$ also is bijective.
Proof: Let $y_{1}:=f\left(x_{1}\right)$ and $y_{2}:=f\left(x_{2}\right)$. Let $f^{-1}: Y \rightarrow X$ be the inverse function of $f$ and define $G: Y \rightarrow X$ as follows

$$
G(y):= \begin{cases}f^{-1}\left(y_{2}\right) & \text { if } y=y_{1},  \tag{6.49}\\ f^{-1}\left(y_{1}\right) & \text { if } y=y_{2}, \\ f^{-1}(y) & \text { if } y \neq y_{1}, y_{2}\end{cases}
$$

We show that $G$ satisfies $G \circ g=i d_{X}$ and $g \circ G=i d_{Y}$, i.e., $g$ has $G$ as its inverse. This suffices to prove bijectivity of $g$.

$$
\begin{aligned}
y \neq y_{1}, y_{2} \Rightarrow g \circ G(y) & =g\left(f^{-1}(y)\right)=f\left(f^{-1}(y)\right)=y \text { as } f^{-1}(y) \neq x_{1}, x_{2}, \\
g \circ G\left(y_{1}\right) & =g\left(f^{-1}\left(y_{2}\right)\right)=g\left(x_{2}\right)=f\left(x_{1}\right)=y_{1} \text { as } f^{-1}\left(y_{2}\right)=x_{2}, \\
g \circ G\left(y_{2}\right) & =g\left(f^{-1}\left(y_{1}\right)\right)=g\left(x_{1}\right)=f\left(x_{2}\right)=y_{2} \text { as } f^{-1}\left(y_{1}\right)=x_{1} .
\end{aligned}
$$

[^44]Further,

$$
\begin{aligned}
x \neq x_{1}, x_{2} \Rightarrow \quad G \circ g(x) & \left.=G(f(x))=f^{-1} f(x)\right)=y \text { as } f(x) \neq y_{1}, y_{2}, \\
G \circ g\left(x_{1}\right) & \left.=G\left(f\left(x_{2}\right)\right)=G\left(y_{2}\right)=f^{-1}\left(y_{1}\right)\right)=x_{1} \text { as } f\left(x_{1}\right)=y_{1}, \\
G \circ g\left(x_{2}\right) & \left.=G\left(f\left(x_{1}\right)\right)=G\left(y_{1}\right)=f^{-1}\left(y_{2}\right)\right)=x_{2} \text { as } f\left(x_{2}\right)=y_{2} .
\end{aligned}
$$

We have proved that $g$ has an inverse, the function $G$.
Note that the validity of $G \circ g=i d_{X}$ and $g \circ G=i d_{Y}$ is obvious without the use of any formalism: $g$ differs from $f$ only in that it switches around the function values $f\left(x_{1}\right)$ and $f\left(x_{2}\right)$. and $G$ differs from $f^{-1}$ only in that this switch is reverted.

A more general version of the above shows how to preserve surjectivity if two function values need to be switched around.

Proposition 6.11. Let $X, Y \neq \emptyset$ and assume that $Y$ contains at least two elements $y_{1}$ and $y_{2}$. Let $f: X \rightarrow$ $Y$ be surjective.

$$
\text { Let } A_{1}:=f^{-1}\left\{y_{1}\right\}, A_{2}:=f^{-1}\left\{y_{2}\right\} \text {, and } B:=X \backslash\left(A_{1} \cup A_{2}\right) \text {. Let }
$$

$$
g(x):= \begin{cases}y_{2} & \text { if } x \in A_{1}  \tag{6.50}\\ y_{1} & \text { if } x \in A_{2} \\ f(x) & \text { if } x \in B\end{cases}
$$

In other words, everything that $f$ maps to $y_{1}$ is now mapped to $y_{2}$ and everything that $f$ maps to $y_{2}$ is now mapped to $y_{1}$. Then $g: X \rightarrow Y$ also is surjective.

## PROOF:

We notice that $A_{1}, A_{2}$, $B$ partition $X$ into three mutually exclusive parts: $X=B \biguplus A_{1} \biguplus A_{2}$
and that the sets $f\left(A_{1}\right)=\left\{y_{1}\right\}, f\left(A_{2}\right)=\left\{y_{2}\right\}, f(B)=Y \backslash\left\{y_{1}, y_{2}\right\}$
partition $Y$ into $Y=f(B) \biguplus f\left(A_{1}\right) \biguplus f\left(A_{2}\right)$. (Do you see why $f(B)=Y \backslash\left\{y_{1}, y_{2}\right\}$ ?)
$B$ and hence $f(B)$ might be empty but none of the other four sets are.
It follows that there is indeed a function value $g(x)$ for each $x \in X$ and there is exactly one such value, i.e., $g$ in fact defines a mapping from $X$ to $Y$.

The surjectivity of $g$ follows from that of $f$ and the fact that

$$
\begin{equation*}
Y=f(B) \cup f\left(A_{1}\right) \cup f\left(A_{2}\right)=g(B) \cup g\left(A_{2}\right) \cup g\left(A_{1}\right) \tag{6.51}
\end{equation*}
$$

(see (6.19) on p. 105 in prop. 6.3 (Properties of the direct image)).

### 6.2 Indicator Functions

Sometimes it is advantageous to think of the subsets of a universal set $\Omega$ as "binary" functions $\Omega \rightarrow\{0,1\}$.
Definition 6.2 (indicator function for a set). Let $\Omega$ be "the" universal set, i.e., we restrict our scope of interest to subsets of $\Omega$. Let $A \subseteq \Omega$. Let $1_{A}: \Omega \rightarrow\{0,1\}$ be the function defined as

$$
1_{A}(\omega):= \begin{cases}1 & \text { if } \omega \in A  \tag{6.52}\\ 0 & \text { if } \omega \notin A\end{cases}
$$

$1_{A}$ is called the indicator function of the set $A .{ }^{59}$
The above association of a subset $A$ of $\Omega$ with its indicator function is a bijection:
Proposition 6.12. Let $\mathscr{F}(\Omega,\{0,1\}):=\{0,1\}^{\Omega}$ denote the set of all functions $f: \Omega \rightarrow\{0,1\}$, i.e., all functions $f$ with domain $\Omega$ for which the only possible function values $f(\omega)$ are zero or one. ${ }^{60}$
a. The mapping

$$
\begin{equation*}
F: 2^{\Omega} \rightarrow \mathscr{F}(\Omega,\{0,1\}), \text { defined as } F(A):=1_{A} \tag{6.53}
\end{equation*}
$$

which assigns to each subset of $\Omega$ its indicator function is injective.
b. Let $f \in \mathscr{F}(\Omega,\{0,1\})$. Further, let $A:=\{f=1\}=f^{-1}(\{1\})=\{a \in A: f(a)=1\}$. Then $f=1_{A}$.
c. The function $F$ above is bijective and its inverse function is

$$
\begin{equation*}
G: \mathscr{F}(\Omega,\{0,1\}) \rightarrow 2^{\Omega}, \quad \text { defined as } G(f):=\{f=1\} . \tag{6.54}
\end{equation*}
$$

PROOF of $\boldsymbol{a}$ : This follows from $c$ which will be proved below.
PROOF of $\boldsymbol{b}$ : We have

$$
\begin{aligned}
f(\omega)=1 & \Leftrightarrow \omega \in\{f=1\} \quad \text { (def. of inverse image) } \\
& \Leftrightarrow \omega \in A \quad \text { (because } A=\{f=1\}) \\
& \Leftrightarrow 1_{A}(\omega)=1 \quad \text { (def. of indicator function). }
\end{aligned}
$$

It follows that $f(\omega)=1$ if and only if $1_{A}(\omega)=1$. As the only other possible function value is 0 we conclude that $f(\omega)=0$ if and only if $1_{A}(\omega)=0$. It follows that $f(\omega)=1_{A}(\omega)$ for all $\omega \in \Omega$, i.e., $f=1_{A}$. This proves $b$.
PROOF of $\boldsymbol{c}$ : According to theorem 4.1 on $p .80$ about the characterization of inverse functions $\boldsymbol{c}$ is proved if we can demonstrate that $F$ and $G$ are inverse to each other. To prove this it suffices to show that

$$
\begin{equation*}
G \circ F=i d_{2^{\Omega}} \quad \text { and } \quad F \circ G=i d_{\mathscr{F}(\Omega,\{0,1\})} \tag{6.55}
\end{equation*}
$$

Let $A \subseteq \Omega$ and $A \in \Omega$. Then

$$
G \circ F(A)=G\left(1_{A}\right)=\left\{1_{A}=1\right\}=\left\{\omega \in \Omega: 1_{A}(\omega)=1\right\}=\{\omega \in \Omega: \omega \in A\}=A .
$$

This proves $G \circ F=i d_{2} \Omega$. Let $\omega \in \Omega$. Then

$$
\begin{aligned}
(F \circ G(f))(\omega) & =F(\{f=1\})(\omega)=1_{\{f=1\}}(\omega) \\
& =\left\{\begin{array}{ll}
1 & \text { iff } \omega \in\{f=1\}, \\
0 & \text { iff } \omega \notin\{f=1\}
\end{array}=\left\{\begin{array}{ll}
1 & \text { iff } f(\omega)=1, \\
0 & \text { iff } f(\omega) \neq 1
\end{array}=\left\{\begin{array}{ll}
1 & \text { iff } f(\omega)=1, \\
0 & \text { iff } f(\omega)=0
\end{array}=f(\omega) .\right.\right.\right.
\end{aligned}
$$

The equation next to the last results from the fact that the only possible function values for $f$ are 0 and 1 . It follows that $F \circ G(f)=f=i d_{\mathscr{F}(\Omega,\{0,1\})}(f)$ for all $f \in \mathscr{F}(\Omega,\{0,1\})$, hence $F \circ G=i d_{\mathscr{F}(\Omega,\{0,1\})}$. We have proved (6.55) and hence $\boldsymbol{c}$.

The following definition is a special case of "modular arithmetic". We will only use it in the context of indicator functions of set differences. For further information we refer to [1] B/G, ch.6.3.

[^45]Definition $6.3(\bmod 2$ addition). Let $m, n \in \mathbb{Z}$. We define

$$
m+n \bmod 2:= \begin{cases}0 & \text { if } m+n \text { is even }  \tag{6.56}\\ 1 & \text { if } m+n \text { is odd }\end{cases}
$$

and we call $m+n \bmod 2$ the sum $\bmod 2$ of $m$ and $n$.
Proposition 6.13. Let $m, n, p \in \mathbb{Z}$. Then addition $\bmod 2$ is associative, i.e.,

$$
\begin{equation*}
(m+n \bmod 2)+p \bmod 2=m+(n+p \bmod 2) \bmod 2 \tag{6.57}
\end{equation*}
$$

The proof is left as exercise 6.7.
Proposition 6.14. Let $A, B, C$ be subsets of $\Omega$. Then

$$
\begin{align*}
\mathbb{1}_{A \cup B} & =\max \left(\mathbb{1}_{A}, \mathbb{1}_{B}\right),  \tag{6.58}\\
\mathbb{1}_{A \cap B} & =\min \left(\mathbb{1}_{A}, \mathbb{1}_{B}\right),  \tag{6.59}\\
\mathbb{1}_{A^{C}} & =1-\mathbb{1}_{A},  \tag{6.60}\\
\mathbb{1}_{A \triangle B} & =\mathbb{1}_{A}+\mathbb{1}_{B} \quad \bmod 2 . \tag{6.61}
\end{align*}
$$

PROOF: The proof of the first three equations is left as an exercise.
PROOF of (6.61): This follows easily from the the fact that

$$
(A \triangle B)^{\complement}=\{\omega \in \Omega: \text { either } \omega \in A \cap B \text { or } \omega \in \text { neither } A \text { nor } B\}
$$

Proposition 6.15 (Symmetric set differences $A \triangle B$ are associative). Let $A, B, C \subseteq \Omega$. Then

$$
\begin{equation*}
(A \triangle B) \triangle C=A \triangle(B \triangle C) \tag{6.62}
\end{equation*}
$$

PROOF: This follows easily from (6.61) and prop. 6.13 as follows. Let $\omega \in \Omega$. Then

$$
\omega \in(A \triangle B) \triangle C \Leftrightarrow \mathbb{1}_{(A \triangle B) \Delta C}=1 \Leftrightarrow \mathbb{1}_{A \triangle(B \triangle C)}=1 \Leftrightarrow \omega \in A \triangle(B \triangle C)
$$

We obtained the equivalence in the middle from prop.6.13.

### 6.3 Exercises for Ch. 6

Exercise 6.1. Prove prop.6.1 on p.104:
a. $\quad f(\emptyset)=f^{-1}(\emptyset)=\emptyset$
b. $A_{1} \subseteq A_{2} \subseteq X \Rightarrow f\left(A_{1}\right) \subseteq f\left(A_{2}\right)$
c. $B_{1} \subseteq B_{2} \subseteq Y \Rightarrow f^{-1}\left(B_{1}\right) \subseteq f^{-1}\left(B_{2}\right)$
d. $\quad x \in X \Rightarrow f(\{x\})=\{f(x)\}$
e. $f(X)=Y \Leftrightarrow f$ is surjective
f. $\quad f^{-1}(Y)=X \quad$ always!

Exercise 6.2. Prove (6.31) of prop.6.4 on p.106: Let $X, Y, Z$ be an arbitrary, nonempty sets. Let $f: X \rightarrow Y$ and $g: Y \rightarrow Z$ and let $W \subseteq Z$. Then $(g \circ f)(U)=g(f(U))$ for all $U \subseteq X$.

Exercise 6.3. Let $X, Y$ be nonempty sets and let $f: X \rightarrow Y$ be a function. We define on the domain $X$ a relation " $\sim$ " as follows:

$$
x_{1} \sim x_{2} \Leftrightarrow f\left(x_{1}\right)=f\left(x_{2}\right)
$$

(See prop.6.5 (Indirect image and fibers of $f$ ) on p.107). Prove that " $\sim$ " is an equivalence relation.
Exercise 6.4. Prove prop. 6.7 on p.108.
Hint: The main tools you need are prop. 6.5 on p.107, prop6.6 on p.107, and their corollaries.
Exercise 6.5. Prove prop.6.8 on p.108.
Hint: Work with the inverse of $f$ and apply prop. 6.2 on p.104.
Exercise 6.6. Prove prop. 6.9 on p. 108.
Hint: To prove a, use prop.6.3 on p.105.
Exercise 6.7. Prove prop. 6.13 on p.112: Let $m, n, p \in \mathbb{Z}$. Then

$$
(m+n \bmod 2)+p \bmod 2=m+(n+p \bmod 2) \bmod 2 .
$$

Hint: There are eight possible combinations of zeroes and ones for the functions

$$
(m, n, p) \rightarrow(m+n \bmod 2)+p \bmod 2 \quad \text { and } \quad(m, n, p) \rightarrow m+(n+p \bmod 2) \bmod 2 .
$$

Complete the entries in the table below and show that the entries in the two rightmost columns match. To save space, write $m \oplus n$ for $m+n \bmod 2$. To get you started, the row for $m=1, n=$ $0, p=0$ has been already completed.

| $m$ | $n$ | $p$ | $m \oplus n$ | $n \oplus p$ | $(m \oplus n) \oplus p$ | $m \oplus(n \oplus p$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 0 | 0 | 0 |  |  |  |  |
| 0 | 0 | 1 |  |  |  |  |
| 0 | 1 | 0 |  |  |  |  |
| 0 | 1 | 1 |  |  |  |  |
| 1 | 0 | 0 | 1 | 0 | 1 | 1 |
| 1 | 0 | 1 |  |  |  |  |
| 1 | 1 | 0 |  |  |  |  |
| 1 | 1 | 1 |  |  |  |  |

Exercise 6.8. See [1] B/G project 6.8 on p. 58 for the following.
Prove that the following are equivalence relations on $\mathbb{R}^{2}$.
a. $\quad(x, y) \sim(u, v) \Leftrightarrow \sqrt{x^{2}+y^{2}}=\sqrt{u^{2}+v^{2}}$.
b. $\quad X=\mathbb{R}_{\neq 0}=\{x \in \mathbb{R}: x \neq 0\}$ and $x \sim y \Leftrightarrow|x y|>0$.
c. $\quad X=\mathbb{R}^{3}$ and $(x, y, z) \sim(u, v, w) \Leftrightarrow z \sin (x y)=w \sin (u v)$.

Hint: See example 6.5 on p.114.

### 6.4 Addenda to Ch. 6

Example 6.5. The following are equivalence relations on the set $X$.
a. $\quad X=\mathbb{R}$ and $x \sim y \Leftrightarrow|x|=|y|$.
b. $\quad X=\mathbb{R}_{\neq 0}=\{x \in \mathbb{R}: x \neq 0\}$ and $x \sim y \Leftrightarrow|x y|>0$.
c. $\quad X=\mathbb{R}^{3}$ and $(x, y, z) \sim(u, v, w) \Leftrightarrow z \sin (x y)=w \sin (u v)$.

You can verify this by brute force, but here is an elegant way. Rewrite the equivalence relations as $\alpha \sim \beta \Leftrightarrow F(\alpha)=F(\beta)$ for a suitable function $F(\cdot)$, then apply prop.6.5 (Indirect image and fibers of $f$ ) on p.107.

For the above examples you do this as follows:
a. $\quad F: X \rightarrow \mathbb{R}, \quad x \mapsto|x|$.
b. $\quad G: X \rightarrow\{-1,1\}, \quad x \mapsto \frac{x}{|x|}$.
b. $\quad H: X \rightarrow \mathbb{R}, \quad(x, y, z) \mapsto z \sin (x y)$.

The following belongs at the end of ch. 6 (Sets and Functions, Direct and Indirect Images).

Proposition 6.16. Let $X, Y$ be two nonempty sets and let $f: X \rightarrow Y$ be surjective. Let $\emptyset \neq B \subsetneq Y$ so that $Y=B \biguplus B^{\complement}$ is a partitioning of $Y$ into two nonempty subsets $B$ and $B^{\complement}$. Let $A:=\{f \in B\}$. Then the restrictions $f_{1}:=\left.f\right|_{A}: A \rightarrow B$ and $f_{2}:=\left.f\right|_{A^{\complement}}: A^{\complement} \rightarrow B^{\complement}$ of $f$ to $A$ and $A^{\complement}$ are surjections.

PROOF: It follows from (6.11) on p. 104 that $A^{\complement}=\left(f^{-1}(B)\right)^{\complement}=f^{-1}\left(B^{\complement}\right)$. Hence $B^{\complement} \neq \emptyset$ and surjectivity of $f$ imply that $A^{\complement} \neq \emptyset$. Likewise $A=f^{-1}(B)$ is not empty because it is the preimage of the nonempty set $B$ under the surjection $f$. It follows from prop.6.7.b on $p .108$ that

$$
\begin{aligned}
f_{1}(A) & =f(A)=f \circ f^{-1}(B)=B \\
f_{2}\left(A^{\complement}\right) & =f\left(A^{\complement}\right)=f \circ f^{-1}\left(B^{\complement}\right)=B^{\complement}
\end{aligned}
$$

Because $B$ is the codomain of $f_{1}, f_{1}$ is surjective and because $B^{\complement}$ is the codomain of $f_{2}, f_{2}$ is surjective.

## 7 Some Miscellaneous Topics

Although this chapter only contains a single topic at this time (cardinality), additional topics are planned in the future.

### 7.1 Cardinality - Alternate Approach to Beck/Geoghegan (Study this!)

Notation: In this entire chapter on cardinality, if $n \in \mathbb{N}$, the symbol $[n]$ does not denote an equivalence class of any kind but the set $\{1,2, \ldots, n\}$ of the first $n$ natural numbers.

Definition 7.1 (Cardinality Comparisons). Given are two arbitrary sets $X$ and $Y$. We say that
a. $\quad X, Y$ have same cardinality and we write $\operatorname{card}(X)=\boldsymbol{\operatorname { c a r d }}(Y)$ if there is a bijective mapping $f: X \xrightarrow{\sim} Y$. or if both $X$ and $Y$ are empty.
b. the cardinality of $X$ is less than or equal to the cardinality of $Y$ and we write $\boldsymbol{\operatorname { c a r d }}(X) \leqq \boldsymbol{\operatorname { c a r d }}(Y)$ if there is an injective mapping $f: X \rightarrow Y$ or if $X$ is empty.
c. the cardinality of $X$ is less than the cardinality of $Y$ and we write $\operatorname{card}(X)<\boldsymbol{\operatorname { c a r d }}(Y)$ if both $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ and $\operatorname{card}(Y) \neq \operatorname{card}(X)$, i.e., there is an injective mapping but not a bijection $f: X \rightarrow Y$.
d. the cardinality of $X$ is greater than or equal to the cardinality of $Y$ and we write $\operatorname{card}(X) \geqq \operatorname{card}(Y)$ if $\operatorname{card}(Y) \leqq \operatorname{card}(X)$.
e. the cardinality of $X$ is greater than the cardinality of $Y$ and we write $\operatorname{card}(X)>\operatorname{card}(Y)$ if $\operatorname{card}(Y)<\operatorname{card}(X)$.

Remark 7.1. Note the following.
a. The above definition does not specify how $\operatorname{card}(X)$ itself is defined. This will be done in def.7.4 on p. 123.
b. The above definition includes the cases $A=\emptyset$ and $B=\emptyset$.

Example 7.1. Let $A, B$ be two sets such that $A \subseteq B$. Then $\operatorname{card}(A) \leqq \operatorname{card}(B)$.

## PROOF:

Case 1: $A=\emptyset$. It then is true by definition that $\operatorname{card}(\emptyset) \leqq \operatorname{card}(B)$ for any set $B$.
Case 2: $A \neq \emptyset$. It follows from $B \supseteq A$ that $B \neq \emptyset$. Further the mapping $x \mapsto x$ is injective. This proves $\operatorname{card}(A) \leqq \operatorname{card}(B)$.
Definition 7.2 (Finite and infinite sets). Let $n \in \mathbb{N}$. we say that
a. A set $X$ has size $n$ and we write $|X|=n$ if there is a bijective mapping between $X$ and the set $[n]:=\{1,2, \ldots, n\}$. In other words, a set $X$ of size $n$ is one whose elements can be enumerated as $x_{1}, x_{2}, \ldots, x_{n}$ and its size simply is the number of elements it contains.
b. We say that the empty set $\emptyset$ has size $|\emptyset|=0$.
c. A set $X$ is finite if there exists $n \in \mathbb{Z}_{\geqq}$such that $X$ has size $n$. We say that $X$ is infinite and we write $|X|=\infty$ if $X$ is not finite.

As strange as this may seem, there are ways to classify the degree of infinity when looking at infinite sets. The "smallest degree of infinity" is found in sets that can be compared, in a sense, to the set $\mathbb{N}$ of all natural numbers. We have a special name for infinite sets whose elements can be mapped bijectively to $\mathbb{N}$.
d. Let $X$ be a set such that there is a bijection $f: \mathbb{N} \xrightarrow{\sim} X$. This means that all of the elements of $X$ can be arranged in a sequence

$$
X=\left\{x_{1}=f(1), x_{2}=f(2), x_{3}=f(3), \ldots\right\} .
$$

which is infinite, i.e., we rule out the case of sets with finitely many members. $X$ is called a countably infinite set.
e. We call a set $X$ that is either finite or countably infinite a countable set.
f. A set that is neither finite nor countably infinite is called uncountable or not countable.

Example 7.2. Let $X:=\{2 n: n \in \mathbb{N}\}$ be the set of all even natural numbers. Then $X$ is countably infinite because the function

$$
f: \mathbb{N} \longrightarrow X ; n \mapsto 2 n \quad \text { has the function } \quad g: X \longrightarrow \mathbb{N} ; k \mapsto \frac{k}{2}
$$

as its inverse and hence is bijective.
Note that because the even natural numbers are divisible by $2, \frac{k}{2} \in \mathbb{N}$ if $k \in X$.

## Proposition 7.1.

Let $X$ and $Y$ be two nonempty sets and $f: X \xrightarrow{\sim} Y$. Then
a. $Y$ is finite if and only if $X$ is finite,
b. $Y$ is countably infinite if and only if $X$ is countably infinite,
c. $Y$ is countable if and only if $X$ is countable,
d. $Y$ is uncountable if and only if $X$ is uncountable.
e. $|Y|=|X|$.

PROOF: The proof of $\boldsymbol{a}, \boldsymbol{b}$ and $\boldsymbol{e}$ is based on prop.4.1.c on $p .81$ which states that the composition of two bijective functions is bijective.

We only need to prove the " $\Rightarrow$ " directions because we obtain " $\Leftarrow$ " by switching the roles of $X$ and $Y$.
PROOF of $\boldsymbol{a}$ and $\boldsymbol{e}$. If $X$ is finite then there exists $n \in \mathbb{N}$ and a bijection $g: X \xrightarrow{\sim}[n] . Y \xrightarrow{f^{-1} \circ g}[n]$ is bijective according to prop.4.1.c on p.81. This proves both that $Y$ is finite and $|Y|=n=|X|$.
PROOF of $\boldsymbol{b}$. If $X$ is countably infinite then there exists a bijection $g: X \xrightarrow{\sim} \mathbb{N}$. Because $Y \xrightarrow{f^{-1} \circ g} \mathbb{N}$ is bijective, $Y$ also is countably infinite.
PROOF of $c$. If $X$ is countable then this set is either finite or countably infinite. If $X$ is finite then $Y$ is finite according to part $\boldsymbol{a}$; if $X$ is countably infinite then $Y$ is countably infinite according to part $\boldsymbol{b}$. This proves that $Y$ is countable.
PROOF of $\boldsymbol{d}$. Assume to the contrary that $X$ is uncountable and $Y$ is countable. It follows from part $\boldsymbol{c}$ that $X$ is countable and we have reached a contradiction. This proves that $Y$ is uncountable.

Placeholder for A. thm.7.7 on p.126:
Let $X$ be a set. Then $\operatorname{card}(X)<\operatorname{card}\left(2^{X}\right)$.
and B. lemma 7.2 on p. 126
Proposition 7.2. Let $n \in \mathbb{N}$. Let $\emptyset \neq A \subsetneq[n]$ be a proper, nonempty subset of $[n]$. Then there is no surjection from $A$ onto $[n]$.

PROOF by induction on $n$ :
Base case: Let $n=2$. The only proper subsets of $[2]$ are the singleton sets $\{1\}$ and $\{2\}$. For any function $f:\{1\} \rightarrow[2]$ we have either $f(1)=1$ in which case $2 \notin f(\{1\})$ or $f(1)=2$ in which case $1 \notin f(\{1\})$. It follows in either case that $f$ is not surjective. The proof for functions $\{2\} \rightarrow[2]$ is similar. This proves the base case.
Induction assumption: Let $n \in \mathbb{N}$ such that if $\emptyset \neq \Gamma \subsetneq[n-1]$ then there is no surjection $\Gamma \rightarrow[n-1]$.
We must prove that if $\emptyset \neq A \subsetneq[n]$ then there is no surjection from $A$ to $[n]$.
We assume to the contrary that a surjective $f: A \rightarrow[n]$ exists.
case 1: $n \notin A$ :
Then $A \subseteq[n-1]$.
Let $\Gamma:=A \backslash\{f=n\}$. Because $f$ is surjective, $\{f=n\} \neq \emptyset$, hence $\Gamma \subsetneq[n-1]$. It follows from prop.6.16 on p. 114 that $f_{\Gamma}:=\left.f\right|_{\Gamma}: \Gamma \rightarrow[n] \backslash\{n\}=[n-1]$ is surjective.

But the induction assumption implies that there is no surjection from the proper subset $\Gamma$ of $[n-1]$ onto [ $n-1]$. We have reached a contradiction.
case 2: $n \in A$ :
Because $A \subsetneq[n]$ there exists $j \in[n-1]$ such that $j \notin A$. Let $A^{\prime}:=(A \backslash\{n\}) \biguplus\{j\}$. It follows from lemma 7.2 that there is a bijection $g: A^{\prime} \xrightarrow{\sim} A$. Hence $f \circ g: A^{\prime} \rightarrow[n]$ is surjective as the composition of a surjection with a bijection (see cor.4.1.b on p.81). But $A^{\prime}$ satisfies the conditions of case 1 and we have already proven that such a surjection cannot exist.

Corollary 7.1 (No bijection from $[m]$ to $[n]$ exists). B/G Thm.13.4: Let $m, n \in \mathbb{N}$. If $m \neq n$ then there is no bijective $f:[m] \xrightarrow{\sim}[n]$.

PROOF: We may assume $m<n$ and can now apply prop. 7.2 with $A:=[m]$.
Corollary 7.2 (Pigeonhole Principle). B/G Prop.13.5: Let $m, n \in \mathbb{N}$. If $m<n$ then there is no injective $f:[n] \rightarrow[m]$.

PROOF: Otherwise, by thm.4.2 on $p .83, f$ would have a (surjective) left inverse $g:[m] \rightarrow[n]$ in contradiction to the preceding proposition.

Proposition 7.3 (B/G Prop.13.6, p.122: Subsets of finite sets are finite). Let $\emptyset \neq B \subseteq A$ and let $A$ be finite. Then $B$ is finite.

PROOF: Done by induction on the size $n$ of $A$. We may assume that $B \subset A$ because bijectivity is reflexive (the identity mapping $i d_{A}$ maps $A$ bijectively to itself).

Base case: $n=1$ or $n=2$ : Proof obvious.
Induction assumption: Assume that all subsets of sets of size less than $n$ are finite.

Now let $A$ be a set with size $|A|=n$. There is a bijection $a(\cdot):[n] \xrightarrow{\sim} A$. Let $B \subseteq A$.
Case 1: $a(n) \in B$ : Let $B_{n}:=B \backslash\{a(n)\}$ and $A_{n}:=A \backslash\{a(n)\}$.
Then the restriction $\left.a(\cdot)\right|_{[n-1]}$ of $a(\cdot)$ to $[n-1]$ is a bijection $[n-1] \xrightarrow{\sim} A_{n}$ according to cor.6.3 on p.109.
As $\left|A_{n}\right|=n-1$ and $B_{n} \subseteq A_{n}$ it follows from the induction assumption that $B_{n}$ is finite: there exists $m \in \mathbb{N}$ and a bijection $b(\cdot):[m] \xrightarrow{\sim} B_{n}$.

We now extend $b(\cdot)$ to $[m+1]$ by defining $b(m+1):=a(n)$. It follows that this extension remains injective and it is also surjective if we choose as codomain $B_{n} \cup\{a(n)\}=B$.

It follows that $B$ is finite.
Case 2: $a(n) \notin B$ : We pick an arbitrary $b \in B$. Let $j:=a^{-1}(b)$. Clearly $j \in[n]$.
We modify the mapping $a(\cdot)$ by switching the function values for $j$ and $n$. We obtain another bijection $f:[n] \xrightarrow{\sim} A$ (see prop.6.10 on $p$. 109) for which $f(n)=a(j)=b \in B$.

We apply to $f$ what was proved in case 1 and obtain that $B$ is finite.
Our next major result is that subsets of countable sets are countable.
Lemma 7.1. Any subset of $\mathbb{N}$ is countable.
PROOF: Let $B \subseteq \mathbb{N}$. Finite sets are countable by definition, so we may assume that $B$ is infinite. We extract $j_{1}<j_{3}<j_{2}<\ldots$ from $\mathbb{N}$ as follows.

$$
\begin{aligned}
& j_{1}:=\min \{j \in \mathbb{N}: j \in B\}, \\
& j_{2}:=\min \left\{j \in B: j>j_{1}\right\}, \\
& \ldots \ldots \ldots \ldots \ldots \ldots \cdots \\
& j_{k+1}:=\min \left\{j \in B: j>j_{k}\right\}
\end{aligned}
$$

Note that none of the sets $B_{k}:=\left\{j \in B: j>j_{k}\right\}$ are empty because $B_{k}=\emptyset$ would imply $B \subseteq$ $\left\{1,2, \ldots, j_{k}\right\}$, i.e., $B$ is a subset of the finite set $\left[j_{k}\right]$. This is a contradiction to prop.7.3 (Subsets of finite sets are finite) on p.117.

Let $B^{\prime}:=\left\{j_{k}: k \in \mathbb{N}\right\}$. Note that $j_{k}<j_{k+1}$ for all $k \in \mathbb{N}$ by construction. Hence

$$
\begin{equation*}
j_{k} \geqq k \quad \text { for all } k \in \mathbb{N} \tag{7.1}
\end{equation*}
$$

In particular, $B^{\prime}$ is not bounded above.
We claim that $B=B^{\prime}$. It follows from $j_{k} \in B$ for all $k$ that $B^{\prime} \subseteq B$. We now show that $B \subseteq B^{\prime}$. Otherwise there would be $n^{\prime} \in B$ such that $n^{\prime} \neq j_{k}$ for all natural numbers $k$. It follows from 7.1 that $n^{\prime} \leqq j_{n^{\prime}}$. Hence there is $m \leqq n^{\prime}$ such that $j_{m-1}<n^{\prime}<j_{m}$.

This contradicts $j_{m}=\min \left\{j \in B: j>j_{m-1}\right\}$.
It follows from 7.1 that the function

$$
\psi: \mathbb{N} \rightarrow B ; \quad k \mapsto j_{k}
$$

is injective. $\psi$ also is surjective because its codomain is $B^{\prime}=\left\{b_{k}: k \in \mathbb{N}\right\}$. Hence $B$ is countably infinite, hence $B$ is countable.

Theorem 7.1. Subsets of countable sets are countable.

PROOF: Let $X$ be countable and $A \subseteq X$.
Before we continue, note that no assumption is made about the nature of $X$ and $A$. We may not deal with natural numbers or any other kind of numbers. $X$ and hence $A$ might, e.g., be sets of functions or sets of sets.

We may assume that $X$ is infinite because otherwise $A$ is finite, hence countable, as the subset of a finite set. See prop.7.3 on p.117. $X$ is countably infinite, hence there exists a bijection

$$
\phi: \mathbb{N} \xrightarrow{\sim} X ; \quad n \mapsto x_{n}:=x(n) .
$$

The function $\left.\quad \phi^{-1}\right|_{A}: A \rightarrow \phi^{-1}(A), \quad x \mapsto \phi^{-1}(x)$,
is injective as a restriction of the bijective, hence injective function $\phi^{-1} .\left.\phi^{-1}\right|_{A}$ also is surjective because the codomain is $\phi^{-1}(A)=\left.\phi^{-1}\right|_{A}(A)$. It follows that $\left.\phi^{-1}\right|_{A}$ is bijective.

The diagram to the right shows the role of $\phi$ and its restriction to $A$. The bijectivity of those functions allows us to shift our problem from a subset $A$ of an abstract set $X$ to a subset of $\mathbb{N}$ whose cardinality equals that of $A$.


According to lemma 7.1, $\phi^{-1}(A)=\{n \in \mathbb{N}: \phi(x) \in A\}$ is countable as a subset of $\mathbb{N}$. It follows from the bijectivity of $\left.\phi^{-1}\right|_{A}$ and prop.7.1.c on p. 116 that $A$ is countable.

Corollary 7.3. Supersets of uncountable sets are uncountable.
Proof: Assume to the contrary that $X$ is uncountable, $Y \supseteq X$, and $Y$ is countable. This can be worded as follows: $Y$ is countable, $X \subseteq Y$, and $X$ is uncountable. We have a contradiction to theorem 7.1.

The next proposition is a major stepping stone for proving that countable unions of countable sets are countable.

Proposition 7.4 (B/G Cor.13.16, p.122). $\mathbb{N}^{2}$ is countable.
PROOF: ${ }^{61}$ Done by directly specifiying a bijection $F: \mathbb{N}^{2} \xrightarrow{\sim} \mathbb{N}$.
The following definitions and observations will make it easier to understand this proof. Let

$$
s_{0}:=0 ; \quad s_{n}:=\sum_{j=1}^{n} j=1+2+\cdots+n=\frac{n(n+1)}{2} .
$$

For the last equality see [1] B/G prop.4.11 on p.37. We note that

$$
\begin{equation*}
s_{n-1}+n=s_{n} . \tag{7.2}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
A_{n}:=\left\{j \in \mathbb{N}: s_{n-1}<j \leqq s_{n}\right\}=\left\{s_{n-1}+1, s_{n-1}+2, \ldots, s_{n-1}+n\right\} \quad(n \in \mathbb{N}) \tag{7.3}
\end{equation*}
$$

[^46]is a partition of $\mathbb{N}$ (see example 5.2.d on $p .97$ ). For $n \in \mathbb{N} \backslash\{1\}=\{2,3,4, \ldots\}$ let
$$
D_{n}:=\left\{(i, j) \in \mathbb{N}^{2}: i+j=n\right\}
$$
be the set of all pairs of natural numbers whose sum equals $n$. Clearly, $\left(D_{n}\right)_{n \geqq 2}$ is a partition of $\mathbb{N}^{2}$. For $n \in \mathbb{N}$ let
\[

$$
\begin{equation*}
f_{n}: A_{n} \rightarrow D_{n+1}, \quad s_{n-1}+k \mapsto f_{n}\left(s_{n-1}+k\right):=(n+1-k, k) \quad(1 \leqq k \leqq n) \tag{7.4}
\end{equation*}
$$

\]

We see from the second equality in (7.3) that the set of argument values $s_{n-1}+k(1 \leqq k \leqq n)$ coincides with the domain $A_{n}$ and it follows that (7.4) indeed defines a function $A_{n} \rightarrow \mathbb{N}^{2}$.

It is immediate that $D_{n+1}=\{(n+1-k, k): 1 \leqq k \leqq n\}$. We conclude that $\left.f_{n}\left(A_{n}\right)\right)=D_{n+1}$ and we have proved surjectivity of $f_{n}$.

To prove injectivity of $f_{n}$, we observe that if $1 \leqq k, k^{\prime} \leqq n$ and $k \neq k^{\prime}$ then

$$
f_{n}\left(s_{n-1}+k\right)=(n+1-k, k) \neq\left(n+1-k^{\prime}, k^{\prime}\right)=f_{n}\left(s_{n-1}+k^{\prime}\right) .
$$

We "glue together" the functions $f_{n}$ to obtain a function $f$ with domain $\bigcup\left[A_{n}: n \in \mathbb{N}\right]=\mathbb{N}$ and codomain $\bigcup\left[D_{n+1}: n \in \mathbb{N}\right]=\mathbb{N}^{2}$ as follows:
$f(m):=f_{n}(m)$ for $m \in A_{n}$, i.e.,

$$
f\left(s_{n-1}+k\right)=f_{n}\left(s_{n-1}+k\right)=(n+1-k, k) \text { for } 1 \leqq k \leqq n
$$

$f$ inherits injectivity from the individual $f_{n}$ as the ranges $f_{n}\left(A_{n}\right)=D_{n+1}$ are mutually disjoint for different values of $n$ and $f$ is surjective because

$$
\begin{aligned}
f(\mathbb{N}) & =f\left(\bigcup\left[A_{n}: n \in \mathbb{N}\right]\right)=\bigcup\left[f\left(A_{n}\right): n \in \mathbb{N}\right] \\
& =\bigcup\left[f_{n}\left(A_{n}\right): n \in \mathbb{N}\right]=\bigcup\left[D_{n+1}: n \in \mathbb{N}\right]=\mathbb{N}^{2} .
\end{aligned}
$$

To summarize, we have proved that $f$ is a bijective mapping between $\mathbb{N}$ and $\mathbb{N}^{2}$ and this proves that $\mathbb{N}^{2}$ is countable.
Remark 7.2. The following will help to visualize the proof just given. We think of $\mathbb{N}^{2}$ as a matrix with "infinitely many rows and columns"

$$
\begin{align*}
& (1,1)(1,2)(1,3) \ldots  \tag{7.5}\\
& (2,1)(2,2)(2,3) \ldots  \tag{7.6}\\
& (3,1)(3,2)(3,3) \ldots \tag{7.7}
\end{align*}
$$

We reorganize this matrix into an ordinary sequence $(f(j))_{j \in \mathbb{N}}$ as follows:

$$
\begin{align*}
& f(1)=f_{1}(1)=(1,1)  \tag{7.8}\\
& f(2)=f_{2}(2)=(1,2), f(3)=f_{2}(3)=(2,1),  \tag{7.9}\\
& f(4)=f_{3}(4)=(1,3), f(5)=f_{3}(5)=(2,2), f(6)=f_{3}(6)=(3,1),  \tag{7.10}\\
& f(7)=f_{4}(7)=(1,4), f(8)=f_{4}(8)=(2,3), f(9)=f_{4}(9)=(3,2), f(10)=f_{4}(10)=(4,1), \tag{7.12}
\end{align*}
$$

In other words, we traverse first $D_{2}$, then $D_{3}$, then $D_{4}, \ldots$ starting for each $D_{n}$ at the upper right $(1, n-1)$ and ending at the lower left $(n-1,1))$.

Theorem 7.2 (B/G prop.13.6: Countable unions of countable sets). The union of countably many countable sets is countable.

PROOF: Let the sets

$$
A_{1}, A_{2}, A_{3}, \ldots \text { be countable and let } A:=\bigcup_{n \in \mathbb{N}} A_{i} \text {. }
$$

We assume that at least one of those sets is not empty: otherwise their union is empty, hence finite, hence countable and we are done.

We may assume that the sets $A_{i}$ are mutually disjoint, i.e., any two different sets $A_{i}, A_{j}$ have intersection $A_{i} \cap A_{j}=\emptyset$. Otherwise we replace each $A_{n}$ with $A_{n} \backslash \bigcup\left[A_{j}: j<n\right]$ (see prop.5.3 on $p$. pagerefprop-x:math330-union-as-disjoint-u).
A. As each of the nonempty $A_{i}$ is countable, either $A_{i}$ is finite and we have an $N_{i} \in \mathbb{N}$ and a bijective mapping $a_{i}(\cdot): A_{i} \xrightarrow{\sim}\left[N_{i}\right]$, or $A_{i}$ is countably infinite and we have a bijective mapping $a_{i}(\cdot): A_{i} \xrightarrow{\sim} \mathbb{N}$.

We now define the mapping $f: A \rightarrow \mathbb{N}^{2}, a \mapsto\left(i_{a}, j_{a}\right)$ as follows: Let $a \in A$. As the $A_{j}$ are disjoint there is a unique index $i_{a}$ such that $a \in A_{i_{a}}$ and, as sets do not contain duplicates of their elements, there is a unique index $j_{a}$ such that $a=a_{\left(i_{a}, j_{a}\right)}$.

In other words, for any $a \in A$ there exists a unique pair $\left(i_{a}, j_{a}\right) \in \mathbb{N}^{2}$ such that $a=a_{\left(i_{a}, j_{a}\right)}$ and the assignment $a \mapsto\left(i_{a}, j_{a}\right)$ defines the injective function $f: A \rightarrow \mathbb{N}^{2}$ we were looking to create.

We shrink the codomain from $\mathbb{N}^{2}$ to $f(A)$ and the assignment $a \mapsto\left(i_{a}, j_{a}\right)$ gives us a bijective function $F: A \xrightarrow{\sim} f(A)$.
$f(A)$ is countable as a subset of the countable set $\mathbb{N}^{2}$ and this proves the theorem because any subset of a countable set is countable (see prop.7.1 on p.118).
Corollary 7.4. Let the set $X$ be uncountable and let $A \subseteq X$ be countable. Then the complement $A^{\complement}$ of $A$ is uncountable.

The proof is left as exercise 7.3.
Here is another corollary to thm.7.2.
Corollary 7.5. The set $Z$ of all integers is countable.
PROOF: The set $-\mathbb{N}$ is countable because the function $n \mapsto-n$ is a bijection $\mathbb{N} \xrightarrow{\sim}-\mathbb{N}$, hence

$$
Z=\mathbb{N} \cup-\mathbb{N} \cup\{0\}
$$

is countable as the union of three countable sets.

Placeholder for cor.7.8 on p.127: The rational numbers are countable.

We now examine the cardinality of cartesian products.
Theorem 7.3 (Finite Cartesians of countable sets are countable). The Cartesian product of finitely many countable sets is countable.

Proof by induction: Let $X:=X_{1} \times \cdots \times X_{n}$ We may assume that none of the factor sets $X_{j}$ is empty: Otherwise the Cartesian is empty too and there is nothing to prove.

The proof is a triviality for $k=1$. It is more instructive to choose $k=2$ for the base case instead.
So let $X_{1}, X_{2}$ be two nonempty countable sets. We now prove that $X_{1} \times X_{2}$ is countable.
For fixed $x_{1} \in X_{1}$ the function $F_{2}: X_{2} \rightarrow\left\{x_{1}\right\} \times X_{2} ; \quad x_{2} \mapsto\left(x_{1}, x_{2}\right)$ is bijective because it has as an inverse the function $G_{2}:\left\{x_{1}\right\} \times X_{2} \rightarrow X_{2} ; \quad\left(x_{1}, x_{2}\right) \mapsto x_{2}$. It follows that $\left\{x_{1}\right\} \times X_{2}$ is countable.

Hence $X_{1} \times X_{2}=\bigcup_{x \in X_{1}}\left\{x_{1}\right\} \times X_{2}$ is countable according to thm.7.2 on $p .121$. We have proved the base case.

Our induction assumption is that $X_{1} \times \cdots \times X_{k}$ is countable. We must prove that $X_{1} \times \cdots \times X_{k+1}$ is countable. We can "identify"

$$
\begin{equation*}
X_{1} \times \cdots \times X_{k+1}=\left(X_{1} \times \cdots \times X_{k}\right) \times X_{k+1} \tag{7.13}
\end{equation*}
$$

by means of the bijection $\left(x_{1}, \ldots, x_{n}, x_{n+1}\right) \mapsto\left(\left(x_{1}, \ldots, x_{n}\right), x_{n+1}\right)$. According to the induction assumption the set $X_{1} \times \cdots \times X_{k}$ is countable.

The proof for the base case shows that $X_{1} \times \cdots \times X_{k+1}$ as the Cartesian product of the two sets $X_{1} \times$ $\cdots \times X_{k}$ and $X_{k+1}$ is countable. This finishes the proof of the induction step.

Corollary 7.6. Let $n \in \mathbb{N}$. The sets $\mathbb{Q}^{n}$ and $\mathbb{Z}^{n}$ are countable.
PROOF: This follows from the preceding theorem because the sets $\mathbb{Q}$ and $\mathbb{Z}$ are countable.
Theorem 7.4. The real numbers are uncountable.
Proof by contradiction: This proof follows the one given in B/G [1] (thm.13.22, p.125).
According to cor. 7.3 on $p .119$ it is enough to prove that the subset $A$ of all decimals $x=0 . d_{1} d_{2} d_{3} \ldots$ where each $d_{n}$ is a digit $0,1,2, \ldots, 8$ is uncountable. ${ }^{62}$

Let us assume assume that $A$ is countable, i.e., there exists a finite or infinite sequence of decimals $x_{n}=$ $0 . d_{n, 1} d_{n, 2} d_{n, 3} \ldots$ such that $0 \leqq d_{n, j} \leqq 8$ for each digit $d_{n, j}$ as follows: for each $x \in A$ there exists some $k=k(x) \in \mathbb{N}$
such that $x=x_{k}$. We will find a contradiction.
For $n \in \mathbb{N}$ we define the digit $\tilde{d}_{n}$ as

$$
\tilde{d}_{n}:= \begin{cases}3 & \text { if } d_{n, n} \neq 3, \\ 4 & \text { if } d_{n, n}=3 .\end{cases}
$$

Let $y:=0 . \tilde{d}_{1} \tilde{d}_{2} \tilde{d}_{3} \ldots$. Clearly $y \in A$. It follows from $d_{1,1} \neq \tilde{d}_{1}$ that $y \neq x_{1}$. It follows from $d_{2,2} \neq \tilde{d}_{2}$ that $y \neq x_{2}$. You get the idea: Let $k \in \mathbb{N}$. It follows from $d_{k, k} \neq \tilde{d}_{k}$ that $y \neq x_{k}$. This contradicts the assumption that for each $x \in A$ there exists some $k=k(x) \in \mathbb{N}$ such that $x=x_{k}$.

> Here is the proper spot for thm. 7.8 (all intervals of real numbers have the cardinality ofR) on p. 127 in ch. 7.3 (Addenda to Ch.7)

Definition 7.3 (algebraic numbers). Let $x \in \mathbb{R}$ be the root (zero) of a polynomial with integer coefficients. We call such $x$ an algebraic number and we call any real number that is not algebraic a transcendental number.

[^47]Proposition 7.5 (B/G Prop.13.21, p.125). The set of all algebraic numbers is countable.
PROOF: Let $P$ be the set of all integer polynomials and $Z$ the set of zeroes for all such polynomials. Let

$$
\begin{equation*}
P_{n}:=\left\{\text { polynomials } p(x)=\sum_{j=0}^{n} a_{j} x^{j}: a_{j} \in \mathbb{Z} \text { and }-n \leqq a_{j} \leqq n\right\} . \tag{7.14}
\end{equation*}
$$

Then $P_{n}$ is finite and

$$
\begin{equation*}
Z_{n}:=\left\{x \in \mathbb{R}: p(x)=0 \text { for some } p \in P_{n}\right\} \tag{7.15}
\end{equation*}
$$

also is finite as a polynomial of degree $n$ has at most $n$ zeroes (see [1] B/G (Beck/Geoghegan), prop.6.20).
$Z$ is the countable union of the sets $Z_{n}$. It follows that $Z$ is countable.

Placeholder for a. prop.7.10 on p.127: If $k, m, n \in \mathbb{N}$. then $\sqrt[k]{\frac{m}{n}}$ is algebraic and b. prop.7.11 on p.127: Rational numbers are algebraic.

Here are some trivial consequences of the fact that $\mathbb{R}$ is uncountable (see thm. 7.4, p.7.4 and $B / G$ Thm.13.22).

Proposition 7.6. All transcendental numbers are uncountable.
PROOF: the uncountable real numbers are the disjoint union of the countable algebraic numbers with the transcendentals. The assertion follows from cor.7.4.

### 7.1.1 Cardinality as a Partial Ordering

We assume in this subchapter that all sets are subsets of a universal set $\Omega$. Having such a universal set allows us to declare on its power set $2^{\Omega}$ equivalence relations. If we work with specific sets, e.g. the set $\mathbb{R}$ of all real numbers, we assume implicitly that those sets are contained in $\Omega$.

We defined in def.7.1 on p. 115 the meaning of $\operatorname{card}(X)=\operatorname{card}(Y)$ and $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ for two sets $X$ and $Y$ but we never defined the expression card $(X)$ per se. This will be done now.

Definition 7.4 (Cardinality as an Equivalence Class). Let $X, Y \subseteq \Omega$. We say that $X$ and $Y$ are equivalent and we write $X \sim Y$ if and only if $X=Y$ or both $X$ and $Y$ are nonempty and there is a bijective function $f: X \rightarrow Y,{ }^{63}$ i.e., if and only if either $X$ and $Y$ are empty or $\operatorname{card}(X)=\operatorname{card}(Y)$.

The proposition following this definition shows that " $\sim$ " is indeed an equivalence relation on $2^{\Omega}$.

We now define for a set $X \subseteq \Omega$ its cardinality as follows:

$$
\begin{equation*}
\operatorname{card}(X):=\{Y \subseteq \Omega: \exists \text { bijection } X \rightarrow Y\} \tag{7.16}
\end{equation*}
$$

In other words, $\operatorname{card}(X)$ is the equivalence class $[X]$ of $X$ for the relation " $\sim$ ".
Proposition 7.7. $X \sim Y$ as defined above is an equivalence relation on $2^{\Omega}$.

[^48]The proof is left as exercise 7.4.
The following theorem allows us to prove that the relation $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ on $\mathscr{A}:=\{\operatorname{card}(X)$ : $X \subseteq \Omega\}$ is antisymmetric.

Theorem 7.5 (Cantor-Schröder-Bernstein's Theorem). Let $X, Y \subseteq \Omega$ be nonempty. Let there be functions $f^{\prime}: X \rightarrow Y$ and $g^{\prime}: Y \rightarrow X$ which both are injective. Then there exists a bijection $X \rightarrow Y$.

PROOF: The proof given here follows closely that of [1] B/G (Beck/Geoghegan), Further Topics F: Cardinal Number and Ordinal Number.

We have no interest in any particulars of the sets $X$ and $Y$. We only are interested in establishing the existence of a bijection $h: X \rightarrow Y$. We hence may assume that $X$ and $Y$ are mutually disjoint, replacing $X$ with $\{1\} \times X$ and $Y$ with $\{2\} \times Y$ if necessary (see remark 4.12 on $p .83$ ).

Let $f: X \xrightarrow{\sim} f^{\prime}(X)$ and $g: Y \xrightarrow{\sim} g^{\prime}(Y)$ be the bijective functions obtained from the injections $f^{\prime}$ and $g^{\prime}$ by restricting their codomains to the images of their domains. We note that the subsets $f^{\prime}(X) \subseteq Y$ and $g^{\prime}(Y) \subseteq X$ also are disjoint and that $f(X)=f^{\prime}(X), g(Y)=g^{\prime}(Y)$. Let

$$
\sigma: f(X) \biguplus g(Y) \rightarrow X \biguplus Y ; \quad z \mapsto \begin{cases}f^{-1}(z) & \text { if } z \in f(X),  \tag{7.17}\\ g^{-1}(z) & \text { if } z \in g(Y),\end{cases}
$$

i.e., $\left.\sigma\right|_{f(X)}=f^{-1}$ and $\left.\sigma\right|_{g(Y)}=g^{-1}$. Note that if $y \in f(X)$ then $\sigma(y) \in X$; if $x \in g(Y)$ then $\sigma(x) \in Y$. We can create iterates

$$
\sigma^{2}(z)=\sigma(\sigma(z)), \quad \sigma^{3}(z)=\sigma\left(\sigma^{2}(z)\right), \ldots, \sigma^{n+1}(z)=\sigma\left(\sigma^{n}(z)\right), \ldots
$$

just as long as $\sigma^{n}(z) \in f(X) \biguplus g(Y)$. We further define $\sigma^{0}$ for all $z \in f(X) \biguplus g(Y)$ as $\sigma^{0}(z):=z$. We associate with each $z \in X \biguplus Y$ a "score" $N(z) \in \mathbb{Z}_{\geqq 0} \cup\{\infty\}$ as follows.
a. If $\sigma^{k}(z) \in f(X) \biguplus g(Y)$ for all $k \in \mathbb{N}$ then $N(z):=\infty$.
b. If $\sigma^{k}(z) \notin f(X) \biguplus g(Y)$ for some $k \in \mathbb{N}$ then $N(z):=\min \left\{j \geqq 0: \sigma^{j}(z) \notin f(X) \biguplus g(Y)\right\}$.

Note that b implies the following: If $z=\sigma^{0}(z) \notin f(X) \biguplus g(Y)$ then $N(z)=0$.
Depending on whether we start out with $x \in g(Y)$ or $y \in f(X)$, we obtain the following finite or infinite sequences:

$$
\begin{aligned}
& \text { if } x \in g(Y): x \xrightarrow{\sigma} f^{-1}(x) \xrightarrow[\rightarrow]{\sigma} g^{-1}\left(f^{-1}(x)\right) \xrightarrow[\rightarrow]{\sigma} f^{-1}\left(g^{-1}\left(f^{-1}(x)\right)\right) \xrightarrow[\rightarrow]{\sigma} \ldots, \\
& \text { if } y \in f(X): y \xrightarrow{\sigma} g^{-1}(y) \xrightarrow{\sigma} f^{-1}\left(g^{-1}(y)\right) \xrightarrow{\sigma} g^{-1}\left(f^{-1}\left(g^{-1}(y)\right)\right) \xrightarrow{\sigma} \ldots .
\end{aligned}
$$

If $N(z)<\infty$ then the sequence will terminate after $N(z)$ iterations. Let

$$
\begin{aligned}
& X_{E}:=\{x \in X: N(x) \text { is even }\}, X_{O}:=\{x \in X: N(x) \text { is odd }\}, X_{\infty}:=\{x \in X: N(x)=\infty\}, \\
& Y_{E}:=\{y \in Y: N(y) \text { is even }\}, Y_{O}:=\{y \in Y: N(y) \text { is odd }\}, Y_{\infty}:=\{y \in Y: N(y)=\infty\} .
\end{aligned}
$$

The above defines partitions $X=X_{E} \biguplus X_{O} \biguplus X_{\infty}$ and $Y=Y_{E} \biguplus Y_{O} \biguplus Y_{\infty}$ of $X$ and $Y$.
Each of the functions $f, g, \sigma$ changes the score of its argument from odd to even and from even to odd. Hence

$$
\begin{align*}
f\left(X_{E}\right) \subseteq Y_{O}, f^{-1}\left(Y_{O}\right) \subseteq X_{E}, & f\left(X_{O}\right) \subseteq Y_{E}, \quad f^{-1}\left(Y_{E}\right) \subseteq X_{O}, \\
g\left(Y_{E}\right) \subseteq X_{O}, g^{-1}\left(X_{O}\right) \subseteq Y_{E}, & g\left(Y_{O}\right) \subseteq X_{E}, \quad g^{-1}\left(X_{E}\right) \subseteq Y_{O},  \tag{7.18}\\
f\left(X_{\infty}\right) \subseteq Y_{\infty}, f^{-1}\left(Y_{\infty}\right) \subseteq X_{\infty}, & g\left(Y_{\infty}\right) \subseteq X_{\infty}, \quad g^{-1}\left(X_{\infty}\right) \subseteq Y_{\infty},
\end{align*}
$$

We define a bijection $h: X \xrightarrow{\sim} Y$ as follows:

$$
h: X \rightarrow Y ; \quad x \mapsto \begin{cases}\sigma(x)=g^{-1}(x) & \text { if } x \in X_{O} \biguplus X_{\infty}, \\ f(x) & \text { if } x \in X_{E}\end{cases}
$$

Note that $g^{-1}(x)$ is defined for all $x \in X_{O} \biguplus X_{\infty}$ because we then have $N(x)>0$.
We show that $h$ is injective: Let $x_{1}, x_{2} \in X$ such that $x_{1} \neq x_{2}$. There are four cases.
Case 1: Both $x_{1}, x_{2} \in X_{O} \biguplus X_{\infty}$. Then $h\left(x_{1}\right)=g^{-1}\left(x_{1}\right) \neq g^{-1}\left(x_{2}\right)=h\left(x_{2}\right)$ because the bijectivity of $g$ implies that of $g^{-1}$. In particular, $g^{-1}$ is injective.
Case 2: Both $x_{1}, x_{2} \in X_{E}$. Then $h\left(x_{1}\right)=f\left(x_{1}\right) \neq f\left(x_{2}\right)=h\left(x_{2}\right)$ because $f$ is injective.
Case 3: $x_{1} \in X_{O} \biguplus X_{\infty}$ and $x_{2} \in X_{E}$. It follows from (7.18) that $h\left(x_{1}\right)=g^{-1}\left(x_{1}\right) \in Y_{E} \uplus Y_{\infty}$ and that $h\left(x_{2}\right)=f\left(x_{2}\right) \in Y_{O}$. Because $Y_{E} \uplus Y_{\infty}$ and $Y_{O}$ have no elements in common, it follows that $h\left(x_{1}\right) \neq h\left(x_{2}\right)$. We have proved that $h$ is injective.
Case 4: $x_{2} \in X_{O} \biguplus X_{\infty}$ and $x_{1} \in X_{E}$. Injectivity of $h$ follows from case 3 because we can switch the roles of $x_{1}$ and $x_{2}$.

We finally show that $h$ is surjective: Let $y \in Y$. There are two cases.
Case 1: $y \in Y_{E} \biguplus Y_{\infty}$. It follows from (7.18) that $g(y) \in X_{O} \biguplus X_{\infty}$, hence

$$
h(g(y))=\sigma(g(y))=g^{-1}(g(y))=y .
$$

Here the second equation follows from (7.18). We have found an item in $X$ which is mapped by $h$ to $y$. and this proves that $h$ is surjective.
Case 2: $y \in Y_{O}$. It follows from (7.18) that $f^{-1}(y) \in X_{E}$, hence $h\left(f^{-1}(y)\right)=f\left(f^{-1}(y)\right)=y$. Again we have found an item in $X$ which is mapped by $h$ to $y$. We have proved that $h$ is surjective also in this case.

Corollary 7.7. The relation $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ on the set $\mathscr{A}:=\{\operatorname{card}(X): X \subseteq \Omega\}$ is a partial ordering.

Proof if none of the sets involved is empty: Reflexivity is obvious, antisymmetry follows from Cantor-Schröder-Bernstein and transitivity follows from prop.4.1.a: The composition of two injective functions is injective.
Proof if at least one of the sets involved is empty: Use the fact that if $\operatorname{card}(A) \leqq \operatorname{card}(B)$ and $B=\emptyset$ then $A=\emptyset$.

As an example we prove antisymmetry: If $X=\emptyset$ and $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ and $\operatorname{card}(Y) \leqq \operatorname{card}(X)$ then the second " $\leqq$ " implies $Y=\emptyset$, i.e., $X=Y$. On the other hand, If $Y=\emptyset$ and $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ and $\operatorname{card}(Y) \leqq \operatorname{card}(X)$ then the first "§" implies $X=\emptyset$, i.e., $X=Y$.

Theorem 7.6. Let $X, Y \subseteq \Omega$. Then $\operatorname{card}(X) \leqq \operatorname{card}(Y)$ or $\operatorname{card}(Y) \leqq \operatorname{card}(X)$.
In other words, " $\leqq$ " is a total ordering on the set of all cardinalities for subsets of $\Omega$.
PROOF: The proof will be given in thm.12.2, p.261, of ch. 12 (Applications of Zorn's Lemma).

### 7.2 Exercises for Ch. 7

Exercise 7.1. Prove prop.7.9 on p.126:
Let $X, Y$ be two sets such that $\operatorname{card}(X)=\operatorname{card}(Y)$. Then $\operatorname{card}\left(2^{X}\right)=\operatorname{card}\left(2^{Y}\right)$.

Exercise 7.2. Prove prop.7.8 on p.126:
Let $n \in \mathbb{Z}_{\geq 0}$. Let $X$ be a set such that $|X|=n$. Then its powerset has size $\left|2^{X}\right|=2^{n}$.
Exercise 7.3. Prove cor.7.4 on p.121: Let the set $X$ be uncountable and let $A \subseteq X$ be countable. Then the complement $A^{\complement}$ of $A$ is uncountable.

Exercise 7.4. Prove prop.7.7 on p.123: The relation $X \sim Y$ if and only if $X=Y$ or both $X$ and $Y$ are nonempty and there is a bijective function $f: X \rightarrow Y$ is an equivalence relation on $2^{\Omega}$.

### 7.3 Addenda to Ch. 7

Proposition 7.8. Let $n \in \mathbb{Z}_{\geq 0}$. Let $X$ be a set such that $|X|=n$. Then its powerset has size $\left|2^{X}\right|=2^{n}$.
The proof is left as exercise 7.2.
Proposition 7.9. Let $X, Y$ be two sets such that $\operatorname{card}(X)=\operatorname{card}(Y)$. Then $\operatorname{card}\left(2^{X}\right)=\operatorname{card}\left(2^{Y}\right)$.
The proof is left as exercise 7.1.
Lemma 7.2. Let $X$ be an arbitrary set and let $A$ be a nonempty proper subset of $X$. so that $X=A \biguplus A^{\complement}$ is a partitioning of $X$ into two nonempty subsets $A$ and $A^{\complement}$. Let $a \in A, a_{0} \in A^{\complement}$ and $A^{\prime}:=(A \backslash\{a\}) \biguplus\left\{a_{0}\right\}$. Then the function

$$
\varphi: A^{\prime} \xrightarrow{\sim} A ; \quad x \mapsto \begin{cases}x & \text { if } x \neq a_{0}, \\ a & \text { if } x=a_{0}\end{cases}
$$

is a bijection.
PROOF: The function

$$
\psi: A \xrightarrow{\sim} A^{\prime} ; \quad x \mapsto \begin{cases}x & \text { if } x \neq a, \\ a_{0} & \text { if } x=a\end{cases}
$$

is the inverse of $\varphi$.
The following theorem belongs before prop.7.2 on p.117.
Theorem 7.7 ( $\mathrm{B} / \mathrm{G}$ thm.13.31). Let $X$ be a set. Then $\operatorname{card}(X)<\operatorname{card}\left(2^{X}\right)$.
In other words, $X$ can be injected into $2^{X}$, but it is not possible to find bijective $f: X \xrightarrow{\sim} 2^{X}$.
Proof: We note that the function $x \mapsto\{x\}$ is an injection from $X$ into $2^{X}$. It follows that card $(X) \leqq$ $\operatorname{card}\left(2^{X}\right)$.

It remains to show that there is no bijective function $f: X \xrightarrow{\sim} 2^{X}$. We will show that it is not even possible to find surjective $f$ with domain $X$ and codomain $2^{X}$.

We assume to the contrary that such $f$ exists. Let

$$
\Gamma:=\{x \in X: x \notin f(x)\} .
$$

Obviously $\Gamma \subseteq X$, i.e., $\Gamma \in 2^{X} . f$ is surjective, hence there exists $x_{0} \in X$ such that $f\left(x_{0}\right)=\Gamma$.
Case 1: Assume $x_{0} \in \Gamma$. Then $x_{0} \notin f\left(x_{0}\right)$, i.e., $x_{0} \notin \Gamma$. We have a contradiction.
Case 2: Assume $x_{0} \notin \Gamma$. Then $x_{0} \in f\left(x_{0}\right)$, i.e., $x_{0} \in \Gamma$. Again, we have a contradiction.
We conclude that there is no surjective

The following theorem belongs after thm. 7.4 on $p .122(\mathbb{R}$ is countable $)$.
Theorem 7.8. Let $a, b \in \mathbb{R}$ such that $a<b$. Let $A$ be one of $] a, b[] a, b,],[a, b[,[a, b]$. Then $\operatorname{card}(A)=$ $\operatorname{card}(\mathbb{R})$.
PROOF:
a. $F: \mathbb{R} \rightarrow]-1,1\left[; \quad x \mapsto \frac{x}{|x|+1}\right.$ has the function $G(y):=\left\{\begin{array}{ll}\frac{y}{1-y} & \text { if } x \geqq 0, \\ \frac{y}{1+y} & \text { if } x<0 .\end{array}\right.$ as an inverse. The proof of this is tedious but elementary if one observes that $y \geqq 0$ if and only if $x \geqq 0$. This makes it easy to solve $y=\frac{x}{|x|+1}$ for $x$. The details are left to the reader. It follows that $\operatorname{card}(]-1,1[)=\operatorname{card}(\mathbb{R})$.
b. Let $c>0$. The function $x \mapsto c x$ is a bijection from $]-1,1[$ to $]-c, c\left[\right.$ because it has the function $y \mapsto \frac{y}{c}$ as an inverse. It follows from part a that $\operatorname{card}(]-c, c[)=\operatorname{card}(\mathbb{R})$.
c. If $\lambda \in \mathbb{R}$ then $x \mapsto x+\lambda$ bijects $A$ to $A+\lambda$ (Inverse: $y \mapsto y-\lambda$ ). Let $c:=\frac{b-a}{2}$ and $B:=A-\frac{a+b}{2}$. Then $B$ is the interval with endpoints $-c$ and $c$ where either endpoint of $B$ is included/excluded if and only if such is the case for the corresponding endpoint of $A$. Note that $\operatorname{card}(B)=\operatorname{card}(A)$ because $B$ is the image of $A$ under the bijection $x \mapsto x-\frac{b-a}{2}$.
d. It follows from $]-c, c[\subseteq B \subseteq]-2 c, 2 c[$ and part $\boldsymbol{b}$ that

$$
\operatorname{card}(\mathbb{R})=\operatorname{card}(]-c, c[) \leqq \operatorname{card}(B) \leqq \operatorname{card}(]-2 c, 2 c[)=\operatorname{card}(\mathbb{R})
$$

It is a consequence of the Cantor-Schröder-Bernstein Theorem wich is formulated and proved later on (ch.7.1.1 ${ }^{64}$, thm.7.5 on p.124) that there exists a bijection between $B$ and $\mathbb{R}$. We saw in part $c$ that card $(B)$ $=\operatorname{card}(A)$. This proves $\operatorname{card}(A)=\operatorname{card}(\mathbb{R})$.

The following corollary belongs after cor.7.5 (The set $Z$ of all integers is countable) on p.121.

Corollary 7.8. The rational numbers are countable.
PROOF: Let $n \in \mathbb{N}$ and $Q_{n}:=\left\{\frac{m}{n}: m \in \mathbb{Z}\right\}$. Then $f_{n}: Q_{n} \rightarrow \mathbb{Z}, \frac{m}{n} \mapsto m$ is a bijection because it has as an inverse the function $m \mapsto \frac{m}{n}$. It follows from cor.7.5 that $Q_{n}$ is countable. By thm.7.2, $\mathbb{Q}=\bigcup\left[Q_{n}: n \in \mathbb{N}\right]$ is countable as the union of countably many sets.

The next two theorems belong after prop.7.5 on p.123: The set of all algebraic numbers is countable.

Proposition 7.10. Let $k, m, n \in \mathbb{N}$. Then $\sqrt[k]{\frac{m}{n}}$ is algebraic.
PROOF: Let $p(x):=n x^{k}-m$. Then

$$
p(\sqrt[k]{m / n})=n(\sqrt[k]{m / n})^{k}-m=\frac{n m}{n}-m=0
$$

Proposition 7.11. Let $r \in \mathbb{Q}$. Then $r$ is algebraic.
PROOF: $r \in \mathbb{Q}$. There exist $m, n \in \mathbb{Z}$ such that $n \neq 0$ and $r=\frac{m}{b}$. Let $p(x):=n x-m$. Then

$$
p(m / n)=n \cdot \frac{m}{n}-m=0 .
$$

Note that if $m, n>0$ then the above is a special case of prop.7.10 (let $k:=1$ ).

[^49]
## 8 The Completeness of the Real Numbers System (Study this!)

### 8.1 Minima, Maxima, Infima and Suprema

Definition 8.1 (Upper and lower bounds, maxima and minima). ${ }^{65}$ Let $A \subseteq \mathbb{R}$. Let $l, u \in \mathbb{R}$. We call $l$ a lower bound of $A$ if $l \leqq a$ for all $a \in A$. We call $u$ an upper bound of $A$ if $u \geqq a$ for all $a \in A$.

We call $A$ bounded above if this set has an upper bound and we call $A$ bounded below if $A$ has a lower bound. We call $A$ bounded if $A$ is both bounded above and bounded below.

A minimum (min) of $A$ is a lower bound $l$ of $A$ such that $l \in A$. A maximum (max) of $A$ is an upper bound $u$ of $A$ such that $u \in A$.

The next proposition shows that min and max are unique if they exist. This makes it possible to write $\min (A)$ or $\min A$ for the minimum of $A$ and $\max (A)$ or $\max A$ for the maximum of $A$.

Proposition 8.1. Let $A \subseteq \mathbb{R}$. If $A$ has a maximum then it is unique. If $A$ has a minimum then it is unique.
Proof for maxima: Let $u_{1}$ and $u_{2}$ be two maxima of $A$ : both are upper bounds of $A$ and both belong to $A$. As $u_{1}$ is an upper bound, it follows that $a \leqq u_{1}$ for all $a \in A$. Hence $u_{2} \leqq u_{1}$. As $u_{2}$ is an upper bound, it follows that $u_{1} \leqq u_{2}$ and we have equality $u_{1}=u_{2}$. The proof for minima is similar.

Definition 8.2. Let $A \subseteq \mathbb{R}$. We define

$$
\begin{align*}
& A_{\text {low } 6}:=\{l \in \mathbb{R}: l \text { is lower bound of } A\} \\
& A_{u p p}:=\{u \in \mathbb{R}: u \text { is upper bound of } A\} . \tag{8.1}
\end{align*}
$$

Remark 8.1. Note that $A$ is bounded above if and only if $A_{\text {upp } 6} \neq \emptyset$ and bounded below if and only if $A_{\text {low } 6} \neq \emptyset$.

Axiom 8.1. (see [1] B/G axiom 8.52, p.83).
Completeness axiom for $\mathbb{R}$ : Let $A \subseteq \mathbb{R}$. If its set of upper bounds $A_{\text {upp } 6}$ is not empty then $A_{\text {upp } 6}$ has a minimum.

The above has the status of an axiom due to the fact that the real numbers usually are given axiomatically as an "archimedian ordered field" which satisfies the completeness axiom just stated.

Remark 8.2. $A_{\text {fow } 6}$ and/or $A_{\text {upp } 6}$ may be empty. Examples are $A=\mathbb{R}, A=\mathbb{R}_{>0}, A=\mathbb{R}_{<0}$.
Definition 8.3. Let $A \subseteq \mathbb{R}$. If $A_{\text {upp }}$ is not empty then $\min \left(A_{u p p 6}\right)$ exists by axiom 8.1 and it is unique by prop. 8.1. We write $\sup (A)$ or l.u.b. $(A)$ for $\min \left(A_{u p p b}\right)$ and call this element of $\mathbb{R}$ the supremum or least upper bound of $A$.

We will see in cor.8.1 that, if $A_{\text {low } 6}$ is not empty, then $\max \left(A_{\text {fow } 6}\right)$ exists and is unique by prop. 8.1. We write $\inf (A)$ or g.l.b. $(A)$ for $\max \left(A_{\text {lowb }}\right)$ and call this number the infimum or greatest lower bound of $A$.

[^50]Proposition 8.2 (Duality of upper and lower bounds, min and max, inf and sup). Let $A \subseteq \mathbb{R}$ and $x \in \mathbb{R}$. Then the following is true for $-x$ and $-A=\{-y: y \in A\}$ :

$$
\begin{align*}
& -x \text { is a lower bound of } A \Leftrightarrow x \text { is an upper bound of }-A, \\
& -x=\sup (A) \Leftrightarrow x=\inf (-A),  \tag{8.2}\\
& -x=\max (A) \Leftrightarrow x=\min (-A) .
\end{align*}
$$

We switch the roles of $x$ and $-x$ and obtain three more equations:

$$
\begin{align*}
& -x \text { is an upper bound of } A \Leftrightarrow x \text { is a lower bound of }-A, \\
& -x=\inf (A) \Leftrightarrow x=\sup (-A),  \tag{8.3}\\
& -x=\min (A) \Leftrightarrow x=\max (-A) .
\end{align*}
$$

PROOF: A simple consequence of

$$
-x \leq y \Leftrightarrow x \geqq-y \text { and }-x \geqq y \Leftrightarrow x \leqq-y .
$$

Corollary 8.1. Let $A \subseteq \mathbb{R}$. If $A$ has lower bounds then $\inf (A)$ exists.
PROOF: According to the duality proposition prop.8.2, if A has lower bounds then $(-A)$ has upper bounds. It follows from the completeness axiom that $\sup (-A)$ exists. We apply once more prop.8.2 to prove that $\inf (A)$ exists: $\inf (A)=-\sup (-A)$.

Here are some examples. We define for all three of them $f(x):=-x$ and $g(x):=x$.
Example 8.1 (Example a: Maximum exists). Let $X_{1}:=\{t \in \mathbb{R}: 0 \leqq t \leqq 1\}$.
For each $x \in X_{1}$ we have $|f(x)-g(x)|=g(x)-f(x)=2 x$ and the biggest possible such difference is $g(1)-f(1)=2$. So $\max \left(X_{1}\right)$ exists and equals $\max \left(X_{1}\right)=2$.

Example 8.2 (Example b: Supremum is finite). Let $X_{2}:=\{t \in \mathbb{R}: 0 \leqq t<1\}$, i.e., we now exclude the right end point 1 at which the maximum difference was attained. For each $x \in X$ we have

$$
|f(x)-g(x)|=g(x)-f(x)=2 x
$$

and the biggest possible such difference is certainly bigger than

$$
g(0.9999999999)-f(0.9999999999)=1.9999999998 .
$$

If you keep adding $5,0009 \mathrm{~s}$ to the right of the argument $x$, then you get the same amount of 9 s inserted into the result $2 x$, so $2 x$ comes closer than anything you can imagine to the number 2 , without actually being allowed to reach it.

The supremum is still considered in a case like this to be 2 . This precisely is the difference in behavior between the supremum $s:=\sup (A)$ and the maximum $m:=\max (A)$ of a set $A \subseteq \mathbb{R}$ of real numbers: For the maximum there must actually be at least one element $a \in A$ so that $a=\max (A)$.

For the supremum it is sufficient that there is a sequence $a_{1} \leqq a_{2} \leqq \ldots$ which approximates $s$ from below in the sense that the difference $s-a_{n}$ "drops down to zero" as $n$ approaches infinity. We will not be more exact than this because doing so would require us to delve into the concepts of convergence and contact points.

Example 8.3 (Example c: Supremum is infinite). Let $X_{3}:=\mathbb{R}_{\geq 0}=\{t \in \mathbb{R}: 0 \leqq t\}$. For each $x \in X_{3}$ we have again $|f(x)-g(x)|=g(x)-f(x)=2 x$. But there is no more limit to the right for the values of $x$. The difference $2 x$ will exceed all bounds and that means that the only reasonable value for $\sup \left\{|f(x)-g(x)|: x \in X_{3}\right\}$ is $+\infty$.

As in case b above, the max does not exist because there is no $x_{0} \in X_{3}$ such that $\left|f\left(x_{0}\right)-g\left(x_{0}\right)\right|$ attains the highest possible value among all $x \in X_{3}$.

You should understand that even though $\sup (A)$ as best approximation of the largest value of $A \subseteq \mathbb{R}$ is allowed to take the "value" $+\infty$ or $-\infty$ this cannot be allowed for $\max (A)$.

How so? The infinity values are not real numbers, but, by definition of the maximum, if $\alpha:=\max (A)$ exists, then $\alpha \in A$. In particular, the max must be a real number.

That last example motivates the following definition.
Definition 8.4 (Supremum and Infimum of unbounded and empty sets). If $A$ is not bounded from above then we define

$$
\begin{equation*}
\sup A=\infty \tag{8.4}
\end{equation*}
$$

If $A$ is not bounded from below then we define

$$
\begin{equation*}
\inf A=-\infty \tag{8.5}
\end{equation*}
$$

Finally we define

$$
\begin{equation*}
\sup \emptyset=-\infty, \quad \inf \emptyset=+\infty \tag{8.6}
\end{equation*}
$$

Proposition 8.3. Let $A \subseteq B \subseteq \mathbb{R}$.
a. Then $\inf (A) \geqq \inf (B)$ and $\sup (A) \leqq \sup (B)$.
b. If both $A$ and $B$ have a min then $\min (A) \geqq \min (B)$. If both $A$ and $B$ have a max then $\max (A) \leqq$ $\max (B)$.

## PROOF:

$\boldsymbol{b}$ follows from $\boldsymbol{a}$ because if a set has a minimum then it equals its infimum and if a set has a maximum then it equals its supremum

We now prove a for suprema. The proof for infima is similar. First we note that if at least one of the sets is empty or not bounded above then the proof is trivial. ${ }^{66}$ We may assume that both $A$ and $B$ are not empty and bounded above.

It follows from $A \subseteq B$ that any upper bound of $B$ also is an upper bound of $A$.
In particular $\sup (B) \in A_{\text {uppb }}$, hence $\sup (A)=\min \left(A_{\text {upp }}\right) \leqq \sup (B)$.
Definition 8.5 (Translation and dilation of sets of real numbers). Let $A \subseteq \mathbb{R}^{67}$ and $\alpha, b \in \mathbb{R}$. We define

$$
\begin{equation*}
\lambda A+b:=\{\lambda a+b: a \in A\} . \tag{8.7}
\end{equation*}
$$

[^51]In particular, for $\lambda= \pm 1$, we obtain

$$
\begin{align*}
A+b & =\{a+b: a \in A\},  \tag{8.8}\\
-A & =\{-a: a \in A\} . \tag{8.9}
\end{align*}
$$

Proposition 8.4 (The sup of a set is positively homogeneous). Let $A$ be a nonempty subset of $\mathbb{R}$ and let $\lambda \in \mathbb{R}_{\geq 0}$. If $\lambda>0$ or if $\lambda=0$ and $\sup (A)<\infty$ then

$$
\begin{equation*}
\sup (\lambda A)=\lambda \sup (A) \tag{8.10}
\end{equation*}
$$

PROOF: (8.10) holds for $\lambda=0$ and $\sup (A)<\infty$ because

$$
\sup (0 A)=\sup (\{0\})=0=0 \cdot \sup (A)
$$

Let $\lambda>0$. If $B \subseteq \mathbb{R}$, let $B_{u p p}:=\{u \in \mathbb{R}: u$ is upper bound of $B\}$. Note that, because $\lambda>0$,

$$
\begin{equation*}
u \in A_{u p p 6} \Leftrightarrow u \geqq a \forall a \in A \Leftrightarrow \lambda u \geqq \lambda a \forall a \in A \Leftrightarrow \lambda u \in(\lambda A)_{u p p b} . \tag{8.11}
\end{equation*}
$$

It follows from $\sup (A) \in A_{\text {upp }}$ that $\lambda \sup (A) \in(\lambda A)_{\text {upp }}$, hence $\lambda \sup (A) \geqq \min \left((\lambda A)_{\text {uppb }}\right)=\sup (\lambda A)$.
It remains to show that $\lambda \sup (A) \leqq \sup (\lambda A)$.
We substitute $\frac{v}{\lambda}$ for $u$ in (8.11) and obtain $\frac{v}{\lambda} \in A_{\text {upp }} \Leftrightarrow v \in(\lambda A)_{\text {upp }}$.
It follows from $\sup (\lambda A) \in(\lambda A)_{\text {upp }}$ that $\frac{\sup (\lambda A)}{\lambda} \in A_{\text {upp } 6,}$, hence $\frac{\sup (\lambda A)}{\lambda} \geqq \min \left(A_{\text {upp }}\right)=\sup (A)$. This proves $\lambda \sup (A) \leqq \sup (\lambda A)$.
Definition 8.6 (bounded functions). Given is a nonempty set $X$. A real-valued function $f(\cdot)$ with domain $X$ is called bounded from above if there exists a (possibly very large) number $\gamma_{1}>0$ such that

$$
\begin{equation*}
f(x)<\gamma_{1} \quad \text { for all arguments } x \text {. } \tag{8.12}
\end{equation*}
$$

It is called bounded from below if there exists a (possibly very large) number $\gamma_{2}>0$ such that

$$
\begin{equation*}
f(x) \quad>-\gamma_{2} \quad \text { for all arguments } x \tag{8.13}
\end{equation*}
$$

It is called a bounded function if it is both bounded from above and below. It is obvious that if you set $\gamma:=\max \left(\gamma_{1}, \gamma_{2}\right)$ then bounded functions are exactly those that satisfy the inequality

$$
\begin{equation*}
|f(x)|<\gamma \quad \text { for all arguments } x \text {. } \tag{8.14}
\end{equation*}
$$

We note that $f$ is bounded if and only if its range $f(X)$ is a bounded subset of $\mathbb{R}$. We further note that we have defined infimum and supremum for any kind of set: empty or not, bounded above or below or not. We use those definitions to define infimum and supremum for functions, sequences and indexed families.

Definition 8.7 (supremum and infimum of functions). Let $X$ be an arbitrary set, $A \subseteq X$ a subset of $X, f: X \rightarrow \mathbb{R}$ a real valued function on $X$. Look at the set $f(A)=\{f(x): x \in A\}$, i.e., the image of $A$ under $f(\cdot)$.

The supremum of $f(\cdot)$ on $A$ is then defined as

$$
\begin{equation*}
\sup _{A} f:=\sup _{x \in A} f(x):=\sup f(A) \tag{8.15}
\end{equation*}
$$

The infimum of $f(\cdot)$ on $A$ is then defined as

$$
\begin{equation*}
\inf _{A} f:=\inf _{x \in A} f(x):=\inf f(A) \tag{8.16}
\end{equation*}
$$

Definition 8.8 (supremum and infimum of families). Let $\left(x_{i}\right)_{i \in I}$ be an indexed family of real numbers $x_{i}$.

The supremum of $\left(x_{i}\right)_{i \in I}$ is then defined as

$$
\begin{equation*}
\sup \left(x_{i}\right):=\sup _{i}\left(x_{i}\right):=\sup \left(x_{i}\right)_{i}:=\sup \left(x_{i}\right)_{i \in I}:=\sup _{i \in I} x_{i}:=\sup \left\{x_{i}: i \in I\right\} \tag{8.17}
\end{equation*}
$$

The infimum of $\left(x_{i}\right)_{i \in I}$ is then defined as

$$
\begin{equation*}
\inf \left(x_{i}\right):=\inf _{i}\left(x_{i}\right):=\inf \left(x_{i}\right)_{i}:=\inf \left(x_{i}\right)_{i \in I}:=\inf _{i \in I} x_{i}:=\inf \left\{x_{i}: i \in I\right\} . \tag{8.18}
\end{equation*}
$$

The definition above for families extends to sequences. Let $I=\left\{k \in \mathbb{Z}: k \geq k_{0}\right.$ for some $k_{0} \in \mathbb{Z}$.
Definition 8.9 (supremum and infimum of sequences). Let $\left(x_{n}\right)_{n \in I}$ be a sequence of real numbers $x_{n}$. The supremum of $\left(x_{n}\right)_{n \in I}$ is then defined as

$$
\begin{equation*}
\sup \left(x_{n}\right):=\sup \left(x_{n}\right)_{n \in I}:=\sup _{n \in I} x_{n}=\sup \left\{x_{n}: n \in I\right\} \tag{8.19}
\end{equation*}
$$

The infimum of $\left(x_{n}\right)_{n \in I}$ is then defined as

$$
\begin{equation*}
\inf \left(x_{n}\right):=\inf \left(x_{n}\right)_{n \in I}:=\inf _{n \in I} x_{n}=\inf \left\{x_{n}: n \in I\right\} \tag{8.20}
\end{equation*}
$$

We note that the "duality principle" for min and max, sup and inf (8.2 on p.129) is true in all cases above: You flip the sign of the items you examine and the sup/max of one becomes the inf/min of the other and vice versa.

Proposition 8.5. $X$ be a nonempty set and $\varphi, \psi: X \rightarrow \mathbb{R}$ be two real valued functions on $X$. Let $A \subseteq X$. Then

$$
\begin{align*}
\sup \{\varphi(x)+\psi(x): x \in A\} & \leqq \sup \{\varphi(y): y \in A\}+\sup \{\psi(z): z \in A\},  \tag{8.21}\\
\inf \{\varphi(x)+\psi(x): x \in A\} & \geqq \inf \{\varphi(y): y \in A\}+\inf \{\psi(z): z \in A\} . \tag{8.22}
\end{align*}
$$

PROOF:
We only prove (8.21). The proof of (8.22) is similar. ${ }^{68}$
Let $U:=\{\varphi(x)+\psi(x): x \in A\}, V:=\{\varphi(y): y \in A\}, W:=\{\psi(z): z \in A\}$. Let $x \in A$.
Then $\sup (V)$ is an upper bound of $\varphi(x)$ and $\sup (W)$ is an upper bound of $\psi(x)$,
hence $\sup (V)+\sup (W) \geqq \varphi(x)+\psi(x)$.
This is true for all $x \in A$, hence $\sup (V)+\sup (W)$ is an upper bound of $U$.
It follows that $\sup (V)+\sup (W)$ dominates the least upper bound $\sup (U)$ of $U$ and this proves (8.21).

### 8.2 Convergence and Continuity in $\mathbb{R}$

You are familiar with the concepts of convergent sequences and continuous functions whose domain and codomain both are sets of real numbers from calculus. We discuss them here in a more rigorous fashion. Convergence and continuity will be generalized in later chapters from $\mathbb{R}$ to so-called metric spaces.

Definition 8.10 (convergence of sequences of real numbers). We say that a sequence $\left(x_{n}\right)$ of real numbers converges ${ }^{69}$ to $a \in \mathbb{R}$ for $n \rightarrow \infty$ if almost all of the $x_{n}$ will come arbitrarily close to $a$ in the following sense:

For any $\delta \in \mathbb{R}_{>0}$ (no matter how small) there exists (possibly extremely large) $n_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\left|a-x_{j}\right|<\delta \text { for all } j \geqq n_{0} . \tag{8.23}
\end{equation*}
$$

We write either of

$$
\begin{equation*}
a=\lim _{n \rightarrow \infty} x_{n} \quad \text { or } \quad x_{n} \rightarrow a \tag{8.24}
\end{equation*}
$$

and we call $a$ the limit of the sequence $\left(x_{n}\right)$.

There is an equivalent way of expressing convergence towards $a \in \mathbb{R}$ : No matter how small a "neighborhood" ] $a-\delta, a+\delta\left[\right.$ of a you choose: at most finitely many of the $x_{n}$ will be located outside that neighborhood.

Example 8.4. Some simple examples for convergence to a real number:
a. Let $x_{n}:=1 / n(n \in \mathbb{N})$. Then $x_{n} \rightarrow 0$ as $n \rightarrow \infty$.
b. Let $\alpha \in \mathbb{R}$ and $z_{n}:=\alpha^{2} \pi(n \in \mathbb{N})$. Then the sequence $\left(z_{n}\right)_{n}$ has limit $\alpha^{2} \pi$.
c. More generally let $z_{n}:=x_{0}$ for some $x_{0} \in \mathbb{R}(n \in \mathbb{N})$. Then $\lim _{n \rightarrow \infty} z_{n}=x_{0}$.

PROOF of $\boldsymbol{a}$ : If $\delta>0$, let $n_{0}:=$ some integer larger than $1 / \delta$. Such a number exists because the natural numbers are not bounded above. It follows for $n \geqq n_{0}$ that

$$
\left|x_{n}-0\right|=1 / n \leqq 1 / n_{0}<\delta .
$$

[^52]PROOF of $\boldsymbol{b}$ and $\boldsymbol{c}$ : Left as an exercise.
Convergence is an extremely important concept in Mathematics, but it excludes the case of sequences such as $x_{n}:=n$ and $y_{n}:=-n(n \in \mathbb{N})$. Intuition tells us that $x_{n}$ converges to $\infty$ and $y_{n}$ converges to $-\infty$ because we think of very big numbers as being very close to $+\infty$ and very small numbers (i.e., very big ones with a minus sign) as being very close to $-\infty$.

Definition 8.11 (Limit infinity). Given a real number $K>0$, we define

$$
\begin{align*}
N_{K}(\infty) & :=\{x \in \mathbb{R}: x>K\}  \tag{8.25a}\\
N_{K}(-\infty) & :=\{x \in \mathbb{R}: x<-K\} \tag{8.25b}
\end{align*}
$$

We call $N_{K}(\infty)$ the K-neighborhood of $\infty$ and $N_{K}(-\infty)$ the K-neighborhood of $-\infty$. We say that a sequence $\left(x_{n}\right)$ has limit $\infty$ and we write either of

$$
\begin{equation*}
x_{n} \rightarrow \infty \quad \text { or } \quad \lim _{n \rightarrow \infty} x_{n}=\infty \tag{8.26}
\end{equation*}
$$

if the following is true for any (big) $K$ : There is a (possibly extremely large) integer $n_{0}$ such that all $x_{j}$ belong to $N_{K}(\infty)$ just as long as $j \geqq n_{0}$.
We say that the sequence $\left(x_{n}\right)$ has limit $-\infty$ and we write either of

$$
\begin{equation*}
x_{n} \rightarrow-\infty \quad \text { or } \quad \lim _{n \rightarrow \infty} x_{n}=-\infty \tag{8.27}
\end{equation*}
$$

if the following is true for any (big) $K$ : There is a (possibly extremely large) integer $n_{0}$ such that all $x_{j}$ belong to $N_{K}(-\infty)$ just as long as $j \geqq n_{0}$.
$\boldsymbol{a}$. There is an equivalent way of stating that the sequence $\left(x_{n}\right)$ has limit $\infty$ : No matter how big a threshold $K>0$ you choose: at most finitely many of the $x_{n}$ will be located above that threshold.
b. $x_{n} \rightarrow-\infty$ can also be expressed as follows: No matter how big a threshold $K>0$ you choose: at most finitely many of the $x_{n}$ will be located below $-K$.

## REMARK:

The majority of mathematicians agrees that there is no "convergence to $\infty$ " or "divergence to $\infty$ ". Rather, they will state that a sequence has the limit $\infty$. We follow this convention

Note 8.1 (Notation for limits of monotone sequences). Let $\left(x_{n}\right)$ be a non-decreasing sequence of real numbers and let $y_{n}$ be a non-increasing sequence. If $\xi=\lim _{k \rightarrow \infty} x_{k}$ (that limit might be $+\infty$ ) then we write suggestively

$$
x_{n} \nearrow \xi \quad(n \rightarrow \infty)
$$

If $\eta=\lim _{j \rightarrow \infty} y_{j}$ (that limit might be $-\infty$ ) then we write suggestively

$$
y_{j} \searrow \eta \quad(j \rightarrow \infty)
$$

The following proposition shows that the limit behavior of a sequence is a property of its tail, i.e., it does not depend on the first finitely many indices.
Proposition 8.6. Let $x_{n}, y_{n} \in \mathbb{R}$ be two sequences of real numbers. Assume there is $K \in \mathbb{N}$ such that $x_{n}=y_{n}$ for all $n \geqq K$. Let $L \in \mathbb{R}$. Then

$$
\lim _{n \rightarrow \infty} x_{n}=L \Leftrightarrow \lim _{n \rightarrow \infty} y_{n}=L, \quad \lim _{n \rightarrow \infty} x_{n}= \pm \infty \Leftrightarrow \lim _{n \rightarrow \infty} y_{n}= \pm \infty .
$$

PROOF:
Case 1: $\lim _{n \rightarrow \infty} x_{n}=L$.
Let $\delta>0$. Then there exists $N^{\prime} \in \mathbb{N}$ such that

$$
\begin{equation*}
\left|x_{j}-L\right|<\delta \text { for all } j \geqq N^{\prime} . \tag{8.28}
\end{equation*}
$$

Let $N:=\max \left(K, N^{\prime}\right)$. It follows from $N \geqq K, N \geqq N^{\prime} x_{n}=y_{n}$ for all $n \geqq K$ and (8.28) that

$$
\begin{equation*}
\left|y_{j}-L\right|=\left|x_{j}-L\right|<\delta \text { for all } j \geqq N^{\prime} . \tag{8.29}
\end{equation*}
$$

This proves that $y_{n}$ converges to $L$.
Case 2: $\lim _{n \rightarrow \infty} x_{n}=\infty$.
Let $M \in \mathbb{R}$. It follows from $x_{n} \rightarrow \infty$ and def.8.11 that there exists $N^{\prime} \in \mathbb{N}$ such that

$$
\begin{equation*}
x_{j}>M \text { for all } j \geqq N^{\prime} \tag{8.30}
\end{equation*}
$$

Let $N:=\max \left(K, N^{\prime}\right)$. It follows from $N \geqq K, N \geqq N^{\prime} x_{n}=y_{n}$ for all $n \geqq K$ and (8.30) that

$$
\begin{equation*}
y_{j}=x_{j}>M \text { for all } j \geqq N^{\prime} \tag{8.31}
\end{equation*}
$$

This proves that $\lim _{n \rightarrow \infty} y_{n}=\infty$.
Case 3: $\lim _{n \rightarrow \infty} x_{n}=-\infty$.
This is true according to the already proven case 2, applied to the sequences $\left(-x_{n}\right)_{n}$ and $\left(-y_{n}\right)_{n}$.
Proposition 8.7. [See B/G prop.10.16]
Let $\left(x_{n}\right)_{n}$ be a sequence of real numbers such that $\lim _{n \rightarrow \infty} x_{n}$ exists. Let $K \in \mathbb{N}$. For $n \in \mathbb{N}$ let $y_{n}:=x_{n+K}$. Then $\left(y_{n}\right)_{n}$ has the same limit.

PROOF:
Case 1: $\lim _{n \rightarrow \infty} x_{n}=L \in \mathbb{R}$.
The proof is left as exercise 8.3.
Case 2: $\lim _{n \rightarrow \infty} x_{n}=\infty$.
Let $M \in \mathbb{R}$. It follows from $x_{n} \rightarrow \infty$ and def.8.11 that there exists $N \in \mathbb{N}$ such that

$$
\begin{equation*}
x_{j}>M \text { for all } j \geqq N . \tag{8.32}
\end{equation*}
$$

Let $j \geqq N$. Then $j+K \geqq N$ and it follows from (8.32) that

$$
\begin{equation*}
y_{j}=x_{j+K}>M \tag{8.33}
\end{equation*}
$$

This proves that $\lim _{n \rightarrow \infty} y_{n}=\infty$.
Case 3: $\lim _{n \rightarrow \infty} x_{n}=-\infty$.
This is true according to the already proven case 2, applied to the sequences $\left(-x_{n}\right)_{n}$ and $\left(-y_{n}\right)_{n}$.

Proposition 8.8 (convergent $\Rightarrow$ bounded). Let $\left(x_{n}\right)_{n}$ be a sequence in $\mathbb{R}$ with limit $a \in \mathbb{R}$. Then this sequence is bounded.

PROOF: We must prove that there exists $K \in \mathbb{R}$ such that $\left|x_{j}\right| \leqq K$ for all $j \in \mathbb{N}$.
Let $\delta=1$ in (8.23). It follows that there is $n_{0} \in \mathbb{N}$ such that $\left|x_{j}-a\right|<1$ for all $j \geqq n_{0}$. But then

$$
\left|x_{j}\right|=\left|\left(x_{j}-a\right)+a\right| \leqq\left|x_{j}-a\right|+|a| \leqq|a|+1 \text { for all } j \geqq n_{0} \text {. }
$$

Let $K:=\max \left(1,\left|x_{1}\right|,\left|x_{2}\right|, \ldots,\left|x_{n_{0}-1}\right|\right.$ It follows that $\left|x_{j}\right| \leqq K+a$ for all $j \in \mathbb{N}$.
The following proposition states that the product of a sequence which converges to zero and a bounded sequence converges to zero.

Proposition 8.9 (bounded times zero-convergent is zero-convergent). Let $\left(x_{n}\right)_{n}$ and $\left(\alpha_{n}\right)_{n}$ be two sequences in $\mathbb{R}$ and let $\alpha \in \mathbb{R}$. If $\lim _{n \rightarrow \infty} x_{n}=0$ and if $\left|\alpha_{j}\right| \leqq \alpha$ for all $j \in \mathbb{N}$ then

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left(\alpha_{j} x_{j}\right)=0 \tag{8.34}
\end{equation*}
$$

PROOF:
Case 1: $\alpha=0$. Then $\alpha_{j}=0$ for all $j \in \mathbb{N}$, hence $\alpha_{j} x_{j}=0$. For any $\delta>0$ let $n_{0}=1$. Then

$$
\left|\alpha_{j} x_{j}-0\right|=\left|\alpha_{j} x_{j}\right|=0<\delta \text { for all } j \in \mathbb{N} \text { such that } j \geqq n_{0} \text {. }
$$

This proves convergence $\alpha_{j} x_{j} \rightarrow 0$.
Case 2: $\alpha \neq 0$, i.e., $|\alpha|>0$. Let $\delta>0$. We must show that
there is $n_{0} \in \mathbb{N}$ such that $\left|\alpha_{j} x_{j}\right|<\delta$ for all $j \in \mathbb{N}$ such that $j \geqq n_{0}$.
Let $\varepsilon:=\delta /|\alpha|$. Then $\varepsilon>0$ and it follows from $\lim _{j \rightarrow \infty} x_{j}=0$ that

$$
\begin{equation*}
\text { there is } N \in \mathbb{N} \text { such that }\left|x_{j}\right|<\varepsilon \text { for all } j \in \mathbb{N} \text { such that } j \geqq N \text {. } \tag{8.36}
\end{equation*}
$$

But then $\left|\alpha_{j} x_{j}\right|=\left|\alpha_{j}\right| \cdot\left|x_{j}\right|<|\alpha| \cdot \varepsilon=\delta$ for all $j \in \mathbb{N}$ such that $j \geqq N$. We choose $n_{0}:=N$ and (8.35) follows.

It is very rare that you need to apply def. 8.10 on $p .133$ to compute a limit. Rather, the previous proposition and the following set of rules are employed.
Proposition 8.10 (Rules of arithmetic for limits ${ }^{70}$ ). Let $\left(x_{n}\right)_{n}$ and $\left(y_{n}\right)_{n}$ be two sequences in $\mathbb{R}$ and let $x, y, \alpha \in \mathbb{R}$.
Let $\lim _{j \rightarrow \infty} x_{j}=x$ and $\lim _{j \rightarrow \infty} y_{j}=y$. Then
a. $\lim _{j \rightarrow \infty} \alpha=\alpha$,
b. $\lim _{j \rightarrow \infty}\left(\alpha \cdot x_{j}\right)=\alpha \cdot x$,
c. $\lim _{j \rightarrow \infty}\left(x_{j}+y_{j}\right)=x+y$,
d. $\lim _{j \rightarrow \infty}\left(x_{j} \cdot y_{j}\right)=x \cdot y$,
e. if $x \neq 0$ then $\lim _{j \rightarrow \infty} \frac{1}{x_{j}}=\frac{1}{x}$.

[^53]PROOF of a: Exercise 8.1.
PROOF of $\boldsymbol{b}$ :
Case 1: $\alpha=0$. Then $\alpha x_{j}$ is the constant sequence $0,0, \ldots$ which converges to $0=\alpha x$ - Done.
Case 2: $\alpha \neq 0$. Let $\delta>0$. We must show that

$$
\begin{equation*}
\text { there is } n_{0} \in \mathbb{N} \text { such that }\left|\alpha x_{j}-\alpha x\right|<\delta \text { for all } j \in \mathbb{N} \text { such that } j \geqq n_{0} \text {. } \tag{8.37}
\end{equation*}
$$

Let $\varepsilon:=\delta /|\alpha|$. Then $\varepsilon>0$ and it follows from $\lim _{j \rightarrow \infty} x_{j}=x$ that

$$
\begin{equation*}
\text { there is } N \in \mathbb{N} \text { such that }\left|x_{j}-x\right|<\varepsilon \text { for all } j \in \mathbb{N} \text { such that } j \geqq N \text {. } \tag{8.38}
\end{equation*}
$$

But then $\left|\alpha x_{j}-\alpha x\right|=|\alpha| \cdot\left|x_{j}-x\right|<|\alpha| \cdot \varepsilon=\delta$ for all $j \in \mathbb{N}$ such that $j \geqq N$. We choose $n_{0}:=N$ and (8.37) is proved.
PROOF of $c$ :
Let $\delta>0$. It follows from $\lim _{j \rightarrow \infty} x_{j}=x$ and $\lim _{j \rightarrow \infty} y_{j}=y$ that there exist $N_{1}, N_{2} \in \mathbb{N}$ such that

$$
\begin{equation*}
\text { if } j \geqq N_{1} \text { then }\left|x_{j}-x\right|<\delta / 2 \text { and if } j \geqq N_{2} \text { then }\left|y_{j}-y\right|<\delta / 2 \text {. } \tag{8.39}
\end{equation*}
$$

It follows from the triangle inequality $|A+B| \leqq|A|+|B|$ (prop.2.3 on p.19) and from (8.39) that

$$
\begin{equation*}
\left.\left|\left(x_{j}+y_{j}\right)-(x+y)\right|=\left|\left(x_{j}-x\right)+\left(y_{j}-y\right)\right| \leqq\left|x_{j}-x\right|+\mid y_{j}-y\right) \mid<\delta / 2+\delta / 2=\delta \tag{8.40}
\end{equation*}
$$

for all $j \geqq \max \left(N_{1}, N_{2}\right)$. Let $n_{0}:=\max \left(N_{1}, N_{2}\right)$. It follows from (8.40) that $\left|\left(x_{j}+y_{j}\right)-(x+y)\right|<\delta$ for all $j \geqq n_{0}$. This proves $\boldsymbol{c}$.

## PROOF of $d$ :

Let $u_{j}:=\left(x_{j}-x\right) y_{j}$ and $v_{j}:=x\left(y_{j}-y\right)(j \in \mathbb{N})$. Then

$$
\begin{equation*}
x_{j} y_{j}-x y=\left(x_{j} y_{j}-x y_{j}\right)+\left(x y_{j}-x y\right)=\left(x_{j}-x\right) y_{j}+x\left(y_{j}-y\right)=u_{j}+v_{j} . \tag{8.41}
\end{equation*}
$$

The convergent sequences $\left(y_{n}\right)_{n}$ and $(x)_{n}$ (constant sequence!) are bounded by prop.8.8 and it follows from parts $c$ and $\boldsymbol{a}$ that

$$
\begin{aligned}
\lim _{j \rightarrow \infty}\left(x_{j}-x\right) & =\lim _{j \rightarrow \infty} x_{j}-\lim _{j \rightarrow \infty} x=x-x=0, \\
\lim _{j \rightarrow \infty}\left(y_{j}-y\right) & =\lim _{j \rightarrow \infty} y_{j}-\lim _{j \rightarrow \infty} y=y-y=0,
\end{aligned}
$$

i.e., the sequences $x_{j}-x$ and $y_{j}-y$ converge to zero.

It now follows from prop.8.9 that $\lim _{j \rightarrow \infty} u_{j}=0$ and $\lim _{j \rightarrow \infty} v_{j}=0$. According to (8.41) $x_{j} y_{j}=x y+$ $u_{j}+v_{j}$ is the limit of three convergent sequences. ${ }^{71}$ It follows from part $\boldsymbol{c}$ that

$$
\lim _{j \rightarrow \infty} x_{n} y_{n}=\lim _{j \rightarrow \infty}(x y)+\lim _{j \rightarrow \infty} u_{j}+\lim _{j \rightarrow \infty} v_{j}=x y+0+0=x y
$$

PROOF of $\boldsymbol{e}$ :
Because $\lim _{j \rightarrow \infty} x_{n}=x$ and $|x|>0$ there exists $N_{1} \in \mathbb{N}$ such that $\left|x_{n}-x\right| \leqq|x| / 2$ for all $j \geqq N_{1}$. Hence

$$
\begin{align*}
& |x|=\left|\left(x-x_{n}\right)+x_{n}\right| \leqq\left|x-x_{n}\right|+\left|x_{n}\right| \leqq \frac{|x|}{2}+\left|x_{n}\right| \\
\Rightarrow & \frac{|x|}{2} \leqq\left|x_{n}\right| \Rightarrow|x|\left|x_{n}\right| \geqq \frac{x^{2}}{2} \Rightarrow \frac{1}{\left|x x_{n}\right|} \leqq \frac{2}{x^{2}} . \tag{8.42}
\end{align*}
$$

[^54]Let $z_{n}:=\left(x x_{n}\right)^{-1}$ and $K:=\max \left(2 / x^{2},\left|z_{1}\right|,\left|z_{2}\right|, \ldots,\left|z_{N_{1}}\right|\right)$. It follows from (8.42) that the sequence $\left(z_{n}\right)_{n}$ is bounded by K. It follows from part a that $\lim _{j \rightarrow \infty} x=x$ and from part $\boldsymbol{c}$ that $\lim _{j \rightarrow \infty}\left(x_{n}-x\right)=\lim _{j \rightarrow \infty}\left(x_{n}\right)-x=0$.

It now follows from prop.8.9 that

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left(\frac{1}{x_{n}}-\frac{1}{x}\right)=\lim _{j \rightarrow \infty} \frac{1}{x x_{n}} \cdot\left(x-x_{n}\right)=\lim _{j \rightarrow \infty} z_{n}\left(x-x_{n}\right)=0 \tag{8.43}
\end{equation*}
$$

Let $\delta>0$. On account of (8.43) there exists $n_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\left|\frac{1}{x_{n}}-\frac{1}{x}\right|=\left|\left(\frac{1}{x_{n}}-\frac{1}{x}\right)-0\right|<\delta \text { for all } j \geqq n_{0} . \tag{8.44}
\end{equation*}
$$

This proves convergence of $1 / x_{n}$ to $1 / x$.

## Proposition 8.11.

a. Let $x_{n}$ be a sequence of real numbers that is non-decreasing, i.e., $x_{n} \leqq x_{n+1}$ for all $n$ (see def. 17.1 on $p .304$ ), and bounded above. Then $\lim _{n \rightarrow \infty} x_{n}$ exists and coincides with $\sup \left\{x_{n}: n \in \mathbb{N}\right\}$
b. Further, if $y_{n}$ is a sequence of real numbers that is non-increasing, i.e., $y_{n} \geqq y_{n+1}$ for all $n$, and bounded below, the analogous result is that $\lim _{n \rightarrow \infty} y_{n}$ exists and coincides with $\inf \left\{y_{n}: n \in \mathbb{N}\right\}$.

PROOF: Let $x:=\sup _{n \in \mathbb{N}}$. This is an upper bound of the sequence, hence $x_{j} \leqq x$ for all $j \in \mathbb{N}$. Let $\varepsilon>0 . x-\frac{\varepsilon}{2}$ is not an upper bound, hence there exists $N \in \mathbb{N}$ such that $x-\frac{\varepsilon}{2} \leqq x_{N}$. Because $\left(x_{n}\right)_{n}$ is non-decreasing, it follows for all $j \geqq N$ that

$$
x-\varepsilon<x-\frac{\varepsilon}{2} \leqq x_{N} \leqq j \leqq x, \text { hence }\left|x_{j}-x\right|=x-x_{j}<\varepsilon \text { for all } j \geqq N .
$$

This proves that $x=\sup _{n \in \mathbb{N}} x_{n}=\lim _{j \rightarrow \infty} x_{j}$.
@@MFx-Auth


The following proposition appears in ch.8.3 (Limit Inferior and Limit Superior) as cor.8.4 (see p.148). It is made available here because we use it in lecture to give an alternate proof of prop. 16.18 on p.299. There is no reference tag so as not to change the references for the subsequent propositions of this chapter.
Proposition (Sandwiching Theorem for Limits) - No reference tag!
Let $x_{n}, y_{n} \in \mathbb{R}$ be two sequences of real numbers both of which have limits. Assume there is $K \in \mathbb{N}$ such that $x_{n} \leqq y_{n}$ for all $n \geqq K$. Then

$$
\lim _{n \rightarrow \infty} x_{n} \leqq \lim _{n \rightarrow \infty} y_{n} .
$$

PROOF: Let $x_{0}:=\lim _{n \rightarrow \infty} x_{n}$ and $y_{0}:=\lim _{n \rightarrow \infty} y_{n}$. We will show that $\lim _{n \rightarrow \infty}\left(y_{n}-x_{n}\right) \geqq 0$. This suffices to prove the proposition because, according to prop. 8.10.c,

$$
\lim _{n \rightarrow \infty}\left(y_{n}-x_{n}\right)=\lim _{n \rightarrow \infty} y_{n}-\lim _{n \rightarrow \infty} x_{n} .
$$

We abbreviate $z_{n}:=y_{n}-x_{n}$ and $z_{0}:=\lim _{n \rightarrow \infty} z_{n}$. We assume to the contrary that $z_{0}<0$. Let $\varepsilon:=-\frac{1}{2} z_{0}$. Then $\varepsilon>0$. It follows from the definition of limits that there exists $N \in \mathbb{N}$ such that

$$
\begin{align*}
& \left|z_{j}-z_{0}\right|<\varepsilon \text {, i.e., }\left|z_{j}+2 \varepsilon\right|<\varepsilon \\
& \text { i.e., }\left|z_{j}+2 \varepsilon\right|<\varepsilon \text {, i.e., }-\varepsilon<z_{j}+2 \varepsilon<\varepsilon \text {, } \\
& \text { i.e., }-3 \varepsilon<z_{j}<-\varepsilon<0 \text {. } \tag{*}
\end{align*}
$$

for all $j \geqq N$. Let $n_{0}:=\max (K, N)$. We obtain for all $j \geqq n_{0}$ from $(\star)$ that $y_{j}-x_{j}=z_{j}<0$. This contradicts our assumption that $x_{j} \leqq y_{j}$ for all such $j$.

We now briefly address continuity of functions which map real numbers to real numbers. This subject will be addressed in more detail and in more general settings in ch.10.2.1 on p.212. Let $A \subseteq \mathbb{R}$. Informally speaking, a continuous function $f: A \rightarrow \mathbb{R}$ is one whose graph in the xy-plane is a continuous line without any disconnections or gaps. This can be stated in slightly more formal terms by saying that, if the $x$-values are closely together then the $f(x)$-values must be closely together too.

Here is the formal definition.
Definition 8.12 (Continuity in $\mathbb{R}$ ). Let $A \subseteq \mathbb{R}, x_{0} \in A$, and let $f: A \rightarrow \mathbb{R}$ be a real-valued function with domain $A$. We say that $f$ is continuous at $x_{0}{ }^{72}$ and we write

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}} f(x)=f\left(x_{0}\right) \tag{8.45}
\end{equation*}
$$

if the following is true for any sequence $\left(x_{n}\right)$ with values in $A$ :

$$
\begin{equation*}
\text { if } x_{n} \rightarrow x_{0} \text { then } f\left(x_{n}\right) \rightarrow f\left(x_{0}\right) . \tag{8.46}
\end{equation*}
$$

In other words, the following must be true for any sequence $\left(x_{n}\right)$ in $A$ and $x_{0} \in A$ :

$$
\begin{equation*}
\lim _{n \rightarrow \infty} x_{n}=x_{0} \Rightarrow \lim _{n \rightarrow \infty} f\left(x_{n}\right)=f\left(\lim _{n \rightarrow \infty} x_{n}\right)=f\left(x_{0}\right) \tag{8.47}
\end{equation*}
$$

We say that $f$ is continuous if $f$ is continuous at $x_{0}$ for all $x_{0} \in A$.
Remark 8.3. Important points to notice:
a) It is not enough for the above to be true for some sequences that converge to $x_{0}$. Rather, it must be true for all such sequences!
b) We restrict our universe to the domain $A$ of $f: x_{0}$ and the entire sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ must belong to $A$ because there must be function values for $x_{0}$ and all $x_{n}$.

Proposition 8.12 (Rules of arithmetic for continuous real-valued functions with domain in $\mathbb{R}$ ). Let $A \subseteq \mathbb{R}$ and $\alpha \in \mathbb{R}$. Assume that the functions

$$
f(\cdot), g(\cdot), f_{1}(\cdot), f_{2}(\cdot), f_{3}(\cdot), \ldots, f_{n}(\cdot): A \longrightarrow \mathbb{R}
$$

all are continuous at $x_{0} \in A$. Then

[^55]a. Constant functions are continuous everywhere on $A$.
b. The product $f g(\cdot): x \mapsto f(x) g(x)$ is continuous at $x_{0}$. Specifically, $\alpha f(\cdot): x \mapsto \alpha \cdot f(x)$ is continuous at $x_{0}$. Hence, setting $\alpha:=-1,-f(\cdot): x \mapsto-f(x)$ is continuous at $x_{0}$.
c. The sum $f+g(\cdot): x \mapsto f(x)+g(x)$ is continuous at $x_{0}$.
d. Any linear combination ${ }^{73} \sum_{j=0}^{n} a_{j} f_{j}(\cdot): x \mapsto \sum_{j=0}^{n} a_{j} f_{j}(x)$ is continuous in $x_{0}$.

PROOF: This proposition will be generalized in prop. 10.24 on p. 217 and a full proof will be given there. Here we only give a proof for the product $f g$ to demonstrate how knowledge about the convergence of sequences can be employed to prove statements concerning continuity.

Let $x_{n} \in A$ for all $n \in \mathbb{N}$ such that $\lim _{n \rightarrow \infty} x_{n}=x_{0}$. We must show $\lim _{n \rightarrow \infty} f\left(x_{n}\right) g\left(x_{n}\right)=f\left(x_{0}\right) g\left(x_{0}\right)$, and this follows from part d of prop.8.10 (Rules of arithmetic for limits) on p.136.
Theorem 8.1. Let $A \subseteq \mathbb{R}, x_{0} \in A$, and let $f: A \rightarrow \mathbb{R}$ be a real-valued function with domain $A$. Then $f$ is continuous at $x_{0}$ if and only if for any $\varepsilon>0$, no matter how small, there exists $\delta>0$ such that either one of the following equivalent statements is satisfied:

$$
\begin{equation*}
f\left(\left\{x \in A:\left|x-x_{0}\right|<\delta\right\}\right) \subseteq\left\{y \in \mathbb{R}:\left|y-f\left(x_{0}\right)\right|<\varepsilon\right\}, \tag{8.48}
\end{equation*}
$$

$$
\begin{equation*}
\left|x-x_{0}\right|<\delta \Rightarrow\left|f(x)-f\left(x_{0}\right)\right|<\varepsilon \text { for all } x \in A \tag{8.49}
\end{equation*}
$$

PROOF: A generalized version of this theorem will be proved in a later chapter. ${ }^{74}$
Figure 8.1: $\varepsilon-\delta$ continuity: Any $(x, f(x))$ such that $x_{0}-\delta<x<x_{0}+\delta$ must be contained in the shaded rectangle.


### 8.3 Limit Inferior and Limit Superior

Definition 8.13 (Tail sets of a sequence). Let $\left(x_{k}\right)_{k \in \mathbb{N}}$ be a sequence in $\mathbb{R}$. Let

$$
\begin{equation*}
T_{n}:=\left\{x_{j}: j \in \mathbb{N} \text { and } j \geqq n\right\}=\left\{x_{n}, x_{n+1}, x_{n+2}, x_{n+3}, \ldots\right\} \tag{8.50}
\end{equation*}
$$

[^56]be what remains in the sequence after we discard the first $n-1$ elements. We call $\left(T_{n}\right)_{n \in \mathbb{N}}$ the $n$-th tail set of the sequence $\left(x_{k}\right)_{k}$.

Remark 8.4. Some simple properties of tail sets:
a. We deal with sets and not with sequences $T_{n}$ : If, e.g., $x_{n}=(-1)^{n}$ then each $T_{n}=\{-1,1\}$ only contains two items and not infinitely many.
b. The tail set sequence $\left(T_{n}\right)_{n \in \mathbb{N}}$ is "decreasing": If $m<n$ then $T_{m} \supseteq T_{n}$.
c. It follows from (b) and prop. 8.3 on p. 130 and prop.8.11 that

$$
\begin{aligned}
\beta_{n} & :=\sup \left(T_{n}\right) \text { is non-increasing, hence } \lim _{n \rightarrow \infty} \beta_{n}=\inf _{n} \beta_{n} \\
\alpha_{n} & :=\inf \left(T_{n}\right) \text { is non-decreasing, hence } \lim _{n \rightarrow \infty} \alpha_{n}=\sup _{n} \alpha_{n} .
\end{aligned}
$$

These limits can also be expressed as follows.

$$
\begin{align*}
\lim _{n \rightarrow \infty}\left(\sup \left\{x_{j}: j \in \mathbb{N}, j \geqq n\right\}\right) & =\lim _{n \rightarrow \infty}\left(\sup \left(T_{n}\right)\right)=\inf \left(\left\{\sup \left(T_{n}\right): n \in \mathbb{N}\right\}\right)  \tag{8.51}\\
\lim _{n \rightarrow \infty}\left(\inf \left\{x_{j}: j \in \mathbb{N}, j \geqq n\right\}\right) & =\lim _{n \rightarrow \infty}\left(\inf \left(T_{n}\right)\right)=\sup \left(\left\{\inf \left(T_{n}\right): n \in \mathbb{N}\right\}\right) .
\end{align*}
$$

An expression like $\sup \left\{x_{j}: j \in \mathbb{N}, j \geqq n\right\}$ can be written more compactly as $\sup _{j \in \mathbb{N},>n}\left\{x_{j}\right\}$. Moreover, when dealing with sequences $\left(x_{n}\right)$, it is understood in most cases that $n \in \mathbb{N}$ or $n \in \mathbb{Z}_{\geq 0}$ and the last expression simplifies to $\sup _{j \geqq n}\left\{x_{j}\right\}$. This can also be written as $\sup _{j \geqq n}\left(x_{j}\right)$ or $\sup _{j \geqq n} x_{j}$.

In other words, (8.51) becomes

$$
\begin{align*}
& \inf _{n \in \mathbb{N}}\left(\sup _{j \geqq n} x_{j}\right)=\inf \left(\left\{\sup \left(T_{n}\right): n \in \mathbb{N}\right\}\right)=\lim _{n \rightarrow \infty}\left(\sup \left(T_{n}\right)\right)=\lim _{n \rightarrow \infty}\left(\sup _{j \geqq n} x_{j}\right),  \tag{8.52}\\
& \sup _{n \in \mathbb{N}}\left(\inf _{j \geqq n} x_{j}\right)=\sup \left(\left\{\inf \left(T_{n}\right): n \in \mathbb{N}\right\}\right)=\lim _{n \rightarrow \infty}\left(\inf \left(T_{n}\right)\right)=\lim _{n \rightarrow \infty}\left(\inf _{j \geqq n} x_{j}\right) .
\end{align*}
$$

The above leads us to the following definition:
Definition 8.14. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence in $\mathbb{R}$ and let $T_{n}=\left\{x_{j}: j \in \mathbb{N}, j \geqq n\right\}$ be the tail set for $x_{n}$. Assume that $T_{n}$ is bounded above for some $n_{0} \in \mathbb{N}$ (and hence for all $n \geqq n_{0}$ ). We call

$$
\limsup _{n \rightarrow \infty} x_{j}:=\lim _{n \rightarrow \infty}\left(\sup _{j \geqq n} x_{j}\right)=\inf _{n \in \mathbb{N}}\left(\sup _{j \geqq n} x_{j}\right)=\inf _{n \in \mathbb{N}}\left(\sup \left(T_{n}\right)\right)
$$

the lim sup or limit superior of the sequence $\left(x_{n}\right)$.
If, for each $n, T_{n}$ is not bounded above then we say $\limsup _{n \rightarrow \infty} x_{j}=\infty$.
Assume that $T_{n}$ is bounded below for some $n_{0}$ (and hence for all $n \geqq n_{0}$ ). We call

$$
\liminf _{n \rightarrow \infty} x_{j}:=\lim _{n \rightarrow \infty}\left(\inf _{j \geqq n} x_{j}\right)=\sup _{n \in \mathbb{N}}\left(\inf _{j \geqq n} x_{j}\right)=\sup _{n \in \mathbb{N}}\left(\inf \left(T_{n}\right)\right)
$$

the lim inf or limit inferior of the sequence $\left(x_{n}\right)$.
If, for each $n, T_{n}$ is not bounded below then we say $\liminf _{n \rightarrow \infty} x_{j}=-\infty$.


$$
\begin{aligned}
& \text { Given } \varepsilon>0 \\
& \text { There is } N_{1}=N_{1}(\varepsilon) \text { : } \\
& \quad k \geqq N_{1} \Rightarrow\left|x_{k}-\beta\right| \leqq \frac{\varepsilon}{2} \\
& \text { There is } N_{2}=N_{2}(\varepsilon): \\
& \quad k \geqq N_{2} \Rightarrow\left|x_{k}-\alpha\right| \leqq \frac{\varepsilon}{2} \\
& k \geqq N \Rightarrow \\
& \quad \text { both }\left|x_{k}-\beta\right| \leqq \frac{\varepsilon}{2},\left|x_{k}-\alpha\right| \leqq \frac{\varepsilon}{2}
\end{aligned}
$$

Hence only finitely many $x_{k}$ above $\beta+\varepsilon$ or below $\alpha-\varepsilon$;
Hence have subsequences $x_{m_{i}}, x_{n_{j}}$ : $\lim _{i \rightarrow \infty} x_{m_{i}}=\beta, \lim _{j \rightarrow \infty} x_{n_{j}}=\alpha$

Proposition 8.13. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence in $\mathbb{R}$ which is bounded above with tail sets $T_{n}$.
A. Let

$$
\begin{align*}
& \mathscr{U}:=\left\{y \in \mathbb{R}: T_{n} \cap[y, \infty[\neq \emptyset \text { for all } n \in \mathbb{N}\},\right. \\
& \mathscr{U}_{1}:=\left\{y \in \mathbb{R}: \text { for all } n \in \mathbb{N} \text { there exists } k \in \mathbb{Z}_{\geqq 0} \text { such that } x_{n+k} \geqq y\right\}, \\
& \mathscr{U}_{2}:=\left\{y \in \mathbb{R}: \exists \text { subsequence } n_{1}<n_{2}<n_{3}<\cdots \in \mathbb{N} \text { such that } x_{n_{j}} \geqq y \text { for all } j \in \mathbb{N}\right\},  \tag{8.53}\\
& \mathscr{U}_{3}:=\left\{y \in \mathbb{R}: x_{n} \geqq y \text { for infinitely many } n \in \mathbb{N}\right\} .
\end{align*}
$$

Then $\mathscr{U}=\mathscr{U}_{1}=\mathscr{U}_{2}=\mathscr{U}_{3}$.
B. There exists $z=z(\mathscr{U}) \in \mathbb{R}$ such that $\mathscr{U}$ is either an interval $]-\infty, z]$ or an interval $]-\infty, z[$.
C. Let $u:=\sup (\mathscr{U})$. Then $u=z=z(\mathscr{U})$ as defined in part B. Further, $u$ is the only real number such that

C1.
(8.54) $u-\varepsilon \in \mathscr{U}$ and $u+\varepsilon \notin \mathscr{U}$ for all $\varepsilon>0$.

C2. There exists a subsequence $\left(n_{j}\right)_{j \in \mathbb{N}}$ of integers such that $u=\lim _{j \rightarrow \infty} x_{n_{j}}$ and $u$ is the largest real number for which such a subsequence exists.

## PROOF of $A$ :

A.1- $\mathscr{U}=\mathscr{U}_{1}$ : This equality is valid by definition of tailsets of a sequence:

$$
x \in T_{n} \Leftrightarrow x=x_{j} \text { for some } j \geqq n \Leftrightarrow x=x_{n+k} \text { for some } k \in \mathbb{Z}_{\geqq 0}
$$

from which it follows that $x \in T_{n} \cap\left[y, \infty\left[\Leftrightarrow x=x_{n+k} \geqq y\right.\right.$ for some $k \geqq 0$.

## A. $2-\mathscr{U}_{1} \subseteq \mathscr{U}_{2}:$

Let $y \in \mathscr{U}_{1}$ and $n \in \mathbb{N}$. We prove the existence of $\left(n_{j}\right)_{j}$ by induction on $j$.
Base case $j=1$ : As $T_{1} \cap\left[y, \infty\left[\neq \emptyset\right.\right.$ there is some $x \in T_{1}$ such that $y \leqq x<\infty$, i.e., $x \geqq y$. Because $x \in T_{1}=\left\{x_{1}, x_{2}, \ldots\right\}$ we have $x=x_{n_{1}}$ for some integer $n_{1} \geqq 1$; we have proved the existence of $n_{1}$.

Induction assumption: Assume that $n_{1}<n_{2}<\cdots<n_{j_{0}}$ have already been picked.
Induction step: As $y \in \mathscr{U}_{1}$ there is $k \in \mathbb{Z}_{\geqq 0}$ such that $x_{\left(n_{j_{0}}+1\right)+k} \geqq y$. We set $n_{j_{0}+1}:=n_{j_{0}}+1+k$. As this index is strictly larger than $n_{j_{0}}$, the induction step has been proved.
A.3- $\mathscr{U}_{2} \subseteq \mathscr{U}_{3}$ : This is trivial: Let $y \in \mathscr{U}_{2}$. The strictly increasing subsequence $n_{1}<n_{2}<n_{3}<\cdots \in \mathbb{N}$ constitutes the infinite set of indices that is required to grant y membership in $\mathscr{U}_{3}$.
A.4- $\mathscr{U}_{3} \subseteq \mathscr{U}:$ Let $y \in \mathscr{U}_{3}$. Fix some $n \in \mathbb{N}$.

Let $J=J(y) \subseteq \mathbb{N}$ be the infinite set of indices $j$ for which $x_{j} \geqq y$. At most finitely many of those $j$ can be less than that given $n$ and there must be (infinitely many) $j \in J$ such that $j \geqq n$

Pick any one of those, say $j^{\prime}$. Then $x_{j^{\prime}} \in T_{n}$ and $x_{j^{\prime}} \geqq y$. It follows that $y \in \mathscr{U}$
We have shown the following sequence of inclusions:

$$
\mathscr{U}=\mathscr{U}_{1} \subseteq \mathscr{U}_{2} \subseteq \mathscr{U}_{3} \subseteq \mathscr{U}
$$

It follows that all four sets are equal and part $A$ of the proposition has been proved.
PROOF of $\boldsymbol{B}$ : Let $y_{1}, y_{2} \in \mathbb{R}$ such that $y_{1}<y_{2}$ and $y_{2} \in \mathscr{U}$.
It follows from $\left[y_{2}, \infty\left[\subseteq\left[y_{1}, \infty\left[\right.\right.\right.\right.$ and $T_{n} \cap\left[y_{2}, \infty\left[\neq \emptyset\right.\right.$ for all $n \in \mathbb{N}$ that $T_{n} \cap\left[y_{1}, \infty[\neq \emptyset\right.$ for all $n \in \mathbb{N}$, i.e., $y_{1} \in \mathscr{U}$.

We conclude that $\mathscr{U}$ is an interval of the form $]-\infty, z]$ or $]-\infty, z[$ for some $z \in \mathbb{R}$.
PROOF of $C$ : Let $z=z(\mathscr{U})$ as defined in part $B$ and $u:=\sup (\mathscr{U})$.
PROOF of C. 1 - (8.54) part 1, $u-\varepsilon \in \mathscr{U}$ :
As $u-\varepsilon$ is smaller than the least upper bound $u$ of $\mathscr{U}, u-\varepsilon$ is not an upper bound of $\mathscr{U}$. Hence there is $y>u-\varepsilon$ such that $y \in \mathscr{U}$. It follows from part B that $u-\varepsilon \in \mathscr{U} . \checkmark$
PROOF of C. 1 - (8.54) part 2, $u+\varepsilon \notin \mathscr{U}$ :
This is trivial as $u+\varepsilon>u=\sup (\mathscr{U})$ implies that $y \leqq u<u+\varepsilon$ for all $y \in \mathscr{U}$.
But then $y \neq u$ for all $y \in \mathscr{U}$, i.e., $u \notin \mathscr{U}$. This proves $u+\varepsilon \notin \mathscr{U}$.
PROOF of C.2: We construct by induction a sequence $n_{1}<n_{2}<\ldots$ of natural numbers such that

$$
\begin{equation*}
u-1 / j \leq x_{n_{j}} \leqq u+1 / j \tag{8.55}
\end{equation*}
$$

Base case: We have proved as part of C. 1 that $x_{n} \geqq u+1$ for at most finitely many indices $n$. Let $K$ be the largest of those.

As $u-1 \in \mathscr{U}_{3}$, there are infinitely many $n$ such that $x_{n} \geqq u-1$. Infinitely many of those $n$ must exceed $K$. We pick one of them and that will be $n_{1}$. Clearly, $n_{1}$ satisfies (8.55) and this proves the base case.

Induction step: Let us now assume that $n_{1}<n_{2}<\cdots<n_{k}$ satisfying (8.55) have been constructed. $x_{n} \geqq u+1 /(k+1)$ is possible for at most finitely many indices $n$. Let $K$ be the largest of those.

As $u-1 /(k+1) \in \mathscr{U}_{3}$, there are infinitely many $n$ such that $x_{n} \geqq u-1 /(k+1)$. Infinitely many of those $n$ must exceed $\max \left(K, n_{k}\right)$. We pick one of them and that will be $n_{k+1}$. Clearly, $n_{k+1}$ satisfies (8.55) and this finishes the proof by induction.

We now show that $\lim _{j \rightarrow \infty} x_{n_{j}}=u$. Given $\varepsilon>0$ there is $N=N(\varepsilon)$ such that $1 / N<\varepsilon$. It follows from (8.55) that $\left|x_{n_{j}}-u\right| \leqq 1 / j<1 / N<\varepsilon$ for all $j \geqq n$ and this proves that $x_{n_{j}} \rightarrow u$ as $j \rightarrow \infty$.

We will be finished with the proof of $\boldsymbol{C} .2$ if we can show that if $w>u$ then there is no sequence $n_{1}<$ $n_{2}<\ldots$ such that $x_{n_{j}} \rightarrow w$ as $j \rightarrow \infty$.

Let $\varepsilon:=(w-u) / 2$. According to (8.54), $u+\varepsilon \notin \mathscr{U}$. But then, by definition of $\mathscr{U}$, there is $n \in \mathbb{N}$ such that $T_{n} \cap[u+\varepsilon, \infty[=\emptyset$.

But $u+\varepsilon=w-\varepsilon$ and we have $T_{n} \cap\left[w-\varepsilon, \infty\left[=\emptyset\right.\right.$. This implies that $\left|w-x_{j}\right| \geqq \varepsilon$ for all $j \geqq n$ and that rules out the possibility of finding $n_{j}$ such that $\lim _{j \rightarrow \infty} x_{n_{j}}=w$.

Corollary 8.2. As in prop.8.13, let $u:=\sup (\mathscr{U})$. Then $\mathscr{U}=]-\infty, u]$ or $\mathscr{U}=]-\infty, u[$.
Further, $u$ is determined by the following property: For any $\varepsilon>0, x_{n}>u-\varepsilon$ for infinitely many $n$ and $x_{n}>u+\varepsilon$ for at most finitely many $n$.

PROOF: This follows from $\mathscr{U}=\mathscr{U}_{3}$ and parts B and C of prop.8.13.
When we form the sequence $y_{n}=-x_{n}$ then the roles of upper bounds and lower bounds, max and min, inf and sup are reversed. Example: $x$ is an upper bound for $\left\{x_{j}: j \geqq n\right\}$ if and only if $-x$ is a lower bound for $\left\{y_{j}: j \geqq n\right\}$.

The following "dual" version of prop. 8.13 is a direct consequence of the duality of upper/lower bounds, $\min /$ max, inf/sup proposition prop.8.2, p.129.

Proposition 8.14. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence in $\mathbb{R}$ which is bounded below with tail sets $T_{n}$.
A. Let

$$
\begin{align*}
& \left.\left.\mathscr{L}:=\left\{y \in \mathbb{R}: T_{n} \cap\right]-\infty, y\right] \neq \emptyset \text { for all } n \in \mathbb{N}\right\}, \\
& \mathscr{L}_{1}:=\left\{y \in \mathbb{R}: \text { for all } n \in \mathbb{N} \text { there exists } k \in \mathbb{Z} \geqq 0 \text { such that } x_{n+k} \leqq y\right\}, \\
& \mathscr{L}_{2}:=\left\{y \in \mathbb{R}: \exists \text { subsequence } n_{1}<n_{2}<n_{3}<\cdots \in \mathbb{N} \text { such that } x_{n_{j}} \leqq y \text { for all } j \in \mathbb{N}\right\},  \tag{8.56}\\
& \mathscr{L}_{3}:=\left\{y \in \mathbb{R}: x_{n} \leqq y \text { for infinitely many } n \in \mathbb{N}\right\} .
\end{align*}
$$

Then $\mathscr{L}=\mathscr{L}_{1}=\mathscr{L}_{2}=\mathscr{L}_{3}$.
B. There exists $z=z(\mathscr{L}) \in \mathbb{R}$ such that $\mathscr{L}$ is either an interval $[z, \infty[$ or an interval $] z, \infty[$.
C. Let $l:=\inf (\mathscr{L})$. Then $l=z=z(\mathscr{L})$ as defined in part B. Further, $l$ is the only real number such that

C1. (8.57) $l+\varepsilon \in \mathscr{L}$ and $l-\varepsilon \notin \mathscr{L}$
C2. There exists a subsequence $\left(n_{j}\right)_{j \in \mathbb{N}}$ of integers such that $l=\lim _{j \rightarrow \infty} x_{n_{j}}$ and $l$ is the smallest real number for which such a subsequence exists.

PROOF: Let $y_{n}=-x_{n}$ and apply prop.8.13.
Proposition 8.15. Let $\left(x_{n}\right)$ be a bounded sequence of real numbers. As in prop. 8.13 and prop 8.14, let

$$
\begin{align*}
u & =\sup (\mathscr{U})=\sup \left\{y \in \mathbb{R}: T_{n} \cap[y, \infty[\neq \emptyset \text { for all } n \in \mathbb{N}\},\right. \\
l & \left.\left.=\inf (\mathscr{L})=\inf \left\{y \in \mathbb{R}: T_{n} \cap\right]-\infty, y\right] \neq \emptyset \text { for all } n \in \mathbb{N}\right\}, \tag{8.58}
\end{align*}
$$

Then $\quad u=\limsup _{n \rightarrow \infty} x_{j}$ and $l=\liminf _{n \rightarrow \infty} x_{j}$.
Proof that $u=\limsup _{n \rightarrow \infty} x_{j}$ : Let

$$
\begin{equation*}
\beta_{n}:=\sup _{j \geqq n} x_{j}, \quad \beta:=\inf _{n} \beta_{n}=\limsup _{n \rightarrow \infty} x_{n} . \tag{8.59}
\end{equation*}
$$

We will prove that $\beta$ has the properties listed in prop.8.13.C that uniquely characterize $u$ : For any $\varepsilon>0$, we have

$$
\beta-\varepsilon \in \mathscr{U} \quad \text { and } \beta+\varepsilon \notin \mathscr{U}
$$

Another way of saying this is that

$$
\begin{equation*}
b \in \mathscr{U} \text { for } b<\beta \text { and } a \notin \mathscr{U} \text { for } a>\beta \text {. } \tag{8.60}
\end{equation*}
$$

We now prove the latter characterization.
Let $a \in \mathbb{R}, a>\beta=\inf \left\{\beta_{n}: n \in \mathbb{N}\right\}$. Then $a$ is not a lower bound of the $\beta_{n}: \beta_{n_{0}}<a$ for some $n_{0} \in \mathbb{N}$.
As the $\beta_{n}$ are not increasing in $n$, this implies strict inequality $\beta_{j}<a$ for all $j \geqq n_{0}$. By definition, $\beta_{j}$ is the least upper bound (hence an upper bound) of the tail set $T_{j}$. We conclude that $x_{j}<a$ for all $j \geqq n_{0}$.

From that we conclude that $T_{n} \cap\left[a, \infty\left[=\emptyset\right.\right.$ for all $j \geqq n_{0}$. It follows that $a \notin \mathscr{U}$.
Now let $b \in \mathbb{R}, b<\beta=$ g.l. $b\left\{\beta_{n}: n \in \mathbb{N}\right\}$. As $\beta \leqq \beta_{n}$ we obtain $b<\beta_{n}$ for all $n$.
In other words, $b<\sup \left(T_{n}\right)$ for all $n$ : It is possible to pick some $x_{k} \in T_{n}$ such that $b<x_{k}$.
But then $T_{n} \cap[b, \infty[\neq \emptyset$ for all $n$ and we conclude that $b \in \mathscr{U}$.
We put everything together and see that $\beta$ has the properties listed in (8.60). This finishes the proof that $u=\limsup _{n \rightarrow \infty} x_{j}$. The proof that $l=\liminf _{n \rightarrow \infty} x_{j}$ follows again by applying what has already been proved to the sequence ( $-x_{n}$ ).

We have collected everything to prove
Theorem 8.2 (Characterization of limsup and liminf). Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a bounded sequence in $\mathbb{R}$. Then
a1. $\limsup x_{n}$ is the largest of all real numbers $x$ for which $n_{1}<n_{2}<\cdots \in \mathbb{N}$ can be found such $\stackrel{n \rightarrow \infty}{n \rightarrow \infty}=\lim _{j \rightarrow \infty} x_{n_{j}}$.
a2. $\limsup _{n \rightarrow \infty} x_{n}$ is the only real number $u$ such that, for all $\varepsilon>0$, the following is true: $\stackrel{n \rightarrow \infty}{x_{n}}>u+\varepsilon$ for at most finitely many $n$ and $x_{n}>u-\varepsilon$ for infinitely many $n$.
b1. $\liminf _{n \rightarrow \infty} x_{n}$ is the smallest of all real numbers $x$ for which $n_{1}<n_{2}<\cdots \in \mathbb{N}$ can be found such that $x=\lim _{j \rightarrow \infty} x_{n_{j}}$.
b2. $\liminf _{n \rightarrow \infty} x_{n}$ is the only real number $l$ such that, for all $\varepsilon>0$, the following is true:
$x_{n}<l-\varepsilon$ for at most finitely many $n$ and $x_{n}<l+\varepsilon$ for infinitely many $n$.
PROOF: We know from prop.8.15 on p. 144 that $\limsup _{n \rightarrow \infty} x_{n}$ is the unique number $u$ described in part $C$ of prop.8.13, p.142:
$u-\varepsilon \in \mathscr{U}$ and $u+\varepsilon \notin \mathscr{U}$ for all $\varepsilon>0$
and $u$ is the largest real number for which there exists a subsequence $\left(n_{j}\right)_{j \in \mathbb{N}}$ of integers such that $u=\lim _{j \rightarrow \infty} x_{n_{j}}$.
$u-\varepsilon \in \mathscr{U}=\mathscr{U}_{3}$ (see part A of prop.8.15) means that there are infinitely many $n$ such that $x_{n} \geqq u-\varepsilon$ and $u+\varepsilon \notin \mathscr{U}=\mathscr{U}_{3}$ means that there are at most finitely many $n$ such that $x_{n} \geqq u+\varepsilon$. This proves a1 and $a 2$.

We also know from prop.8.15 that $\liminf _{n \rightarrow \infty} x_{n}$ is the unique number $l$ described in part $C$ of prop.8.14, p.144: $l+\varepsilon \in \mathscr{L}$ and $l-\varepsilon \notin \mathscr{L}$ for all $\varepsilon>0$ and $l$ is the smallest real number for which there exists a subsequence $\left(n_{j}\right)_{j \in \mathbb{N}}$ of integers such that $u=\lim _{j \rightarrow \infty} x_{n_{j}}$.
$l+\varepsilon \in \mathscr{L}=\mathscr{L}_{3}$ (see part $\boldsymbol{A}$ of prop.8.15) means that there are infinitely many $n$ such that $x_{n} \leqq l+\varepsilon$ and $l-\varepsilon \notin \mathscr{L}=\mathscr{L}_{3}$ means that there are at most finitely many $n$ such that $x_{n} \leqq l-\varepsilon$. This proves $\mathbf{b 1}$ and b2.

PROOF of thm.8.2 without the use of prop.8.15, prop.8.13 and the dual propositions for the liminf.
Step 1: Let $\varepsilon>0$. It follows from $\beta_{n}=\sup \left(T_{n}\right)=\sup \left\{x_{j}: j \geqq n\right\}$ and $\beta_{n} \searrow \beta=\lim \sup _{n} x_{n}$ that $\beta_{n}<\beta+\varepsilon$ for all $n \geqq N$ for a suitable $N=N(\varepsilon) \in \mathbb{N}$. But then $\beta+\varepsilon$ exceeds the upper bound $\beta_{N}$ of $T_{N}$ and follows that all of its elements, i.e., all $x_{n}$ with $n \geqq N$, satisfy $x_{n}<\beta+\varepsilon$. Hence only some or all of the finitely many $x_{1}, x_{2}, \ldots x_{N-1}$ can exceed $\beta+\varepsilon$. It follows that $\beta$ satisfies the first half of a1 of thm.8.2.
Step 2: We create subsequences $\left(x_{n_{j}}\right)_{j}$ and $\left(\beta_{n_{j}}\right)_{j}$ such that

$$
\begin{equation*}
\beta_{n_{j}} \geqq x_{n_{j}}>\beta_{n_{j}}-1 / j \tag{8.61}
\end{equation*}
$$

for all $j \in \mathbb{N}$ as follows.
$\beta_{1}=\sup \left(T_{1}\right)$ is the smallest upper bound for $T_{1}$, hence $\beta_{1}-1$ is not an upper bound and we can find some $k \in \mathbb{N}$ such that $\beta_{1} \geqq x_{k}>\beta_{1}-1$. We set $n_{1}:=k$.

Having constructed $n_{1}<n_{2}<\cdots<n_{k}$ such that $\beta_{n_{j}} \geqq x_{n_{j}}>\beta_{n_{j}}-1 / j$ for all $j \leqq k$ we now find $x_{n_{k+1}}$ with an index $n_{k+1}>n_{k}$ as follows.
$\beta_{n_{k}+1}-\frac{1}{k+1}$ is not an upper bound for $T_{n_{k}+1}$, hence there exists some $i \in \mathbb{N}$ such that $x_{n_{k}+i}$ (which belongs to $T_{n_{k}+1}$ ) satisfies

$$
\begin{equation*}
x_{n_{k}+i}>\beta_{n_{k}+1}-\frac{1}{k+1} . \tag{8.62}
\end{equation*}
$$

Let $n_{k+1}:=n_{k}+i$. The sequence $\beta_{n}$ non-increasing (i.e., decreasing) and it follows from $n_{k+1}=$ $n_{k}+i \geqq n_{k}+1$ that $\beta_{n_{k+1}} \leqq \beta_{n_{k}+1}$. But then (8.62) implies that $x_{n_{k+1}}>\beta_{n_{k+1}}-\frac{1}{k+1}$. We note that $x_{n_{k+1}} \leqq \beta_{n_{k+1}}$ because $x_{n_{k+1}} \in T_{n_{k+1}}$ and $\beta_{n_{k+1}}=\sup \left(T_{n_{k+1}}\right)$ is an upper bound for all elements of $T_{n_{k+1}}$. Together with (8.62) we have

$$
\begin{equation*}
\beta_{n_{k+1}} \geqq x_{n_{k+1}}>\beta_{n_{k}+1}-\frac{1}{k+1} \tag{8.63}
\end{equation*}
$$

It follows that $x_{n_{k+1}}$ satisfies (8.61) and the the proof of step 2 is completed.
Step 3: The sequence $x_{n_{j}}$ we constructed in step 2 converges to $\beta=\lim \sup _{n} x_{n}$. This is true because $\lim _{k} \beta_{n_{k}}=\beta, \lim _{k} \beta_{n_{k}}-1 / k=\lim _{k} \beta_{n_{k}}-\lim _{k} 1 / k=\beta-0=\beta$ and $x_{n_{j}}$ is "sandwiched" between two sequences which both converge to the same limit $\beta$.

No subsequence of $\left(x_{n}\right)$ can converge to a number $u$ bigger than $\beta$ : Let $\varepsilon:=1 / 2(u-\beta)$. It follows from step 1 that all but finitely many $x_{j}$ satisfy $x_{j} \leqq \beta+\varepsilon$, hence $x_{j} \leqq u-\varepsilon$.

We conclude that $\left|u-x_{j}\right| \geqq \varepsilon$ for $j \geqq N$ and no subsequence converging to $u$ can be extracted. This proves al of thm.8.2.
Step 4. We still must prove the missing half of thm.8.2.a2: $x_{n}>\beta-\varepsilon$ for infinitely many $n$.
Let $\varepsilon>0$. and let $j \in \mathbb{N}$ be so big that $1 / j<\varepsilon$. Let $x_{n_{j}}$ be the subsequence constructed in step 2 . It follows from (8.61) and $\beta_{n_{j}} \geqq \beta$ and $1 / j<\varepsilon$ that $x_{n_{j}}>\beta-\varepsilon$. This proves the missing half of thm.8.2.a2.

Uniqueness of $\beta$ : Let $v>\beta$ and $\varepsilon:=(v-\beta) / 3$. Because $v-\varepsilon>\beta+\varepsilon$, at most finitely many $x_{n}$ satisfy $x_{n}>v-\varepsilon$. It follows that $v$ does not satisfy part 2 of thm.8.2.a2.

Finally let $v<\beta$. Let $\varepsilon:=(\beta-v) / 3$. Because $\beta-\varepsilon>v+\varepsilon$, infinitely many $x_{n}$ satisfy $x_{n}>v+\varepsilon$. It follows that $v$ does not satisfy part 1 of thm.8.2.a2. We have proved that $\lim \sup _{n} x_{n}$ is uniquely determined by the inequalities of thm.8.2.a2 and we have shown both $\boldsymbol{a} \mathbf{1}$ and $\boldsymbol{a} \mathbf{2}$ of thm.8.2.

Parts $\boldsymbol{b 1}$ and $\boldsymbol{b} \mathbf{2}$ of thm.8.2 follow now easily from applying parts $\boldsymbol{a} \mathbf{1}$ and $\boldsymbol{a} \mathbf{2}$ to the sequence $y_{n}:=-x_{n}$.

Theorem 8.3 (Characterization of limits via limsup and liminf). Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a bounded sequence in $\mathbb{R}$. Then $\left(x_{n}\right)$ converges to a real number if and only if liminf and limsup for that sequence coincide and we have

$$
\begin{equation*}
\lim _{n \rightarrow \infty} x_{n}=\liminf _{n \rightarrow \infty} x_{n}=\limsup _{n \rightarrow \infty} x_{n} . \tag{8.64}
\end{equation*}
$$

PROOF of " $\Rightarrow$ ": Let $L:=\lim _{n \rightarrow \infty} x_{n}$. It follows from prop.8.2 (Characterization of limsup and liminf) parts $\boldsymbol{a 1}$ and $\boldsymbol{b 1}$ on $p .145$ that subsequences can be found which converge to $\liminf _{n \rightarrow \infty} x_{n}$ and to $\limsup _{n \rightarrow \infty} x_{n}$. We also know from prop. 8.20 (Subsequences of convergent sequences) on p. 154 that any convergent subsequence has limit L. This proves $\liminf _{n \rightarrow \infty} x_{n}=\limsup _{n \rightarrow \infty} x_{n}=L$.
PROOF of " $\Leftarrow$ ": ${ }^{75}$ Let $L:=\liminf _{n \rightarrow \infty} x_{n}=\limsup _{n \rightarrow \infty} x_{n}$. Let $\varepsilon>0, J_{1}:=\left\{j \in \mathbb{N}: x_{j}<L-\frac{\varepsilon}{2}\right\}$, $J_{2}:=\left\{j \in \mathbb{N}: x_{j}>L+\frac{\varepsilon}{2}\right\}, J:=J_{1} \cup j_{2}$.

It follows from thm.8.2 (Characterization of limsup and liminf) parts $\boldsymbol{a} \mathbf{2}$ and $\boldsymbol{b} \mathbf{2}$ that $J_{1}, J_{2}$, and hence $J$ only contain at most finitely many indices.

If $J=\emptyset$ let $n_{0}:=1$, else let $n_{0}:=\max (J)+1$.
Then $L-\frac{\varepsilon}{2} \leqq x_{j} \leqq L+\frac{\varepsilon}{2}$, hence $\left|x_{j}-L\right|<\varepsilon$ for all $j \geqq n_{0}$. It follows that $L=\lim _{n \rightarrow \infty} x_{n}$.
Proposition 8.16. Let $x_{n}, x_{n}^{\prime} \in \mathbb{R}$ be two sequences of real numbers. Assume there is $K \in \mathbb{N}$ such that $x_{n} \leqq x_{n}^{\prime}$ for all $n \geqq K$. Then

$$
\liminf _{n \rightarrow \infty} x_{n} \leqq \liminf _{n \rightarrow \infty} x_{n}^{\prime} \quad \text { and } \quad \underset{n \rightarrow \infty}{\limsup } x_{n} \leqq \limsup _{n \rightarrow \infty} x_{n}^{\prime} .
$$

PROOF:
We only prove prove the limsup inequality because once we have that, we apply it to the sequences $\left(-x_{n}\right)_{n}$ and $\left(-x_{n}^{\prime}\right)_{n}$ which satisfy $-x_{n}^{\prime} \leqq-x_{n}$ for all $n \geqq K$. We obtain

$$
-\liminf _{n \rightarrow \infty} x_{n}^{\prime}=\limsup _{n \rightarrow \infty}\left(-x_{n}^{\prime}\right) \leqq \limsup _{n \rightarrow \infty}\left(-x_{n}\right)=-\liminf _{n \rightarrow \infty} x_{n},
$$

hence $\liminf _{n \rightarrow \infty} x_{n} \leqq \liminf _{n \rightarrow \infty} x_{n}^{\prime}$. and this proves the liminf inequality of the proposition.
Case 1: Both sequences are bounded.
Let $u:=\lim \sup _{n} x_{n}$ and $u^{\prime}:=\lim \sup _{n} x_{n}^{\prime}$. We assume to the contrary that $u>u^{\prime}$. Then $\varepsilon:=\frac{u-u^{\prime}}{2}>$ 0.

According to cor.8.2 on p. 144 there are infinitely many $x_{n_{1}}, x_{n_{2}}, \ldots$ such that $x_{n_{j}}>u-\varepsilon$. At most finitely of those $n_{j}$ can be less than $K$. We discard those and there still are infinitely many $n_{j} \geqq K$ such that $x_{n_{j}}>u-\varepsilon$.

As $x_{i}^{\prime} \geqq x_{i}$ for all $i \geqq K$, it follows that there are infinitely many $n_{j}$ such that

$$
x_{n_{j}}^{\prime} \geqq x_{n_{j}}>u-\varepsilon=u^{\prime}+\varepsilon .
$$

[^57]We employ cor.8.2 a second time. It also states that there are at most finitely many $x_{n_{j}}^{\prime}$ such that $x_{n_{j}}^{\prime} \geqq u^{\prime}+\varepsilon$. We have reached a contradiction.
Case 2: Not both sequences are bounded above.
If both are bounded below, $\lim _{\inf }^{n} x_{n} \leqq \liminf _{n} x_{n}^{\prime}$ is obtained just as in case 1 , otherwise this is covered in case 3. We now observe what happens to the limits superior.
Case 2a: $x_{n}$ is not bounded above.
Then neither is $x_{n}^{\prime}$, hence all tailsets for both sequences have sup $=\infty$, hence $\limsup _{n} x_{n}=$ $\limsup { }_{n} x_{n}^{\prime}=\infty$.
Case 2b: $x_{n}^{\prime}$ is not bounded above.
Then all tailsets for $x_{n}^{\prime}$ have sup $=\infty$, hence $\lim \sup _{n} x_{n}^{\prime}=\infty$, hence $\lim \sup _{n} x_{n} \leqq \lim \sup _{n} x_{n}^{\prime}$.
Case 3: Not both sequences are bounded below.
If both are bounded above, $\limsup _{n} x_{n} \leqq \lim \sup _{n} x_{n}^{\prime}$ is obtained just as in case 1 , otherwise this is covered in case 2. We now observe what happens to the limits inferior.
Case 2a: $x_{n}^{\prime}$ is not bounded below.
Then neither is $x_{n}$ and we that all tailsets for both sequences have inf $=-\infty$, hence $\liminf _{n} x_{n}=$ $\liminf _{n} x_{n}^{\prime}=\infty$.
Case 2b: $x_{n}$ is not bounded above.
Then all tailsets for $x_{n}$ have inf $=-\infty$, hence $\liminf _{n} x_{n}=-\infty$, hence $\liminf _{n} x_{n} \leqq \liminf _{n} x_{n}^{\prime}$.

Here is the first corollary to prop.8.16.
Corollary 8.3. Let $x_{n}, y_{n} \in \mathbb{R}$ be two sequences of real numbers. Assume there is $K \in \mathbb{N}$ such that $x_{n}=y_{n}$ for all $n \geqq K$. Then

$$
\limsup _{n \rightarrow \infty} x_{n}=\limsup _{n \rightarrow \infty} y_{n} \quad \text { and } \quad \liminf _{n \rightarrow \infty} x_{n}=\liminf _{n \rightarrow \infty} y_{n} .
$$

PROOF: Immediate from prop.8.16.
Here is the second corollary to prop.8.16.
Corollary 8.4. Let $x_{n}, y_{n} \in \mathbb{R}$ be two sequences of real numbers both of which have limits. Assume there is $K \in \mathbb{N}$ such that $x_{n} \leqq y_{n}$ for all $n \geqq K$. Then

$$
\lim _{n \rightarrow \infty} x_{n} \leqq \lim _{n \rightarrow \infty} y_{n} .
$$

PROOF: Immediate from prop.8.16.

### 8.4 Sequences of Sets and Indicator functions and their liminf and limsup

Let $\Omega$ be a nonempty set and let $f_{n}: \Omega \rightarrow \mathbb{R}$ be a sequence of real-valued functions. Let $\omega \in \Omega$. Then $\left(f_{n}(\omega)\right)_{n \in \mathbb{N}}$ is a sequence of real numbers for which we can examine $\lim _{\inf _{n}} f_{n}(\omega)$ and $\lim \sup _{n} f_{n}(\omega)$. We will look at those two expressions as functions of $\omega$.

Example 8.5. The following are examples of sequences of real-valued functions.
a. $\quad f_{n}:[0,1] \rightarrow \mathbb{R} ; x \mapsto x^{n}$ is a sequence of real-valued functions.
b. Let $\Omega:=\left\{(x, y) \in \mathbb{R}^{2}: x^{2}+y^{2} \leqq 1\right\}$. be the unit circle in the Euclidean plane. Then $\varphi_{n}: \Omega \rightarrow \mathbb{R} ; \varphi_{n}(x, y):=\sqrt{x^{2}+y^{2}}$ is a sequence of real-valued functions.
c. Let $f:=\mathbb{R} \rightarrow \mathbb{R}$ be a (fixed, but arbitrary) function which is infinitely often differentiable at all its arguments, i.e., $D_{n} f\left(x_{0}\right):=f^{(n)}\left(x_{0}\right)=\left.\frac{d^{n} f}{d x^{n}}\right|_{x=x_{0}}$ exists for all $x_{0} \in \mathbb{R}$ and all $n \in \mathbb{N}$. Then $h_{n}: \mathbb{R} \rightarrow \mathbb{R} ; x \mapsto D_{n} f(x)$ is a sequence of real-valued functions.

Definition 8.15 (limsup and liminf of a sequence of real valued functions). Let $\Omega$ be a nonempty set and let $f_{n}: \Omega \rightarrow \mathbb{R}$ be a sequence of real-valued functions such that $f_{n}(\omega)$ is bounded for all $\omega \in \Omega$. 76

We define

$$
\begin{align*}
\liminf _{n \rightarrow \infty} f_{n}: \Omega \rightarrow \mathbb{R} & \text { as follows: } \omega \mapsto \liminf _{n \rightarrow \infty} f_{n}(\omega),  \tag{8.66}\\
\limsup _{n \rightarrow \infty} f_{n}: \Omega \rightarrow \mathbb{R} & \text { as follows: } \omega \mapsto \limsup _{n \rightarrow \infty} f_{n}(\omega) . \tag{8.67}
\end{align*}
$$

Remark 8.5. We recall from thm.8.2 (Characterization of limsup and liminf) on p. 145 that

$$
\begin{align*}
\liminf _{n \rightarrow \infty} f_{n}(\omega) & =\inf \left\{\alpha \in \mathbb{R}: \lim _{j \rightarrow \infty} f_{n_{j}}(\omega)=\alpha \text { for some subsequence } n_{1}<n_{2}<\ldots\right\}  \tag{8.68}\\
\limsup _{n \rightarrow \infty} f_{n}(\omega) & =\sup \left\{\beta \in \mathbb{R}: \lim _{j \rightarrow \infty} f_{n_{j}}(\omega)=\beta \text { for some subsequence } n_{1}<n_{2}<\ldots\right\} . \tag{8.69}
\end{align*}
$$

We now characterize $\liminf _{n} f_{n}$ and $\limsup _{n} f_{n}$ for functions $f_{n}$ such that $f_{n}(\omega)$ is either zero or one. We have seen in prop. 6.12 on $p .111$ that any such function is the indicator function $1_{A}$ of the set

$$
A:=\left\{f_{n}=1\right\}=f_{n}^{-1}(\{1\})=\left\{\omega \in \Omega: f_{n}(\omega)=1\right\} \subseteq \Omega .
$$

Proposition 8.17 (liminf and limsup of $\{0,1\}$-functions). Let $\Omega \neq \emptyset$ and $f_{n}: \Omega \rightarrow\{0,1\}$. Let $\omega \in \Omega$. Then either $\liminf _{n} f_{n}(\omega)=1$ or $\liminf _{n} f_{n}(\omega)=0$ and either $\lim \sup _{n} f_{n}(\omega)=1$ or $\limsup { }_{n} f_{n}(\omega)=$ 0 . Further

$$
\begin{align*}
\liminf _{n \rightarrow \infty} f_{n}(\omega) & =1 \Leftrightarrow f_{n}(\omega)=1 \text { except for at most finitely many } n \in \mathbb{N}  \tag{8.70}\\
\limsup _{n \rightarrow \infty} f_{n}(\omega) & =1 \Leftrightarrow f_{n}(\omega)=1 \text { for infinitely many } n \in \mathbb{N} \tag{8.71}
\end{align*}
$$

PROOF: It follows from (8.68), (8.69) and $0 \leqq f_{n}(\omega) \leqq 1$ that $0 \leqq \liminf _{n} f_{n}(\omega) \leqq \lim \sup _{n} f_{n}(\omega) \leqq 1$.

[^58]Definition 8.16 (Extended real valued functions). The set $\overline{\mathbb{R}}:=\mathbb{R} \cup\{\infty\} \cup\{-\infty\}$ is called the extended real numbers line. A mapping $F$ whose codomain is a subset of $\overline{\mathbb{R}}$ is called an extended real valued function.

The above allows to define the functions $\lim _{\inf }^{n} f_{n}$ and $\lim \sup _{n} f_{n}$ even if there are arguments $\omega$ for which $\lim _{\inf _{n}} f_{n}(\omega)$ and/or $\limsup _{n} f_{n}(\omega)$ assumes one of the values $\pm \infty$. There are many issues with functions that allow their arguments to be $\pm \infty$. Example 1: if $F(x)=\infty$ and $F(y)=\infty$, what is $F(x)-F(y)$ ?) Example 2: The following rule is applied to products of extended real numbers:

$$
\begin{equation*}
0 \cdot \pm \infty= \pm \infty \cdot 0=0 \tag{8.65}
\end{equation*}
$$

We conclude from (8.68) that $\liminf _{n} f_{n}(\omega)=0$ if a subsequence $n_{1}<n_{2}<\ldots$ can be found such that $f_{n_{j}}(\omega)=0$ for all $j$ and that $\liminf _{n} f_{n}(\omega)=1$ if no such subsequence exists, i.e., if $f_{n}(\omega)=1$ for all except at most finitely many $n$. This proves not only (both directions(!) of) (8.70) but also that either $\liminf _{n} f_{n}(\omega)=1$ or $\liminf _{n} f_{n}(\omega)=0$

We conclude from (8.69) that $\limsup _{n} f_{n}(\omega)=1$ if a subsequence $n_{1}<n_{2}<\ldots$ can be found such that $f_{n_{j}}(\omega)=1$ for all $j$ and that $\lim \sup _{n} f_{n}(\omega)=0$ if no such subsequence exists, i.e., if $f_{n}(\omega)=0$ for all except at most finitely many $n$. This proves not only (both directions(!) of) (8.71) but also that either $\limsup \sup _{n}(\omega)=1$ or $\limsup \sup _{n} f_{n}(\omega)=0$.

We now look at indicator functions $1_{A_{n}}$ of a sequence of sets $A_{n} \subseteq \Omega$. For such a sequence we define

$$
\begin{equation*}
A_{\star}:=\bigcup_{n \in \mathbb{N}} \bigcap_{j \geqq n} A_{j}, \quad A^{\star}:=\bigcap_{n \in \mathbb{N}} \bigcup_{j \geqq n} A_{j} . \tag{8.72}
\end{equation*}
$$

Proposition 8.18. Let $\omega \in \Omega$. Then

$$
\begin{align*}
& \omega \in A_{\star} \Leftrightarrow \omega \in A_{n} \text { for all except at most finitely many } n \in \mathbb{N} .  \tag{8.73}\\
& \omega \in A^{\star} \Leftrightarrow \omega \in A_{n} \text { for infinitely many } n \in \mathbb{N}, \tag{8.74}
\end{align*}
$$

a. Proof that $\omega \in A_{\star} \Rightarrow \omega \in A_{n}$ for all except at most finitely many $n \in \mathbb{N}$ :

We will prove the contrapositive: Assume that there exists $1 \leqq n_{1}<n_{2}<\ldots$ such that $\omega \notin A_{n_{j}}$ for all $j \in \mathbb{N}$. We must show that $\omega \notin A_{\star}$.

Let $k \in \mathbb{N}$. Then $k \leqq n_{k}$ (think!) and it follows from $\omega \notin A_{n_{k}}$ and $A_{n_{k}} \supseteq \bigcap_{j \geqq n_{k}} A_{j} \supseteq \bigcap_{j \geqq k} A_{j}$ that there is no $k \in \mathbb{N}$ such that $\omega \in \bigcap_{j \geqq k} A_{j}$.

But then $\omega \notin \bigcup_{k} \bigcap_{j \geqq k} A_{j}=A_{\star}$ and we are done with the proof of $\boldsymbol{a}$.
b. Proof that $\omega \in A_{n}$ for all except at most finitely many $n \in \mathbb{N}$ implies $\omega \in A_{\star}$ : By assumption there exists some $N=N(\omega) \in \mathbb{N}$ such that $\omega \in A_{n}$ for all $n \geqq N$.

It follows that $\omega \in \bigcap_{n \geqq N} A_{n} \subseteq \bigcup_{m \in \mathbb{N}} \bigcap_{n \geqq m} A_{n}=A_{\star}$ and $\boldsymbol{b}$ has been proved.
c. Proof that $\omega \in A^{\star} \Rightarrow \omega \in A_{n}$ for infinitely many $n \in \mathbb{N}$ :

Let $\omega \in A^{\star}$. We will recursively construct $1 \leqq n_{1}<n_{2}<\ldots$ such that $\omega \in A_{n_{j}}$ for all $j \in \mathbb{N}$.
We observe that $\omega \in \bigcup_{j \geqq n} A_{j}$ for all $n \in \mathbb{N}$. As $\omega \in \bigcup_{j \geqq 1} A_{j}$ there exists $n_{1} \geqq 1$ such that $\omega \in A_{n_{1}}$ and we have constructed the base case.

Let $k \in \mathbb{N}$. If we already have found $n_{1}<n_{2}<\ldots n_{k}$ such that $\omega \in A_{n_{j}}$ for $1 \leqq j \leqq k$ then we find $n_{k+1}$ as follows: As $\omega \in \bigcup_{j \geqq n_{k}+1} A_{j}$ there exists $n_{k+1} \geqq n_{k}+1$ such that $\omega \in A_{n_{k+1}}$. We have constructed our infinite sequence and this finishes the proof of $c$.
d. Proof that if $\omega \in A_{n}$ for infinitely many $n \in \mathbb{N} \Rightarrow \omega \in A^{\star}$ :

For $n \in \mathbb{N}$ we abbreviate $\Gamma_{n}:=\bigcup_{j \geqq n} A_{j}$.
Let $1 \leqq n_{1}<n_{2}<\ldots$ such that $\omega \in A_{n_{j}}$ for all $j \in \mathbb{N}$. Let $k \in \mathbb{N}$.
Then $n_{k} \geqq k$, hence $\omega \in A_{n_{k}} \in \Gamma_{n_{k}} \subseteq \Gamma_{k}$ for all $k \in \mathbb{N}$, hence $\omega \in \bigcap_{k \in \mathbb{N}} \Gamma_{k}=A^{\star}$. We have proved $d$.

Proposition 8.19 (liminf and limsup of indicator functions).

$$
\begin{equation*}
1_{A_{\star}}=\liminf _{n \rightarrow \infty} 1_{A_{n}} \text { and } 1_{A^{\star}}=\limsup _{n \rightarrow \infty} 1_{A_{n}} \tag{8.75}
\end{equation*}
$$

PROOF: Let $\omega \in \Omega$. Then

$$
\begin{align*}
1_{A_{\star}}(\omega)=1 \Leftrightarrow \omega \in A_{\star} & \Leftrightarrow \omega \in A_{n} \text { for all except at most finitely many } n \in \mathbb{N} \\
& \Leftrightarrow 1_{A_{n}}(\omega)=1 \text { for all except at most finitely many } n \in \mathbb{N}  \tag{8.76}\\
& \Leftrightarrow \liminf _{n} 1_{A_{n}}(\omega)=1
\end{align*}
$$

The second equivalence follows from prop.8.18 and the last equivalence follows from prop.8.17 and this proves the first equation. Similarly we have

$$
\begin{align*}
1_{A^{\star}}(\omega)=1 \Leftrightarrow \omega \in A^{\star} & \Leftrightarrow \omega \in A_{n} \text { for infinitely many } n \in \mathbb{N} \\
& \Leftrightarrow 1_{A_{n}}(\omega)=1 \text { for infinitely many } n \in \mathbb{N}  \tag{8.77}\\
& \Leftrightarrow \limsup _{n} 1_{A_{n}}(\omega)=1
\end{align*}
$$

Again the second equivalence follows from prop.8.18 and the last equivalence follows from prop.8.17.
This last proposition is the reason for the following definition.
Definition 8.17 (limsup and liminf of a sequence of sets). Let $\Omega$ be a nonempty set and let $A_{n} \subseteq$ $\Omega(n \in \mathbb{N})$. We define

$$
\begin{align*}
\liminf _{n \rightarrow \infty} A_{n} & :=\bigcup_{n \in \mathbb{N}} \bigcap_{j \geqq n} A_{j},  \tag{8.78}\\
\limsup _{n \rightarrow \infty} A_{n} & :=\bigcap_{n \in \mathbb{N}} \bigcup_{j \geqq n} A_{j} . \tag{8.79}
\end{align*}
$$

We call $\liminf _{n \rightarrow \infty} A_{n}$ the limit inferior and $\limsup _{n \rightarrow \infty} A_{n}$ the limit superior of the sequence $A_{n}$.
We note that $\liminf _{n \rightarrow \infty} A_{n}=\limsup _{n \rightarrow \infty} A_{n}$ if and only if the functions $\liminf _{n \rightarrow \infty} 1_{A_{n}}$ and $\limsup _{n \rightarrow \infty} 1_{A_{n}}$ coincide (prop. 8.19) which is true if and only if the sequence $1_{A_{n}}(\omega)$ has a limit for all $\omega \in \Omega$ (thm.8.3 on p .147 ). In this case we define

$$
\begin{equation*}
\lim _{n \rightarrow \infty} A_{n}:=\liminf _{n \rightarrow \infty} A_{n}=\limsup _{n \rightarrow \infty} A_{n} \tag{8.80}
\end{equation*}
$$

and we call this set the limit of the sequence $A_{n}$.
Note 8.2 (Notation for limits of monotone sequences of sets). Let $\left(A_{n}\right)$ be a non-decreasing sequence of sets, i.e., $A_{1} \subseteq A_{2} \subseteq \ldots$ and let $A:=\bigcup_{n} A_{n}$. Further let $B_{n}$ be a non-increasing sequence of sets, i.e., $B_{1} \supseteq B_{2} \supseteq \ldots$ and let $B:=\bigcap_{n} B_{n}$. We write suggestively ${ }^{77}$

$$
A_{n} \nearrow A \quad(n \rightarrow \infty), \quad B_{n} \searrow B \quad(n \rightarrow \infty) .
$$

[^59]Example 8.6. Let $A_{n} \subseteq \Omega$.

$$
\begin{align*}
& \text { a. If } A_{n} \nearrow \text { then } \liminf _{n \rightarrow \infty} A_{n}=\underset{n \rightarrow \infty}{\limsup } A_{n}=\bigcup_{n \in \mathbb{N}} A_{n} .  \tag{8.81}\\
& \text { b. If } A_{n} \searrow \text { then } \quad \liminf _{n \rightarrow \infty} A_{n}=\limsup _{n \rightarrow \infty} A_{n}=\bigcap_{n \in \mathbb{N}} A_{n} . \tag{8.82}
\end{align*}
$$

Note 8.3 (Liminf and limsup of number sequences vs their tail sets). Let $x_{n} \in \mathbb{R}$ be a sequence of real numbers. We then can associate with this sequence that of its tail sets $T_{n}:=\left\{x_{j}: j \geqq n\right\}$.

a. Do not confuse $\liminf _{n \rightarrow \infty} x_{n}=\sup _{n}\left(\inf \left(T_{n}\right)\right)$ with $\liminf _{n \rightarrow \infty} T_{n}=\bigcup_{n}\left(\bigcap_{k \geq n} T_{k}\right)$.
b. Do not confuse $\limsup _{n \rightarrow \infty} x_{n}=\inf _{n}\left(\sup \left(T_{n}\right)\right)$ with $\limsup _{n \rightarrow \infty} T_{n}=\bigcap_{n}\left(\bigcup_{k \geq n} T_{k}\right)$.

Those two concepts are very different: $\underset{n}{\lim } \inf x_{n}\left(\lim \sup x_{n}\right)$ is a number: it is the lowest possible (highest possible) limit of a convergent subsequence $\left(x_{n_{j}}\right)_{j \in \mathbb{N}}$. On the other hand we deal with a set(!) $\lim \inf _{n} T_{n}=\lim \sup _{n} T_{n}=\bigcap_{n} T_{n}$. The last equality follows from example 8.6 and the fact the the sequence of tailsets $T_{n}$ is always non-increasing.

### 8.5 Sequences that Enumerate Parts of $\mathbb{Q}$

We informally defined the real numbers in ch. 2.2 (Numbers) on p. 16 as the set of all decimals, i.e., all numbers $x$ which can be written as

$$
\begin{align*}
x & =m+\sum_{j=1}^{\infty} d_{j} 10^{-j} \quad \text { where } d_{j} \text { is a digit, i.e., } d_{j}=0,1,2, \ldots, 9,  \tag{8.83}\\
\text { i.e., } x & =\lim _{k \rightarrow \infty} x_{k} \quad \text { where } x_{k}=m+\sum_{j=1}^{k} d_{j} 10^{-j} . \tag{8.84}
\end{align*}
$$

Each $x_{k}$ is a (finite) sum of fractions, hence $x_{k} \in \mathbb{Q}$.
We proved in cor. 7.8 on $p .127$ that $\mathbb{Q}$ and hence all of its subsets can be sequenced: If $A \subseteq \mathbb{Q}$ there is a sequence $\left(q_{n}\right)_{n}$ of fractions such that $A=\left\{q_{n}: n \subseteq \mathbb{N}\right\}$. We apply this to $A:=\mathbb{Q}$ as follows.

Let $x \in \mathbb{R}$ have the representation (8.84). Then $x_{k} \in \mathbb{Q}$ for each $k \in \mathbb{N}$, hence there is some $n \in \mathbb{N}$ such that $x_{k}=q_{n}$. Of course $n$ depends on $k$, i.e., we have a functional dependency $n=n(k)=n_{k}$. It follows from (8.84) that $q_{n_{k}} \rightarrow x$ as $k \rightarrow \infty$. In other words, we have proved the following

Theorem 8.4 (Universal sequence of rational numbers with convergent subsequences to any real number).

There is a sequence $\left(q_{n}\right)_{n \in \mathbb{N}}$ of fractions which satisfies the following: For any $x \in \mathbb{R}$ there is a sequence $n_{1}, n_{2}, n_{3}, \ldots$, of natural numbers such that $x=\lim _{k \rightarrow \infty} q_{n_{k}}$

## Remark 8.6.

a. The above theorem can be phrased as follows: There is a sequence $\left(q_{n}\right)_{n \in \mathbb{N}}$ of fractions such that for any $x \in \mathbb{R}$ one can find a subsequence $\left(q_{n_{j}}\right)_{j \in \mathbb{N}}$ of $\left(q_{n}\right)_{n}$ which converges to $x$.
b. What is remarkable about thm.8.4: A single sequence $\left(q_{n}\right)_{n}$ is so rich that its ingredients can be used to approximate any item in the uncountable! set $\mathbb{R}$
c. Let $A:=\left\{x \in \mathbb{R}: x^{2} \leqq 2\right\}=[-\sqrt{2}, \sqrt{2}]$ and let $A_{\mathbb{Q}}:=A \cap \mathbb{Q}=\left\{q \in \mathbb{Q}: q^{2} \leqq 2\right\}$. $A$ is of such a shape that for any $x \in A$ the partial sums $x_{k}=m+\sum_{j=1}^{k} d_{j} 10^{-j}$ which converge to $x$ belong to $A_{\mathbb{Q}}$. (Why? Especially, why also for $x= \pm \sqrt{2}$ ?)

### 8.6 Exercises for Ch. 8

Exercise 8.1. Prove example 8.4 part c: Let $z_{n}:=x_{0}$ for some $x_{0} \in \mathbb{R}(n \in \mathbb{N})$. Then $\lim _{n \rightarrow \infty} z_{n}=x_{0}$.
If that is too abstract, try to prove the special case $\mathbf{b}$ first.
Exercise 8.2. Prove the assertions of example 8.6 on p.152.
Exercise 8.3. Prove prop. 8.7 on p. 135 for the case that $\lim _{n \rightarrow \infty} x_{n}$ exists in $\mathbb{R}$ :
Let $\left(x_{n}\right)_{n}$ be a sequence of real numbers such that $\lim _{n \rightarrow \infty} x_{n}$ exists. Let $K \in \mathbb{N}$. For $n \in \mathbb{N}$ let $y_{n}:=x_{n+K}$. Then $\left(y_{n}\right)_{n}$ has the same limit.
Exercise 8.4. Let $A \in \mathbb{R}$ be an interval with endpoints $a<b$, i.e., $A$ is either of $] a, b[] a, b,],[a, b[$, $[a, b]$. Let $a<c<b$ and $A_{1}:=\{x \in A: x \leqq c\}$ and $A_{2}:=\{x \in A: x \geqq c\}$.
Let $f_{1}: A_{1} \rightarrow \mathbb{R}$ and $f_{2}: A_{2} \rightarrow \mathbb{R}$ be two continuous, real-valued functions such that $f_{1}(c)=f_{2}(c)$. Prove that the "spliced" function

$$
f(x):= \begin{cases}f_{1}(x) & \text { for } x \in A_{1}  \tag{8.85}\\ f_{2}(x) & \text { for } x \in A_{1}\end{cases}
$$

is continuous on $A$.
Exercise 8.5. Let $a, b, c, d \in \mathbb{R}$ such that $a<b$ and $c<d$. Let $f:] a, b[\rightarrow] c, d[$ be bijective and strictly monotone, i.e., strictly increasing or decreasing. Prove that both $f$ and $f^{-1}$ are continuous.

Hint: Use thm.8.1 on p.140.
Exercise 8.6. Prove cor.8.3 on p.148:
Let $x_{n}, y_{n} \in \mathbb{R}$ be two sequences of real numbers. Assume there is $K \in \mathbb{N}$ such that $x_{n}=y_{n}$ for all $n \geqq K$. Then

$$
\limsup _{n \rightarrow \infty} x_{n}=\limsup _{n \rightarrow \infty} y_{n} \quad \text { and } \quad \liminf _{n \rightarrow \infty} x_{n}=\liminf _{n \rightarrow \infty} y_{n}
$$

Exercise 8.7. Prove cor.8.4 on p.148:
Let $x_{n}, y_{n} \in \mathbb{R}$ be two sequences of real numbers both of which have limits. Assume there is $K \in \mathbb{N}$ such that $x_{n} \leqq y_{n}$ for all $n \geqq K$. Then

$$
\lim _{n \rightarrow \infty} x_{n} \leqq \lim _{n \rightarrow \infty} y_{n} .
$$

Exercise 8.8. Let $x_{n}:=(-1)^{n}$ for $n \in \mathbb{N}$. Prove that $\liminf _{n} x_{n}=-1$ and $\lim \sup _{n} x_{n}=1$ by working with the tailsets of that sequence. Do not use anything after definition 8.14! Hint: What is $\alpha_{n}$ and $\beta_{n}$ ?

### 8.7 Addenda to Ch. 8 (Real functions)

Proposition 8.20 (Subsequences of convergent sequences). Let $\left(x_{n}\right)_{n}$ be a sequence of real numbers with limit $L:=\lim _{n \rightarrow \infty} x_{n}$. Let $\left(x_{n_{j}}\right)$ be a subsequence. Then $\lim _{j \rightarrow \infty} x_{n_{j}}=L$.

## PROOF:

a. Assume that $\left(x_{n}\right)_{n}$ is convergent, i.e., $L \in \mathbb{R}$. Let $\varepsilon>0$ Because the sequence converges, there exists $N \in \mathbb{N}$ such that $\left|x_{j}-L\right|<\varepsilon$ for all $j \geqq N$. As $n_{j} \geqq j$ for all $j$, we conclude that $n_{j} \geqq N$ whenever $j \geqq N$, hence $\left|x_{n_{j}}-L\right|<\varepsilon$ for all $j \geqq N$. It follows that $\left(x_{n_{j}}\right)$ has limit $L$.
b. Assume that $L=\infty$. Let $K \in \mathbb{R}$. Because $\lim _{j \rightarrow \infty} x_{n_{j}}=\infty$, there exists $N \in \mathbb{N}$ such that $x_{j}>K$ for all $j \geqq N$. As $n_{j} \geqq j$ for all $j$ and hence $n_{j} \geqq N$ whenever $j \geqq N$, we obtain $x_{n_{j}}>K$ for all $j \geqq N$. It follows that $\left(x_{n_{j}}\right)$ has limit $\infty$.
The case $L=-\infty$ is proved similarly.

## 9 Vectors and Vector spaces (Understand this!)

## 9.1 $\mathbb{R}^{n}$ : Euclidean Space

Most if not all of the material of this chapter with the exception of ch.9.2.2 (normed Vector Spaces) on p. 168 is familiar to anyone who took a linear algebra course or, in case of two or three dimensional space, to those who took a course in multivariable calculus.

### 9.1.1 $n$-Dimensional Vectors

The following definition of an n-dimensional vector is a special case of a vector, something which will be defined as an element of an abstract "vector space" . ${ }^{78}$
Definition 9.1 ( $n$-dimensional vectors). Let $n \in \mathbb{N}$. An $n$-dimensional vector is a finite, ordered collection $\vec{v}=\left(x_{1}, x_{2}, x_{3}, \ldots, x_{n}\right)$ of real numbers $x_{1}, x_{2}, x_{3}, \ldots, x_{n}$, i.e., an element of $\mathbb{R}^{n}$. ${ }^{79} n$ is called the dimension of the vector $\vec{v}$.

Here are some examples of vectors:
Example 9.1 (Two-dimensional vectors). The two-dimensional vector $\vec{v}$ with coordinates $x=-1.5$ and $y=\sqrt{2}$ is written $(-1.5, \sqrt{2})$ and we have $(-1.5, \sqrt{2}) \in \mathbb{R}^{2}$. Order matters, so this vector is different from $(\sqrt{2},-1.5) \in \mathbb{R}^{2}$.
Example 9.2 (Three-dimensional vectors). $\overrightarrow{v_{t}}=\left(3-t, 15, \sqrt{5 t^{2}+\frac{22}{7}}\right) \in \mathbb{R}^{3}$ with coordinates $x=$ $3-t, y=15$ and $z=\sqrt{5 t^{3}+\frac{22}{7}}$ is an example of a parametrized vector (parametrized by $t$ ). Each specific value of $t$ defines an element of $\in \mathbb{R}^{3}$, e.g., $\vec{v}_{-2}=\left(5,15, \sqrt{20+\frac{22}{7}}\right)$. note that

$$
F: \mathbb{R} \rightarrow \mathbb{R}^{3} \quad t \mapsto F(t)=\overrightarrow{v_{t}}
$$

defines a function from $\mathbb{R}$ into $\mathbb{R}^{3}$ in the sense of definition (4.6) on p.72. Each argument $s$ has assigned to it one and only one argument $\overrightarrow{v_{s}}=\left(3-s, 15, \sqrt{5 s^{2}+\frac{22}{7}}\right) \in \mathbb{R}^{3}$.

Or, is it rather that we have three functions

$$
\begin{aligned}
& x(\cdot): \mathbb{R} \rightarrow \mathbb{R} \quad t \rightarrow x(t)=3-t, \\
& y(\cdot): \mathbb{R} \rightarrow \mathbb{R} \quad t \rightarrow y(t)=15, \\
& z(\cdot): \mathbb{R} \rightarrow \mathbb{R} \quad t \rightarrow z(t)=\sqrt{5 t^{2}+\frac{22}{7}}, \\
& \text { and } t \rightarrow \overrightarrow{v_{t}}=(x(t), y(t), z(t)) \text { is a vector of three real valued functions } x(\cdot), y(\cdot), z(\cdot) ?
\end{aligned}
$$

Both points of view are correct and it depends on the specific circumstances how you want to interpret $\overrightarrow{v_{t}}$.
Example 9.3 (One-dimensional vectors). A one-dimensional vector has a single coordinate.
For example, $\overrightarrow{w_{1}}=(-3) \in \mathbb{R}^{1}$ with coordinate $x=-3 \in \mathbb{R}$ and $\overrightarrow{w_{2}}=(5.7 a) \in \mathbb{R}^{1}$ with coordinate $x=5.7 a \in \mathbb{R}$ are one-dimensional vectors. $\overrightarrow{w_{2}}$ is not a fixed number but parametrized by $a \in \mathbb{R}$.

Mathematicians do not distinguish between the one-dimensional vector $(x)$ and its coordinate value, the real number $x$. For brevity, they will simply write $\overrightarrow{w_{1}}=-3$ and $\overrightarrow{w_{2}}=5.7 a$.

[^60]Example 9.4 (Vectors as functions). An $n$-dimensional vector $\vec{x}=\left(x_{1}, x_{2}, x_{3}, \cdots, x_{n}\right)$ can be interpreted as a real valued function

$$
\begin{align*}
& f_{\vec{x}}(\cdot):\{1,2,3, \cdots, n\} \longrightarrow \mathbb{R} \quad m \mapsto x_{m}, \text { i.e., } \\
& f_{\vec{x}}(1)=x_{1}, f_{\vec{x}}(2)=x_{2}, \cdots, f_{\vec{x}}(n)=x_{n}, \tag{9.1}
\end{align*}
$$

This can be done in reverse. Any real valued function $f(\cdot):\{1,2,3, \cdots, n\} \longrightarrow \mathbb{R}$ can be associated with the vector $\vec{v}_{f(\cdot)}$ that lists the function values $f(j)$ :

$$
\begin{equation*}
\vec{v}_{f(\cdot)}:=(f(1), f(2), f(3), \cdots, f(n)) \in \mathbb{R}^{n} . \tag{9.2}
\end{equation*}
$$

### 9.1.2 Addition and Scalar Multiplication for $n$-Dimensional Vectors

Definition 9.2 (Addition and scalar multiplication in $\mathbb{R}^{n}$ ). Given are two $n$-dimensional vectors $\vec{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ and $\vec{y}=\left(y_{1}, y_{2}, \ldots, y_{n}\right)$ and a real number $\alpha$.

We define the sum $\vec{x}+\vec{y}$ of $\vec{x}$ and $\vec{y}$ as the vector $\vec{z}$ with the components

$$
\begin{equation*}
z_{1}=x_{1}+y_{1} ; \quad z_{2}=x_{2}+y_{2} ; \quad \ldots ; \quad z_{n}=x_{n}+y_{n} ; \tag{9.3}
\end{equation*}
$$

We define the scalar product $\alpha \vec{x}$ of $\alpha$ and $\vec{x}$ as the vector $\vec{w}$ with the components

$$
\begin{equation*}
w_{1}=\alpha x_{1} ; \quad w_{2}=\alpha x_{2} ; \quad \ldots ; \quad w_{n}=\alpha x_{n} . \tag{9.4}
\end{equation*}
$$

Figure 9.1 below describes vector addition.


Figure 9.1: Adding two vectors.

Adding two vectors $\vec{v}$ and $\vec{w}$ means that you take one of them, say $\vec{v}$, and shift it in parallel (without rotating it in any way or flipping its direction), so that its starting point moves from the origin to the endpoint of the other vector $\vec{w}$. Look at the picture and you see that the vectors $\vec{v}, \vec{w}$ and $\vec{v}$ shifted form three pages of a parallelogram. $\vec{v}+\vec{w}$ is then the diagonal of this parallelogram which starts at the origin and ends at the endpoint of $\vec{v}$ shifted.

### 9.1.3 Length of $n$-Dimensional Vectors and the Euclidean Norm

It is customary to write $\|\vec{v}\|_{2}$ for the length, often also called the Euclidean norm, of the vector $\vec{v}$.
Example 9.5 (Length of one-dimensional vectors). For a vector $\vec{v}=x \in \mathbb{R}$ its length is its absolute value $\|\vec{v}\|_{2}=|x|$. This means that $\|-3.57\|_{2}=|-3.57|=3.57$ and $\|\sqrt{2}\|_{2}=|\sqrt{2}| \approx 1.414$.

Example 9.6 (Length of two-dimensional vectors). We start with an example. Look at $\vec{v}=(4,-3)$. Think of an xy-coordinate system with origin (the spot where $x$-axis and $y$-axis intersect) ( 0,0 ). Then $\vec{v}$ is represented by an arrow which starts at the origin and ends at the point with coordinates $x=4$ and $y=-3$ (see figure 9.2). How long is that arrow?


Figure 9.2: Length of a 2-dimensional vectors.

Think of it as the hypothenuse of a right angle triangle whose two other sides are the horizontal arrow from $(0,0)$ to $(4,0)$ (the vector $\vec{a}=(4,0))$ and the vertical line $\boldsymbol{B}$ between $(4,0)$ and $(4,-3)$. Note that $B$ is not a vector because it does not start at the origin! Obviously (I hope this is obvious) we have $\|\vec{a}\|_{2}=4$ and length-of $(\boldsymbol{B})=3$. Pythagoras tells us that

$$
\|\vec{v}\|_{2}^{2}=\|\vec{a}\|_{2}^{2}+(\text { length-of- } \boldsymbol{B})^{2}
$$

and we obtain for the vector $(4,-3)$ that $\|\vec{v}\|_{2}=\sqrt{16+9}=5$.
The above argument holds for any vector $\vec{v}=(x, y)$ with arbitrary $x, y \in \mathbb{R}$. The horizontal leg on the x -axis is then $\vec{a}=(x, 0)$ with length $|x|=\sqrt{x^{2}}$ and the vertical leg on the y -axis is a line equal in length to $\vec{b}=(0, y)$ the length of which is $|y|=\sqrt{y^{2}}$. The theorem of Pythagoras yields $\|(x, y)\|_{2}^{2}=x^{2}+y^{2}$ which becomes, after taking square roots on both sides,

$$
\begin{equation*}
\|(x, y)\|_{2}=\sqrt{x^{2}+y^{2}} \tag{9.5}
\end{equation*}
$$

Example 9.7 (Length of three-dimensional vectors). This is not so different from the twodimensional case. We build on the previous example. Let $\vec{v}=(4,-3,12)$. Think of an xyzcoordinate system with origin (the spot where $x$-axis, $y$-axis and $z$-axis intersect) $(0,0,0)$. Then $\vec{v}$ is represented by an arrow which starts at the origin and ends at the point with coordinates $x=4$, $y=-3$ and $z=12$. How long is that arrow?

Remember what the standard 3-dimensional coordinate system looks like: The x-axis goes from west to east, the $y$-axis goes from south to north and the $z$-axis goes vertically from down below to the sky. Now drop a vertical line $\boldsymbol{B}$ from the point with coordinates $(4,-3,12)$ to the xy -plane which is "spanned" by the $x$-axis and $y$-axis. This line will intersect the $x y$-plane at the point with coordinates $x=4$ and $y=-3$ (and $z=0$. Why?)

Note that $\boldsymbol{B}$ is not a vector because it does not start at the origin! It should be clear that length-of $(\boldsymbol{B})=|z|=12$.

We connect the origin $(0,0,0)$ with the point $(4,-3,0)$ in the $x y$-plane (the endpoint of $\boldsymbol{B})$.
We can apply what we know 2 -dimensional vectors because this arrow is contained in the xy plane. Matter of fact, we have a genuine two-dimensional vector $\vec{a}=(4,-3)$ because the line starts at the origin. Observe that $\vec{a}$ has the same values 4 and -3 for its $x$ - and $y$-coordinates as the original vector $\vec{v} .{ }^{80}$ We know from the previous example about two-dimensional vectors that

$$
\|\vec{a}\|_{2}^{2}=\|(x, y)\|_{2}^{2}=x^{2}+y^{2}=16+9=25 .
$$

At this point we have constructed a right angle triangle with a) hypothenuse $\vec{v}=(x, y, z)$ where we have $x=4, y=-3$ and $z=12$, b) a vertical leg with length $|z|=12$ and $\mathbf{c}$ ) a horizontal leg with length $\sqrt{x^{2}+y^{2}}=5$. Pythagoras tells us that

$$
\|\vec{v}\|_{2}^{2}=z^{2}+\|(x, y)\|_{2}^{2}=144+25=169, \quad \text { hence }\|\vec{v}\|_{2}=13
$$

None of what we just did depended on the specific values $4,-3$ and 12 . Any vector $(x, y, z) \in \mathbb{R}^{3}$ is the hypothenuse of a right triangle where the square lengths of the legs are $z^{2}$ and $x^{2}+y^{2}$. We conclude that it is true in general that $\|(x, y, z)\|^{2}=x^{2}+y^{2}+z^{2}$, hence

$$
\begin{equation*}
\|(x, y, z)\|=\sqrt{x^{2}+y^{2}+z^{2}} \tag{9.6}
\end{equation*}
$$

The previous examples show how to extend the concept of "length" to vector spaces of any finite dimension:
Definition 9.3 (Euclidean norm). Let $n \in \mathbb{N}$ and $\vec{v}=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ be an n -dimension vector. The Euclidean norm $\|\vec{v}\|_{2}$ of $\vec{v}$ is defined as follows:

$$
\begin{equation*}
\|\vec{v}\|_{2}=\sqrt{x_{1}^{2}+x_{2}^{2}+\ldots+x_{n}^{2}}=\sqrt{\sum_{j=1}^{n} x_{j}^{2}} . \tag{9.7}
\end{equation*}
$$

This definition is important enough to write the special cases for $n=1,2,3$ where $\|\vec{v}\|_{2}$ coincides with the length of $\vec{v}$ :

$$
\begin{array}{ll}
1-\text {-dimensional }: & \|(x)\|_{2}=\sqrt{x^{2}}=|x| \\
2-\text {-dimensional }: & \|(x, y)\|_{2}=\sqrt{x^{2}+y^{2}}  \tag{9.8}\\
3-\text {-dimensional }: & \|(x, y, z)\|_{2}=\sqrt{x^{2}+y^{2}+z^{2}}
\end{array}
$$

[^61]Proposition 9.1 (Properties of the Euclidean norm). Let $n \in \mathbb{N}$. Then the Euclidean norm, viewed as a function

$$
\|\cdot\|_{2}: \mathbb{R}^{n} \rightarrow \mathbb{R} ; \quad \vec{v}=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto\|\vec{v}\|_{2}=\sqrt{\sum_{j=1}^{n} x_{j}^{2}}
$$

has the following three properties:
(9.9a) $\quad\|\vec{v}\|_{2} \geqq 0 \quad \forall \vec{v} \in \mathbb{R}^{n}$ and $\|\vec{v}\|_{2}=0 \quad \Leftrightarrow \vec{v}=0 \quad$ positive definiteness
(9.9b) $\quad\|\alpha \vec{v}\|_{2}=|\alpha| \cdot\|\vec{v}\|_{2} \quad \forall \vec{v} \in \mathbb{R}^{n}, \forall \alpha \in \mathbb{R}$ absolute homogeneity
(9.9c) $\quad\|\vec{v}+\vec{w}\|_{2} \leqq\|\vec{v}\|_{2}+\|\vec{w}\|_{2} \quad \forall \vec{v}, \vec{w} \in \mathbb{R}^{n}$
triangle inequality

## PROOF:

a. It is certainly true that $\|\vec{v}\|_{2} \geqq 0$ for any $n$-dimensional vector $\vec{v}$ because it is defined as $+\sqrt{K}$ where the quantity $K$ is, as a sum of squares, non-negative. If 0 is the zero vector with coordinates $x_{1}=x_{2}=\ldots=x_{n}=0$ then obviously $\|0\|_{2}=\sqrt{0+\ldots+0}=0$. Conversely, let $\vec{v}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ be a vector in $\mathbb{R}^{n}$ such that $\|\vec{v}\|_{2}=0$. This means that $\sqrt{\sum_{j=1}^{n} x_{j}{ }^{2}}=0$ which is only possible if everyone of the non-negative $x_{j}$ is zero. In other words, $\vec{v}$ must be the zero vector 0 .
b. Let $\vec{v}=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ and $\alpha \in \mathbb{R}$. Then

$$
\begin{aligned}
\|\alpha \vec{v}\|_{2} & =\sqrt{\sum_{j=1}^{n}\left(\alpha x_{j}\right)^{2}}=\sqrt{\sum_{j=1}^{n} \alpha^{2} x_{j}^{2}}=\sqrt{\alpha^{2} \sum_{j=1}^{n} x_{j}^{2}}=\sqrt{\alpha^{2}} \sqrt{\sum_{j=1}^{n} x_{j}^{2}} \\
& =\sqrt{\alpha^{2}}\|\vec{v}\|_{2}=|\alpha| \cdot\|\vec{v}\|_{2}
\end{aligned}
$$

because it is true that $\sqrt{\alpha^{2}}=|\alpha|$ for any real number $\alpha$ (see assumption 2.1 on $p .18$ ).
c. The proof will only be given for $n=1,2,3$.
$\boldsymbol{n}=1$ : (9.9.c) simply is the triangle inequality for real numbers (see (2.2) on 19) and we are done.
$\boldsymbol{n}=\mathbf{2 , 3}$ : Look back at the picture about addition of vectors in the plane or in space (see p.156). Remember that for any two vectors $\vec{v}$ and $\vec{w}$ you can always build a triangle whose sides have length $\|\vec{v}\|_{2},\|\vec{w}\|_{2}$ and $\|\vec{v}+\vec{w}\|_{2}$. It is clear that the length of any one side cannot exceed the sum of the lengths of the other two sides, so we get specifically $\|\vec{v}+\vec{w}\|_{2} \leqq\|\vec{v}\|_{2}+\|\vec{w}\|_{2}$ and we are done.

The geometric argument is not exactly an exact proof but I used it nevertheless because it shows the origin of the term "triangle inequality" for property (9.9.c). An exact proof will be given for arbitrary $n \in \mathbb{N}$ as a consequence of the so-called Cauchy-Schwartz inequality (cor.9.1). The inequality itself is stated and proved in prop. 9.9 on $p .169$ in the section which discusses inner products (dot products) on vector spaces.

### 9.2 General Vector Spaces

### 9.2.1 Vector spaces: Definition and Examples

Part of this follows [3] Brin, Matthew and Marchesi, Gerald: Linear Algebra, a text for Math 304, Spring 2016.

Mathematicians are very fond of looking at different objects and figuring out what they have in common. They then create an abstract concept whose items have those properties and examine what they can conclude. For those of you who have had some exposure to object oriented programming: It's like defining a base class, e.g., "mammal", that possesses the core properties of several concrete items such as "horse", "pig", "whale" (sorry - can't require that all mammals have legs). We have looked at the following items that seem to be quite different:

> real numbers
> n-dimensional vectors
> real valued functions

Well, that was disingenuous. We saw that real numbers and one-dimensional vectors are sort of the same (see 9.3 on $p .155$ ). We also saw that n-dimensional vectors can be thought of as real valued functions with domain $X=\{1,2,3, \cdots, n\}$. (see 9.4 on $p .156$ ). Never mind, I'll introduce you now to vector spaces as sets of objects which you can "add" and multiply with real numbers according to rules which are guided by those that apply to addition and multiplication of ordinary numbers.

Here is quick reminder on how we add $n$-dimensional vectors and multiply them with scalars (real numbers) (see (9.1.2) on $p .156$ ). Given are two $n$-dimensional vectors

$$
\vec{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right) \text { and } \vec{y}=\left(y_{1}, y_{2}, \ldots, y_{n}\right) \text { and a real number } \alpha .
$$

Then the sum $\vec{z}=\vec{x}+\vec{y}$ of $\vec{x}$ and $\vec{y}$ is the vector with the components

$$
z_{1}=x_{1}+y_{1} ; \quad z_{2}=x_{2}+y_{2} ; \quad \ldots \quad z_{n}=x_{n}+y_{n}
$$

and the scalar product $\vec{w}=\alpha \vec{x}$ of $\alpha$ and $\vec{x}$ is the vector with the components

$$
w_{1}=\alpha x_{1} ; \quad w_{2}=\alpha x_{2} . \quad \ldots \quad w_{n}=\alpha x_{n}
$$

Example 9.8 (Vector addition and scalar multiplication). We use $n=2$ in this example:
Let $a=(-3,1 / 5), b=(5, \sqrt{2})$ We add those vectors by adding each of the coordinates separately:

$$
a+b=(2,1 / 5+\sqrt{2})
$$

and we multiply $a$ with a scalar $\lambda \in \mathbb{R}$, e.g. $\lambda=100$, by multiplying each coordinate with $\lambda$ :

$$
100 a=100(-3,1 / 5)=(-300,20)
$$

In the last example I have avoided using the notation " $\vec{x}$ " with the cute little arrows on top for vectors. I did that on purpose, because this notation is not all that popular in Math, even for $n$-dimensional vectors, and definitely not for abstract vectors as elements of a vector space. Here now is the definition of a vector space, taken almost word for word from the book "Introductory Real Analysis" (Kolmogorov/Fomin [8]). This definition is quite lengthy because a set needs to satisfy many rules to be a vector space.

Definition 9.4 (Vector spaces (linear spaces)). A nonempty set V of elements $x, y, z, \ldots$ is called a vector space or linear space if it satisfies the following:
A. Any two elements $x, y \in V$ uniquely determine a third element $x+y \in V$, called the sum of $x$ and $y$ with the following properties:

1. $x+y=y+x$ (commutativity);
2. $(x+y)+z=x+(y+z)$ (associativity );
3. There exists an element $0 \in V$, called the zero element, or zero vector, or null vector, with the property that $x+0=x$ for each $x \in V$;
4. For every $x \in V$, there exists an element $-x \in V$, called the negative of $x$, with the property that $x+(-x)=0$ for each $x \in V$. When adding negatives, then there is a convenient short form. We write $x-y$ as an abbreviation for $x+(-y)$;
B. Any real number $\alpha$ and element $x \in V$ together uniquely determine an element $\alpha x \in V$ (sometimes also written $\alpha \cdot x$ ), called the scalar product of $\alpha$ and $x$. It has the following properties:
5. $\alpha(\beta x)=(\alpha \beta) x$;
6. $1 x=x$;
C. The operations of addition and scalar multiplication obey the two distributive laws
7. $(\alpha+\beta) x=\alpha x+\beta x ;$
8. $\alpha(x+y)=\alpha x+\alpha y$;

The elements of a vector space are called vectors.
Definition 9.5 (Subspaces of vector spaces). Let $V$ be a vector space and let $A \subseteq V$ be a nonempty subset of $V$ with the following property: For any $x, y \in A$ and $\alpha \in \mathbb{R}$ the sum $x+y$ and the scalar product $\alpha x$ also belong to $A$. Then $A$ is called a subspace of $V$.

The set $\{0\}$ which only contains the null vector 0 of $V$ is called the nullspace.

## Remark 9.1 (Closure properties).

a. Note that if $\alpha=0$ then $\alpha x=0$. it follows that the null vector belongs to any subspace.
b. We ruled out the case $A=\emptyset$ but did not require that $A$ be a strict subset of $V((2.3)$ on p.13). In other words, the entire vector space $V$ is a subspace of itself.
c. It is trivial to verify that the nullspace $\{0\}$ is a subspace.

Proposition 9.2 (Subspaces are vector spaces). A subspace of a vector space is a vector space, i.e., it satisfies all requirements of definition (9.4).

PROOF: None of the equations that are part of the definition of a vector space magically ceases to be valid just because we look at a subset. The only thing that could go wrong is that some of the expressions might not belong to A anymore. Such can never be the case. Here is the proof for the second distributive law of part $C$.

We must prove that for any $x, y \in A$ and $\lambda \in \mathbb{R}$

$$
\lambda(x+y)=\lambda x+\lambda y .
$$

First, $x+y \in A$ because a subspace contains the sum of any two of its elements. It follows that $\lambda(x+y)$ as product of a real number with an element of $A$ again belongs to $A$ because it is a subspace. Hence the left hand side of the equation belongs to $A$.

Second, both $\lambda x$ and $\lambda y$ belong to $A$ because each is the scalar product of $\lambda$ with an element of $A$ and this set is a subspace. It follows for the same reason that the right hand side of the equation as the sum of two elements of the subspace $A$ belongs to $A$.

Equality of $\lambda(x+y)$ and $\lambda x+\lambda y$ holds because it holds for $x$ and $y$ as elements of $V$.

Remark 9.2 (Closure properties). If a subset $B$ of a larger set $X$ has the property that certain operations on members of $B$ will always yield elements of $B$, then we say that $B$ is closed with respect to those operations.

A subspace is a subset of a vector space which is closed with respect to vector addition and scalar multiplication.

You have already encountered the following examples of vector spaces:
Example 9.9 (Vector space $\mathbb{R}$ ). The real numbers $\mathbb{R}$ are a vector space if you take the ordinary addition of numbers as "+" and the ordinary multiplication of numbers as scalar multiplication.

Example 9.10 (Vector space $\mathbb{R}^{n}$ ). The sets $\mathbb{R}^{n}$ of $n$-dimensional vectors become vector spaces if addition and scalar multiplication are defined as in (9.2) on p.156.

The following example should be thought of as the definition of the very important function spaces $\mathscr{F}(X, \mathbb{R}), \quad \mathscr{B}(X, \mathbb{R}), \quad \mathscr{C}(X, \mathbb{R})$.

Example 9.11 (Vector spaces of real-valued functions). Let $X$ be an arbitrary, nonempty set. Then

$$
\begin{equation*}
\mathscr{F}(X, \mathbb{R}):=\{f(\cdot): f(\cdot) \text { is a real valued function on } X\} \tag{9.10}
\end{equation*}
$$

denotes the set of all real valued functions with domain $X^{81}$ and

$$
\mathscr{B}(X, \mathbb{R}):=\{g(\cdot): g(\cdot) \text { is a bounded real valued function on } X\}
$$

denotes the subset of all bounded real valued functions with domain $X$.
Let $A \subseteq \mathbb{R}$. Then

$$
\mathscr{C}(A, \mathbb{R}):=\{\psi(\cdot): \psi(\cdot) \text { is a continuous real valued function on } A\}
$$

denotes the set of all real-valued continuous functions with domain $A .{ }^{82}$
We list separately the case $X=[a, b]$ where $a, b \in \mathbb{R}$ such that $a<b$. Then

$$
\mathscr{C}([a, b], \mathbb{R}):=\{h(\cdot): h(\cdot) \text { is a continuous real valued function for } a \leqq x \leqq b\}
$$

denotes the set of all contnuous real valued functions with domain $[a, b]$. Note that, for continuous functions, we had to restrict our choice of domain to subsets of real numbers because there is no notion of continuity for functions on abstract domains (and codomains).

If you define addition and scalar multiplication as in (4.13) on p.84, then each of these sets of real-valued functions becomes a vector space for the following reasons:
I: You can verify properties A, B, C of a vector space by looking at the function values for a specific argument $x \in X$ because then you just deal with ordinary real numbers.

[^62]II: The sum of two bounded functions and the product of a bounded function with a scalar is a bounded function. In other words, " + " associates with any two elements $f, g \in \mathscr{B}(X, \mathbb{R})$ a third item $f+g \in \mathscr{B}(X, \mathbb{R})$ and "." associates with any $f \in \mathscr{B}(X, \mathbb{R})$ and $\alpha \in \mathbb{R}$ a third item $\alpha \cdot f \in \mathscr{B}(X, \mathbb{R})$.
III: Likewise, the sum of two continuous functions and the product of a continuous function with a scalar is a continuous function. As for bounded functions, " + " associates with any two elements $f, g \in \mathscr{C}([a, b], \mathbb{R})$ a third item $f+g \in \mathscr{C}([a, b], \mathbb{R})$ and "." associates with any $f \in \mathscr{C}([a, b], \mathbb{R})$ and $\alpha \in \mathbb{R}$ an item $\alpha \cdot f \in \mathscr{C}([a, b], \mathbb{R})$.

It follows from the above that all three function sets are vector spaces and also that 1) $\mathscr{B}(X, \mathbb{R})$ is a subspace of $\mathscr{F}(X, \mathbb{R}), \mathbf{2}) \mathscr{C}(X, \mathbb{R})$ is a subspace of $\mathscr{F}(X, \mathbb{R})$.

We will see in ch. 11 (Compactness) on p .244 that continuous functions defined on a closed interval are bounded. It follows that

$$
\mathscr{C}([a, b], \mathbb{R}) \subseteq \mathscr{B}([a, b], \mathbb{R}) \subseteq \mathscr{F}([a, b], \mathbb{R}) .
$$

We deduce from this that 3) $\mathscr{C}([a, b], \mathbb{R})$ also is a subspace of $\mathscr{B}([a, b], \mathbb{R})$.
It should be noted that, for example, continuous function need not be bounded on open intervals $] a, b\left[\right.$, as the example $f(x)=\frac{1}{x}$ demonstrates for $a=0$ and $b=1$.

Here are some more examples.
Example 9.12 (Subspace $\{(x, y): x=y\}$ ). The set $V:=\{(x, x): x \in \mathbb{R}\}$ of all vectors in the plane with equal $x$ and $y$ coordinates has the following property: For any two vectors $\vec{x}=(a, a)$ and $\vec{y}=(b, b) \in V(a, b \in \mathbb{R})$ and real number $\alpha$ the sum $\vec{x}+\vec{y}=(a+b, a+b)$ and the scalar product $\alpha \vec{x}=(\alpha a, \alpha a)$ have equal $x$-and $y$-coordinates, i.e., they again belong to $V$. It follows that the subset $L$ of $\mathbb{R}^{2}$ is a subspace of $\mathbb{R}^{2}$ (see (9.5) on p.161).

A proof for the following is omitted even though it is not difficult:
Example 9.13 (Subspace $\{(x, y): y=\alpha x\}$ ). Any subset of the form

$$
V_{\alpha}:=\left\{(x, y) \in \mathbb{R}^{2}: y=\alpha x\right\}
$$

is a subspace of $\mathbb{R}^{2}(\alpha \in \mathbb{R})$. Draw a picture: $V_{\alpha}$ is the straight line through the origin in the $x y$-plane with slope $\alpha$.

Example 9.14 (Embedding of linear subspaces). The last example was about the subspace of a bigger space. Now we switch to the opposite concept, the embedding of a smaller space into a bigger space. We can think of the real numbers $\mathbb{R}$ as a part of the $x y$-plane $\mathbb{R}^{2}$ or even 3-dimensional space $\mathbb{R}^{3}$ by identifying a number $a$ with the two-dimensional vector $(a, 0)$ or the three-dimensional vector $(a, 0,0)$. Let $m<n$. It is not a big step from here that the most natural way to uniquely associate an $n$-dimensional vector with an $m$-dimensional vector $\vec{x}:=\left(x_{1}, x_{2}, \ldots, x_{m}\right)$ by adding zero-coordinates to the right:

$$
\vec{x}:=(x_{1}, x_{2}, \ldots, x_{m}, \underbrace{0,0, \ldots, 0}_{n-m \text { times }})
$$

Example 9.15 (All finite-dimensional vectors). Let

$$
\mathfrak{S}:=\bigcup_{n \in \mathbb{N}} \mathbb{R}^{n}=\mathbb{R}^{1} \cup \mathbb{R}^{2} \cup \ldots \cap \mathbb{R}^{n} \cup \ldots
$$

be the set of all vectors of finite (but unspecified) dimension.
We can define addition for any two elements $\vec{x}, \vec{y} \in \mathfrak{S}$ as follows: If $\vec{x}$ and $\vec{y}$ both happen to have the same dimension $n$ then we add them as usual: the sum will be $x_{1}+y_{1}, x_{2}+y_{2}, \ldots, x_{n}+y_{n}$, If not, then one of them, say $\vec{x}$ will have dimension $m$ smaller than the dimension $n$ of $\vec{y}$. We define the sum $\vec{x}+\vec{y}$ as the vector

$$
\vec{z}:=\left(x_{1}+y_{1}, x_{2}+y_{2}, \ldots, x_{m}+y_{m}, y_{m+1}, y_{m+2}, \ldots, y_{n}\right)
$$

Example 9.16 (All sequences of real numbers). Let $\mathbb{R}^{\mathbb{N}}=\prod_{j \in \mathbb{N}} \mathbb{R}$ (see (5.5) on p.101). Is this the same set as $\mathfrak{S}$ from the previous example? The answer is No for the following reason: Each element $x \in \mathfrak{S}$ is of some finite dimension, say $n$, meaning that that it has no more than $n$ coordinates. Each element $y \in \mathbb{R}^{\mathbb{N}}$ is a collection of numbers $y_{1}, y_{2}, \ldots$ none of which need to be zero. In other words, $\mathbb{R}^{\mathbb{N}}$ is the vector space of all sequences of real numbers. Addition is of course done coordinate by coordinate and scalar multiplication with $\alpha \in \mathbb{R}$ is done by multiplying each coordinate with $\alpha$.

There is again a natural way to embed $\mathfrak{S}$ into $\mathbb{R}^{\mathbb{N}}$ as follows: We transform an $n$-dimensional vector $\left(a_{1}, a_{2}, \ldots, a_{n}\right)$ into an element of $\mathbb{R}^{\mathbb{N}}$ (a sequence $\left.\left(a_{j}\right)_{j \in \mathbb{N}}\right)$ by setting $a_{j}=0$ for $j>n$.

Definition 9.6 (linear combinations). Let $V$ be a vector space and let $x_{1}, x_{2}, x_{3}, \ldots, x_{n} \in V$ be a finite number of vectors in $V$. Let $\alpha_{1}, \alpha_{2}, \alpha_{3}, \ldots, \alpha_{n} \in \mathbb{R}$. We call the finite sum

$$
\begin{equation*}
\sum_{j=0}^{n} \alpha_{j} x_{j}=\alpha_{1} x_{1}+\alpha_{2} x_{2}+\alpha_{3} x_{3}+\ldots+\alpha_{n} x_{n} \tag{9.11}
\end{equation*}
$$

a linear combination of the vectors $x_{j}$. The multipliers $\alpha_{1}, \alpha_{2}, \ldots$ are called scalars.
In other words, linear combinations are sums of scalar multiples of vectors. The expression in (9.11) always is an element of $V$, no matter how big $n \in \mathbb{N}$ was chosen:
Proposition 9.3 (Vector spaces are closed w.r.t. linear combinations). Let $V$ be a vector space
and let $x_{1}, x_{2}, x_{3}, \ldots, x_{n} \in V$ be a finite number of vectors in $V$. Let $\alpha_{1}, \alpha_{2}, \alpha_{3}, \ldots, \alpha_{n} \in \mathbb{R}$. Then the linear combination $\sum_{j=0}^{n} \alpha_{j} x_{j}$ also belongs to $V$. Note that this is also true for subspaces because those are vector spaces, too.

## PROOF: Trivial.

Proposition 9.4. Let $V$ be a vector space and let $\left(W_{i}\right)_{i \in I}$ be a family of subspaces of $V$. Let $W:=\bigcap\left[W_{i}: i \in I\right]$. Then $W$ is a subspace of $V$.

PROOF: It suffices to show that $W$ is not empty and that any linear combination of items in $W$ belongs to $W$. As $0 \in W_{i}$ for each $i \in I$, it follows that $0 \in W$, hence $W \neq \emptyset$.

Let $x_{1}, x_{2}, \ldots x_{k} \in W$ and $\alpha_{1}, \alpha_{2}, \ldots \alpha_{k} \in \mathbb{R}(k \in \mathbb{N})$. Let $x:=\sum_{j=1}^{k} \alpha_{j} x_{j}$. Then $x \in W_{i}$ for all $i$ because each $W_{i}$ is a vector space, hence $x \in W$.

Definition 9.7 (Linear span). Let $V$ be a vector space and $A \subseteq V$. Then the set

$$
\begin{equation*}
\operatorname{span}(A)=\left\{\sum_{j=1}^{k} \alpha_{j} x_{j}: k \in \mathbb{N}, \alpha_{j} \in \mathbb{R}, x_{j} \in A(1 \leqq j \leqq k)\right\} \tag{9.12}
\end{equation*}
$$

of all linear combinations of vectors in $A$ is called the span or linear span of $A$.

Proposition 9.5. Let $V$ be a vector space and $A \subseteq V$. Then $\operatorname{span}(A)$ is a subspace of $V$.
PROOF: Let $y_{j} \in \operatorname{span}(A)$ for $j=1,2, \ldots k$, i.e. $y_{j}$ is a linear combination of vectors $x_{j, 1}, x_{j, 2}, \ldots x_{j, n_{j}} \in A$. But then any linear combination of $y_{1}, y_{2}, \ldots y_{k}$ is a linear combination of the vectors

$$
x_{1,1}, x_{1,2}, \ldots, x_{1, n_{1}}, x_{2,1}, x_{2,2}, \ldots, x_{2, n_{2}}, \ldots, x_{k, 1}, x_{k, 2}, \ldots, x_{k, n_{k}}
$$

Theorem 9.1. Let $V$ be a vector space and $A \subseteq V$.
Let $\mathfrak{V}:=\{W \subseteq V: W \supseteq A$ and $W$ is a subspace of $V\}$. Then $\operatorname{span}(A)=\bigcap[W: W \in \mathfrak{V}]$.
PROOF: Clearly, span $(A) \supseteq A$ It follows from prop.9.5 that $\operatorname{span}(A) \in \mathfrak{V}$, hence $\operatorname{span}(A) \supseteq \bigcap[W: W \in$ $\mathfrak{V}]$.

On the other hand, Any subspace $W$ of $V$ that contains $A$ also contains all its linear combinations, hence $\operatorname{span}(A) \subseteq W$ for all $W \in \mathfrak{V}$. But then $\operatorname{span}(A) \subseteq \bigcap[W: W \in \mathfrak{V}]$.

Remark 9.3 (Linear span $(\mathrm{A})=$ subspace generated by A). Let $V$ be a vector space and $A \subseteq V$. Theorem 9.1 justifies to call $\operatorname{span}(A):=$ the subspace generated by $A$.
Definition 9.8 (linear mappings). Let $V_{1}, V_{2}$ be two vector spaces.
Let $f(\cdot): V_{1} \rightarrow V_{2}$ be a function with the following properties:

$$
\begin{align*}
& f(x+y)=f(x)+f(y) \quad \forall x, y \in V_{1}  \tag{9.13a}\\
& f(\alpha x)=\alpha f(x) \quad \forall x \in V_{1}, \forall \alpha \in \mathbb{R} \tag{9.13b}
\end{align*}
$$

## additivity <br> homogeneity

Then we call $f(\cdot)$ a linear function or linear mapping.
Note 9.1 (Note on homogeneity). We encountered "absolute homogeneity" when examining the properties of the Euclidean norm ((9.9) on p.159). That is not the same concept as homogeneity for linear functions because you had to take the absolute value $|\alpha|$ instead of $\alpha$.

Remark 9.4 (Linear mappings are compatible with linear combinations). We saw in the last proposition that vector spaces are closed with respect to linear combinations. Linear functions and linear combinations work harmoniously in the following sense:

A: The image of the sum is the sum of the images,
B: The image of the scalar multiple is the scalar multiple of the image,
C: The image of the linear combination is the linear combination of the images.
Mathematicians express this by saying that linear mappings preserve or are compatible with linear combinations.

The proof of item C in the previous remark is given in the next proposition.
Proposition 9.6 (Linear mappings preserve linear combinations). Let $V_{1}, V_{2}$ be two vector spaces.
Let $f(\cdot): V_{1} \rightarrow V_{2}$ be a linear map and let $x_{1}, x_{2}, x_{3}, \ldots, x_{n} \in V_{1}$ be a finite number of vectors in the domain $V_{1}$ of $f(\cdot)$. Let $\lambda_{1}, \lambda_{2}, \lambda_{3}, \ldots, \lambda_{n} \in \mathbb{R}$. Then $f(\cdot)$ preserves any such linear combination, i.e.,

$$
\begin{equation*}
f\left(\sum_{j=0}^{n} \lambda_{j} x_{j}\right)=\sum_{j=0}^{n} \lambda_{j} f\left(x_{j}\right) . \tag{9.14}
\end{equation*}
$$

PROOF by induction on $n$ : We first note that $f(\lambda x)=\lambda f(x)$ because linear mappings preserve scalar multiples. This proves the base case $n=1$. Because linear mappings also preserve the addition of any two vectors, the proposition holds for $n=2$. Our induction assumption is

$$
f\left(\sum_{j=0}^{k} \lambda_{j} x_{j}\right)=\sum_{j=0}^{k} \lambda_{j} f\left(x_{j}\right) \quad \text { for all } 1 \leqq k<n .
$$

We use it in the second equation $(k=2)$ and the third equation $(k=n-1)$ of the following:

$$
f\left(\sum_{j=0}^{n} \lambda_{j} x_{j}\right)=f\left(\sum_{j=0}^{n-1} \lambda_{j} x_{j}+\lambda_{n} x_{n}\right)=f\left(\sum_{j=0}^{n-1} \lambda_{j} x_{j}\right)+f\left(\lambda_{n} x_{n}\right)=\sum_{j=0}^{n-1} \lambda_{j} f\left(x_{j}\right)+f\left(\lambda_{n} x_{n}\right)=\sum_{j=0}^{n} \lambda_{j} f\left(x_{j}\right)
$$

Here are some examples of linear mappings.
Example 9.17 (Projection on the first coordinate). Let $n \in \mathbb{N}$. The map

$$
\pi_{1}(\cdot): \mathbb{R}^{n} \rightarrow \mathbb{R} \quad\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto x_{1}
$$

is called the projection on the first coordinate or the first coordinate function.
Example 9.18 (Projections on any coordinate). More generally, let $n \in \mathbb{N}$ and $1 \leqq j \leqq n$.

$$
\pi_{j}(\cdot): \mathbb{R}^{n} \rightarrow \mathbb{R} \quad\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto x_{j}
$$

is called the projection on the $j$ th coordinate or the $j$ th coordinate function.
A specific example for $n=2$ : Let $\vec{v}:=(3.5,-2) \in \mathbb{R}^{2}$. Then $\pi_{1}(\vec{v})=3.5$ and $\pi_{2}(\vec{v})=-2$.
Example 9.19 (Projections on any lower dimensional space). In the last two examples we projected $\mathbb{R}^{n}$ onto a one-dimensional space. More generally, we can project $\mathbb{R}^{n}$ onto a vector space $\mathbb{R}^{m}$ of lower dimension $m$ (i.e., we assume $m<n$ ) by keeping $m$ of the coordinates and throwing away the remaining $n-m$ coordinates. Mathematicians express this as follows:

Let $m, n, i_{1}, i_{2}, \ldots, i_{m} \in \mathbb{N}$ such that $m<n$ and $1 \leqq i_{1}<i_{2}<\cdots<i_{m} \leqq n$. The map

$$
\begin{equation*}
\pi_{i_{1}, i_{2}, \ldots, i_{m}}(\cdot): \mathbb{R}^{n} \rightarrow \mathbb{R}^{m} \quad\left(x_{1}, x_{2}, \ldots, x_{n}\right) \mapsto\left(x_{i_{1}}, x_{i_{2}}, \ldots, x_{i_{m}}\right) \tag{9.15}
\end{equation*}
$$

is called the projection on the coordinates $i_{1}, i_{2}, \ldots, i_{m} .{ }^{83}$
Example 9.20. Let $x_{0} \in A$. The mapping

$$
\begin{equation*}
\varepsilon_{x_{0}}: \mathscr{F}(A, \mathbb{R}) \rightarrow \mathbb{R} ; \quad f(\cdot) \mapsto f\left(x_{0}\right) \tag{9.16}
\end{equation*}
$$

which assigns to any real valued function on $A$ its value at the specific point $x_{0}$ is linear because if $h(\cdot)=\sum_{j=0}^{n} a_{j} f_{j}(\cdot)$ then

$$
\varepsilon_{x_{0}}\left(\sum_{j=0}^{n} a_{j} f_{j}\right)=\varepsilon_{x_{0}}(h)=h\left(x_{0}\right)=\sum_{j=0}^{n} a_{j} f_{j}\left(x_{0}\right)=\sum_{j=0}^{n} a_{j} \varepsilon_{x_{0}}\left(f_{j}\right) .
$$

$\varepsilon_{x_{0}}(\cdot)$ is called the abstract integral with respect to point mass at $x_{0}$.
${ }^{83}$ You previously encountered an example where we made use of the projection

$$
\pi_{1,2}(\cdot): \mathbb{R}^{3} \rightarrow \mathbb{R}^{2} \quad(x, y, z) \mapsto(x, y)
$$

This was in the course of computing the length of a 3-dimensional vector (see (9.5) on p.157).

Lemma $9.1(F \circ$ span $=$ span $\circ F)$. [3] Brin/Marchesi Linear Algebra, general lemma 4.1.7: Let $V, W$ be two vector spaces and $F: V \rightarrow W$ a linear mapping from $V$ to $W$. Let $A \subseteq V$. Then

$$
\begin{equation*}
F(\operatorname{span}(A))=\operatorname{span}(F(A)) . \tag{9.17}
\end{equation*}
$$

Proof: See Brin/Marchesi Linear Algebra, general lemma 4.1.7.
Definition 9.9 (Linear dependence and independence). Let $V$ be a vector space and $A \subseteq V$
a. $\quad A$ is called linearly dependent if the following is true: There exist distinct vectors $x_{1}, x_{2}, \ldots x_{k} \in A$ and scalars $\alpha_{1}, \alpha_{2}, \ldots \alpha_{k} \in \mathbb{R}(k \in \mathbb{N})$ such that not all scalars $\alpha_{j}$ are zero $(1 \leqq j \leqq k)$ and $\sum_{j=1}^{k} \alpha_{j} x_{j}=0$.
b. $\quad A$ is called linearly independent if $A$ is not linearly dependent, i.e., if the following is true: Let $x_{1}, x_{2}, \ldots x_{k} \in A$ and $\alpha_{1}, \alpha_{2}, \ldots \alpha_{k} \in \mathbb{R}(k \in \mathbb{N})$.
If $\sum_{j=1}^{k} \alpha_{j} x_{j}=0$ then $\alpha_{j}=0$ for all $1 \leqq j \leqq k$.
Definition 9.10 (Basis of a vector space). Let $V$ be a vector space and $B \subseteq V . B$ is called a basis of $V$ if $\mathbf{a}$. $B$ is linearly independent and $\mathbf{b} . \operatorname{span}(B)=V$.

Lemma 9.2. Let $V$ be a vector space and $A \subseteq V$ linearly independent. If $\operatorname{span}(A) \subsetneq V$ and $y \in \operatorname{span}(A)^{\complement}$ then $A^{\prime}:=A \cup\{y\}$ is linearly independent.

Proof: Let $x_{1}, x_{2}, \ldots x_{k}$ be distinct elements of $A^{\prime}$ and $\alpha_{1}, \alpha_{2}, \ldots \alpha_{k} \in \mathbb{R}(k \in \mathbb{N})$ such that

$$
\begin{equation*}
\sum_{j=1}^{k} \alpha_{j} x_{j}=0 \tag{9.18}
\end{equation*}
$$

We must show that each $\alpha_{j}$ is zero.
Case 1: $y \neq x_{j}$ for all $j$ :
Then $\left\{x_{1}, \ldots, x_{k}\right\} \subseteq A^{\prime} \backslash\{y\}=A$. It follows from the linear independence of $A$ that each $\alpha_{j}$ is zero.
Case 2: $y=x_{j_{0}}$ for some $1 \leq j_{0} \leqq k$ :
We first show that $\alpha_{j_{0}}=0$. This is true because otherwise

$$
\begin{equation*}
x_{j_{0}}=\sum_{j \neq j_{0}} \frac{-\alpha_{j}}{\alpha_{j_{0}}} x_{j} \tag{9.19}
\end{equation*}
$$

would be a linear combination of elements of $A$, contrary to the assumption that $x_{j_{0}}=y \in \operatorname{span}(A)^{\complement}$.
We deduce from (9.18) and $\alpha_{j_{0}}=0$ that

$$
\begin{equation*}
\sum_{j \neq j_{0}} \alpha_{j} x_{j}=0 \tag{9.20}
\end{equation*}
$$

It follows from $\left\{x_{j}: j \neq j_{0}\right\} \subseteq A^{\prime} \backslash\{y\}=A$ and the linear independence of $A$ that $\alpha_{j}=0$ for all $j \neq j_{0}$.

### 9.2.2 Normed Vector Spaces (Study this!)

Definition 9.3 on p. 158 in ch.9.1.3 (Length of $n$-Dimensional Vectors and the Euclidean Norm) gave the definition of the Euclidean norm $\|\vec{x}\|_{2}=\sum_{j=1}^{n} x_{j}^{2}$ in $\mathbb{R}_{n}$. We saw that in dimensions $n=1,2,3$ that $\|\vec{x}\|_{2}$ equals the length of the vector $\vec{x}$ and that prop.9.1 on $p$. 158 "proved" informally for $n=1,2,3$ that $\|\cdot\|_{2}$ satisfies the following three properties:
a. positive definiteness,
b. absolute homogeneity,
c. triangle inequality.

In this chapter we define the norm $\|x\|$ of a vector $x$ in an abstract vector space as a function which satisfies the above three properties, and hence generalizes the concept of the length of a vector in $n$-dimensional space to more general vector spaces. Before we give give this definition, we first introduce the concept of an inner product $x \bullet y$ of two vectors $x$ and $y$. We will see that some of the most important norms, the Euclidean norm among them, can be derived from inner products.

The following definition of inner products and proof of the Cauchy-Schwartz inequality were taken from "Calculus of Vector Functions" (Williamson/Crowell/Trotter [13]).

Definition 9.11 (Inner products). Let $V$ be a vector space with a function

$$
\bullet(\cdot, \cdot): V \times V \rightarrow \mathbb{R} ; \quad(x, y) \mapsto x \bullet y:=\bullet(x, y)
$$

which satisfies the following properties:

| (9.21a) | $x \bullet x \geqq 0 \quad \forall x \in V$ and $\quad x \bullet x=0 \Leftrightarrow x=0$ | positive definiteness |
| :--- | :--- | :--- |
| (9.21b) | $x \bullet y=y \bullet x \quad \forall x, y \in V$ |  |
| (9.21c) | $(x+y) \bullet z=x \bullet z+y \bullet z \quad \forall x, y, z \in V$ | additivity |
| (9.21d) | $(\lambda x) \bullet y=\lambda(x \bullet y) \quad \forall x, y \in V \quad \forall \lambda \in \mathbb{R}$ | homogeneity |

We call such a function an inner product. ${ }^{84}$
Note that additivity and homogeneity of the mapping $x \mapsto x \bullet y$ for a fixed $y \in V$ imply linearity of that mapping and the symmetry property implies that the mapping $y \mapsto x \bullet y$ for a fixed $x \in V$ is linear too. In other words, an inner product is binear in the following sense:

Definition 9.12 (Bilinearity). Let $V$ be a vector space with a function

$$
F(\cdot, \cdot): V \times V \rightarrow \mathbb{R} ; \quad(x, y) \mapsto F(x, y) .
$$

$F(\cdot, \cdot)$ is called bilinear if it is linear in each component, i.e., the mappings

$$
\begin{array}{ll}
F_{1}: V \rightarrow \mathbb{R} ; & x \mapsto F(x, y) \\
F_{2}: V \rightarrow \mathbb{R} ; & y \mapsto F(x, y)
\end{array}
$$

are both linear.

[^63]Proposition 9.7 (Algebraic properties of the inner product). Let $V$ be a vector space with inner product $\bullet(\cdot, \cdot)$. Let $a, b, x, y \in V$. Then

$$
\begin{align*}
& (a+b) \bullet(x+y)=a \bullet x+b \bullet x+a \bullet y+b \bullet y  \tag{9.22a}\\
& (x+y) \bullet(x+y)=x \bullet x+2(x \bullet y)+y \bullet y  \tag{9.22b}\\
& (x-y) \bullet(x-y)=x \bullet x-2(x \bullet y)+y \bullet y \tag{9.22c}
\end{align*}
$$

PROOF of $\boldsymbol{a}$ :

$$
\begin{aligned}
(a+b) \bullet(x+y) & =(a+b) \bullet x+(a+b) \bullet y \\
& =a \bullet x+b \bullet x+a \bullet y+b \bullet y .
\end{aligned}
$$

We used linearity in the second argument for the first equality and linearity in the first argument for the second equality.
PROOF of $\boldsymbol{b}$ and $\boldsymbol{c}$ : Left as an exercise.
The following is the most important example of an inner product.
Proposition 9.8 (Inner product on $\mathbb{R}^{n}$ )). Let $n \in \mathbb{N}$. Then the real-valued function

$$
\begin{equation*}
(\vec{x}, \vec{y}) \mapsto x_{1} y_{1}+x_{2} y_{2}+\ldots+x_{n} y_{n}=\sum_{j=1}^{n} x_{j} y_{j} \quad\left(\vec{x}=\left(x_{1}, \ldots x_{n}\right), \vec{y}=\left(y_{1}, \ldots y_{n}\right)\right) \tag{9.23}
\end{equation*}
$$

is an inner product on $\mathbb{R}^{n} \times \mathbb{R}^{n}$.

## PROOF:

a. For $\vec{x}=\vec{y}$ we obtain $\vec{x} \bullet \vec{x}=\sum_{j=1}^{n} x_{j}^{2}$ and positive definiteness of the inner product follows from

$$
\sum_{j=1}^{n} x_{j}^{2}=0 \Leftrightarrow x_{j}^{2}=0 \forall j \Leftrightarrow x_{j}=0 \forall j .
$$

b. Symmetry is clear because $x_{j} y_{j}=y_{j} x_{j}$.
c. Let $\vec{z}=\left(z_{1}, \ldots z_{n}\right)$. Additivity follows from the fact that $\left(x_{j}+y_{j}\right) z_{j}=x_{j} z_{j}+y_{j} z_{j}$.
d. Homogeneity follows from the fact that $\left(\lambda x_{j}\right) y_{j}=\lambda\left(x_{j} y_{j}\right)$.

Proposition 9.9 (Cauchy-Schwartz inequality for inner products). Let $V$ be a vector space with an inner product

$$
\bullet(\cdot, \cdot): V \times V \rightarrow \mathbb{R} ; \quad(x, y) \mapsto x \bullet y:=\bullet(x, y)
$$

Then

$$
(x \bullet y)^{2} \leqq(x \bullet x)(y \bullet y)
$$

PROOF:
Step 1: We assume first that $x \bullet x=y \bullet y=1$. Then

$$
\begin{aligned}
0 & \leqq(x-y) \bullet(x-y) \\
& =x \bullet x-2 x \bullet y+y \bullet y=2-2(x \bullet y)
\end{aligned}
$$

where the first equality follows from proposition (9.7) on p. 169.
This means $2(x \bullet y) \leqq 2$, i.e., $x \bullet y \leqq 1=(x \bullet x)(y \bullet y)$ where the last equality is true because we had assumed $x \bullet x=y \bullet y=1$. The Cauchy-Schwartz inequality is thus true under that special assumption.

Step 2: General case: We do not assume anymore that $x \bullet x=y \bullet y=1$. If $x$ or $y$ is zero then the Cauchy-Schwartz inequality is trivially true because, say if $x=0$ then the left hand side becomes

$$
(x \bullet y)^{2}=(0 x \bullet y)^{2}=0(x \bullet y)^{2}=0
$$

whereas the right hand side is, as the product of two non-negative numbers $x \bullet x$ and $y \bullet y$, non-negative.
So we can assume that $x$ and $y$ are not zero. On account of the positive definiteness we have $x \bullet x>0$ and $y \bullet y>0$. This allows us to define $u:=x / \sqrt{x \bullet x}$ and $v:=y / \sqrt{y \bullet y}$. But then

$$
\begin{aligned}
& u \bullet u=(x \bullet x) / \sqrt{x \bullet x}^{2}=1 \\
& v \bullet v=(y \bullet y) / \sqrt{y \bullet y}{ }^{2}=1 .
\end{aligned}
$$

We have already seen in step 1 that $u \bullet v \leqq 1$. It follows that

$$
(x \bullet y) /(\sqrt{x \bullet x} \sqrt{y \bullet y})=(x / \sqrt{x \bullet x}) \bullet(y / \sqrt{y \bullet y}) \leqq 1
$$

We multiply both sides with $\sqrt{x \bullet x} \sqrt{y \bullet y}$,

$$
x \bullet y \leqq \sqrt{x \bullet x} \sqrt{y \bullet y} .
$$

We replace $x$ by $-x$ and obtain

$$
-(x \bullet y) \leqq \sqrt{x \bullet x} \sqrt{y \bullet y} .
$$

Because $|x \bullet y|$ is either of $-(x \bullet y)$ or $(x \bullet y)$, it follows from the last two inequalities that

$$
|x \bullet y| \leqq \sqrt{x \bullet x} \sqrt{y \bullet y}
$$

We square this and obtain

$$
(x \bullet y)^{2} \leqq(x \bullet x)(y \bullet y)
$$

and the Cauchy-Schwartz inequality is proved.
Definition 9.13 (sup-norm of bounded real valued functions). Let $X$ be an arbitrary, nonempty set. Let $f: X \rightarrow \mathbb{R}$ be a bounded real valued function on $X$, i.e., there exists a (possibly very large) number $K \geqq 0$ such that $|f(x)| \leqq K$ for all $x \in X .{ }^{85}$ Let

$$
\begin{equation*}
\|f\|_{\infty}:=\sup \{|f(x)|: x \in X\} \tag{9.24}
\end{equation*}
$$

We call $\|f\|_{\infty}$ the supremum norm or sup-norm of the function $f$.
Proposition 9.10 (Properties of the sup norm). Let $X$ be an arbitrary, nonempty set. Let

$$
\mathscr{B}(X, \mathbb{R}):=\{h(\cdot): h(\cdot) \text { is a bounded real valued function on } X\}
$$

(see example 9.11 on $p .162$ ). Then the function

$$
\|\cdot\|_{\infty}: \mathscr{B}(X, \mathbb{R}) \rightarrow \mathbb{R}_{+}, \quad h \mapsto\|h\|_{\infty}=\sup \{|h(x)|: x \in X\}
$$

which assigns to a bounded function on $X$ its sup-norm satisfies the following:
(9.25a) $\|f\|_{\infty} \geqq 0 \forall f \in \mathscr{B}(X, \mathbb{R})$ and $\|f\|_{\infty}=0 \Leftrightarrow f(\cdot)=0$
(9.25b) $\|\alpha f(\cdot)\|_{\infty}=|\alpha| \cdot\|f(\cdot)\|_{\infty} \forall f \in \mathscr{B}(X, \mathbb{R}), \forall \alpha \in \mathbb{R}$
(9.25c) $\|f(\cdot)+g(\cdot)\|_{\infty} \leqq\|f(\cdot)\|_{\infty}+\|g(\cdot)\|_{\infty} \forall f, g \in \mathscr{B}(X, \mathbb{R})$
positive definiteness absolute homogeneity
triangle inequality

[^64]The proof is left as exercise 9.1.
Note 9.2. We previously discussed the Euclidean norm

$$
\begin{equation*}
\|\vec{x}\|_{2}=\sqrt{\sum_{j=1}^{n} x_{j}^{2}} \tag{9.26}
\end{equation*}
$$

for $n$-dimensional vectors $\vec{x}=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$. You saw in (9.1) on p .158 that it satisfies positive definiteness, absolute homogeneity and the triangle inequality, just like the sup-norm. ${ }^{86}$ Those are properties which you associate with the length or size of an object. A very rich mathematical theory can be developed for a generalized definition of length which is based just on those properties.

As mentioned before, mathematicians like to define new objects that are characterized by a certain set of properties. As an example we had the definition of a vector space which encompasses objects as different as finite-dimensional vectors and real valued functions. Accordingly we give a special name to a function defined on a vector space which satisfies positive definiteness, homogeneity and the triangle inequality.
Definition 9.14 (Normed vector spaces). Let $V$ be a vector space. A norm on $V$ is a real valued function

$$
\|\cdot\|: \quad V \rightarrow \mathbb{R} \quad x \mapsto\|x\|
$$

with the following three properties:
(9.27a) $\quad\|x\| \geqq 0 \quad \forall x \in V \quad$ and $\quad\|x\|=0 \Leftrightarrow x=0$

## positive definiteness absolute homogeneity <br> triangle inequality

$$
\begin{equation*}
\|\alpha x\|=|\alpha| \cdot\|x\| \quad \forall x \in V, \forall \alpha \in \mathbb{R} \tag{9.27b}
\end{equation*}
$$

$$
\|x+y\| \leqq\|x\|+\|y\| \quad \forall x, y \in V
$$

We call $V$ a normed vector space and we write $(V,\|\cdot\|)$ instead of $V$ when we wish to emphasize what norm on $V$ we are dealing with.

Proposition 9.11. Let $(V,\|\cdot\|)$ be a normed vector space and let $\gamma>0$.
Let $p: V \rightarrow \mathbb{R}$ be defined as $p(x):=\gamma\|x\|$. Then $p$ also is a norm.
The proof is left as exercise 9.3.
Definition 9.15 ( $p$-norms for $\mathbb{R}^{n}$ ). Let $p \geqq 1$. It will be proved in prop. 9.14 on p .175 that the function

$$
\begin{equation*}
\vec{x} \mapsto\|\vec{x}\|_{p}:=\left(\sum_{j=1}^{n}\left|x_{j}\right|^{p}\right)^{1 / p} \tag{9.28}
\end{equation*}
$$

is a norm on $\mathbb{R}^{n}$ ). This norm is called the $\mathbf{p}$-norm .
The Euclidean norm is a $p$-norm; it is the 2 -norm.
Remark 9.5. We have seen that a vector space can be endowed with more than one norm.
a. We have seen in prop.9.11 that if $x \mapsto\|x\|$ is a norm on a vector space $V$ and $\beta>0$ then $x \mapsto \beta \cdot\|x\|$ also is a norm on $V$.
b. The $p$-norms define a collection of different norms for $\mathbb{R}^{n}$.

[^65]The following theorem shows that an inner product can be associated in a natural fashion with a norm.
Theorem 9.2 (Inner products define norms). Let $V$ be a vector space with an inner product

$$
\bullet(\cdot, \cdot): V \times V \rightarrow \mathbb{R} ; \quad(x, y) \mapsto x \bullet y
$$

Then

$$
\|\cdot\| \bullet: x \mapsto\|x\|=\sqrt{(x \bullet x)}
$$

defines a norm on $V$

## PROOF:

Positive definiteness: follows immediately from that of the inner product.
Absolute homogeneity : Let $x \in V$ and $\lambda \in \mathbb{R}$. Then

$$
\|\lambda x\|_{\bullet}=\sqrt{(\lambda x) \bullet(\lambda x)}=\sqrt{\lambda \lambda(x \bullet x))}=|\lambda| \sqrt{x \bullet x}=|\lambda|\|x\|_{\bullet} .
$$

Triangle inequality : Let $x, y \in V$. Then

$$
\begin{aligned}
\|x+y\|_{\bullet}^{2} & =(x+y) \bullet(x+y) \\
& =x \bullet x+2(x \bullet y)+y \bullet y \\
& \leqq x \bullet x+2|x \bullet y|+y \bullet y \\
& \leqq x \bullet x+2 \sqrt{x \bullet x} \sqrt{y \bullet y}+y \bullet y \\
& =\|x\|_{\bullet}{ }^{2}+2\|x\| \bullet\|y\|_{\bullet}+\|y\|_{\bullet}^{2} \\
& =\left(\|x\|_{\bullet}+\|y\|_{\bullet}\right)^{2} .
\end{aligned}
$$

The second equation uses bilinearity and symmetry of the inner product. The first inequality expresses the simple fact that $\alpha \leqq|\alpha|$ for any number $\alpha$. The second inequality uses Cauchy-Schwartz. The next equality just substitutes the definition $\|x\|_{\bullet}=\sqrt{(x \bullet x)}$ of the norm. The next and last equality is the binomial expansion $(a+b)^{2}=a^{2}+2 a b+b^{2}$ for the ordinary real numbers $a=\|x\|_{\bullet}$ and $b=\|y\|_{\bullet}$.

We take square roots in the above inequality $\|x+y\|_{\bullet}{ }^{2} \leqq\left(\|x\|_{\bullet}+\|y\|_{\bullet}\right)^{2}$ and obtain $\|x+y\|_{\bullet} \leqq$ $\|x\|_{\bullet}+\|y\|_{\bullet}$, the triangle inequality we set out to prove.

Definition 9.16 (Norm for an inner product). Let $V$ be a vector space with an inner product

$$
\bullet(\cdot, \cdot): V \times V \rightarrow \mathbb{R} ; \quad(x, y) \mapsto x \bullet y
$$

Then

$$
\begin{equation*}
\|\cdot\|_{\bullet}: x \mapsto\|x\|_{\bullet}:=\sqrt{(x \bullet x)} \tag{9.29}
\end{equation*}
$$

is called the norm associated with the inner product $\bullet(\cdot, \cdot)$.
It was stated in prop.9.1 on $p .158$ that the Euclidean norm is in fact a norm but only positive definiteness and homogeneity were proved. We now can easily complete the proof.

Corollary 9.1. The Euclidean norm in $\mathbb{R}^{n}$ defined as $\left\|\left(x_{1}, x_{2}, \ldots, x_{n}\right)\right\|_{2}=\sqrt{\sum_{j=1}^{n} x_{j}{ }^{2}}$ (see def. 9.3 on p.158) is a norm.

PROOF: This follows from the fact that

$$
\vec{x} \bullet \vec{y}=\sum_{j=1}^{n} x_{j} y_{j} \quad \text { where } \vec{x}=\left(x_{1}, \ldots, x_{n}\right) \text { and } \vec{y}=\left(y_{1}, \ldots, y_{n}\right) \in \mathbb{R}^{n}
$$

defines an inner product on $\mathbb{R}^{n} \times \mathbb{R}^{n}$ (see prop.9.8) for which $\left\|\left(x_{1}, x_{2}, \ldots, x_{n}\right)\right\|_{2}$ is the associated norm.
We now look at an inner product on the vector space $\mathscr{C}([a, b], \mathbb{R})$ of all continuous real-valued functions on the interval $[a, b]$ which was defined in example 9.11 (Vector spaces of real valued functions) on $p .162$. We use the terminology of [11] Stewart, J: Single Variable Calculus) for the following.

Definition 9.17. Let $a, b \in \mathbb{R}, a<b$ and assume that $f, g:[a, b] \rightarrow \mathbb{R}$ are integrable functions. (See example 4.19 on p.78.)
a. We call the definite integral $\int_{a}^{b} f(x) d x$ the net area between the graph of $f$, the $x$-axis, and the vertical lines through $(a, 0)(y=a)$ and $(b, 0)(y=b)$. The above integral treats areas above the $x$-axis as positive and below the $x$-axis as negative, i.e., the net area is the difference between the areas above the $x$-axis and those below the $x$-axis.
b. We call $\int_{a}^{b}|f(x)| d x$ the area between the graph of $f$, the $x$-axis, and the vertical lines $y=a$ and $y=b$. Note that $f(x)$ has been replaced by its absolute value $|f(x)|$. In contrast to the net area, areas below the $x$-axis are also counted positive.
c. We call $\int_{a}^{b} f(x)-g(x) d x$ the net area between the graphs of $f$ and $g$ and the vertical lines $y=a$ and $y=b$. We call $\int_{a}^{b}|f(x)-g(x)| d x$ the area between the graphs of $f$ and $g$ and the vertical lines $y=a$ and $y=b$.

Example 9.21. Let $f:[-1,1] ; \quad x \mapsto 4 x^{3}$. The antiderivative (see example 4.19 on p .78 ) of $f(\cdot)$ is $x \mapsto x^{4}$ and we compute net area and area as follows:
a. $\quad$ Net area $=\int_{(-1)}^{1} 4 x^{3} d x=\left.x^{4}\right|_{-1} ^{1}=1-1=0$;
b. Area $=\int_{(-1)}^{1} 4\left|x^{3}\right| d x=\int_{(-1)}^{0}\left(-4 x^{3}\right) d x+\int_{0}^{1} 4 x^{3} d x$

$$
=-\left.x^{4}\right|_{-1} ^{0}+-\left.x^{4}\right|_{0} ^{1}=(0-(-1))+(1-0)=2
$$

Let $a, b \in \mathbb{R}$ such that $a<b$. We remember from example 4.19 on $p .78$ ) that continuous functions are integrable. This allows us to compare for $f \in \mathscr{C}([a, b], \mathbb{R})$ the expressions

$$
\begin{equation*}
\|f\|_{\infty}=\sup \{|f(x)|: x \in X\}, \quad \int_{a}^{b}|f(x)| d x, \text { and } \int_{a}^{b}(f(x))^{2} d x . \tag{9.30}
\end{equation*}
$$

All three expressions give in a sense the size of $f$. The sup-norm measures it as the biggest possible displacement from zero, the integral over the absolute value measures the area between the graphs of the functions $x \mapsto f(x)$ and $x \mapsto 0$, and the last expression does the same with the square of $f$. In many respects the use of areas is superior to using the biggest difference to zero.

Squaring $f(\cdot)$ rather than using its absolute value has some mathematical advantages. One of them is that the function $(f, g) \mapsto \int_{a}^{b} f(x) g(x) d x$ defines an inner product on $\mathscr{C}([a, b], \mathbb{R})$ whose associated norm is $f \mapsto \int_{a}^{b}(f(x))^{2} d x$. We will discuss that now. In preparation we prove the following proposition.
Proposition 9.12. Let $a, b \in \mathbb{R}$ such that $a<b$. and let $f:[a, b] \rightarrow[0, \infty[$ be continuous. Then $\int_{a}^{b} f(x) d x=0$ only if $f(x)=0$ for all $\left.x \in\right] a, b[$.

PROOF: Assume that there is $a<x_{0}<b$ such that $f\left(x_{0}\right) \neq 0$, i.e., $f\left(x_{0}\right)>0$. Let $\varepsilon:=\frac{f\left(x_{0}\right)}{2}$. As $f$ is continuous at $x_{0}$ there exists according to thm.8.1 on $p .140$ some $\delta>0$ such that

$$
\begin{equation*}
\left|f\left(x_{0}\right)-f(x)\right|<\varepsilon, \text { hence } f(x)>f\left(x_{0}\right)-\varepsilon=\frac{f\left(x_{0}\right)}{2}=\varepsilon \text { for all } x_{0}-\delta<x<x_{0}+\delta \tag{9.31}
\end{equation*}
$$

Continuity at $x_{0}$ :
If $\left|x-x_{0}\right|<\delta$ then $\left|f\left(x_{0}\right)-f(x)\right|<\varepsilon:$ The graph of $f$ stays within the rectangle with corners $\left(x_{0} \pm \delta, f\left(x_{0}\right) \pm \varepsilon\right)$.


Let $g:[a, b] \rightarrow \mathbb{R}$ be defined as follows.

$$
g(x, y)= \begin{cases}\varepsilon & \text { if } x_{0}-\delta<x<x_{0}+\delta \\ 0 & \text { else } .\end{cases}
$$

It follows from (9.31) that $f \geqq g$, hence $\int_{a}^{b} f(x) d x \geqq \int_{a}^{b} g(x) d x=(2 \delta) \varepsilon>0$.
Proposition 9.13. Let $a, b \in \mathbb{R}$ such that $a<b$. Then the mapping

$$
\begin{equation*}
(f, g) \mapsto f \bullet g:=\int_{a}^{b} f(x) g(x) d x \tag{9.32}
\end{equation*}
$$

defines an inner product on $f \in \mathscr{C}([a, b], \mathbb{R})$.
PROOF: We must prove positive definiteness, symmetry, and linearity in the left argument. In the following let $f, g, h \in \mathscr{C}([a, b], \mathbb{R})$ and $\lambda \in \mathbb{R}$.
a. Positive definiteness: It follows from $f^{2}(x) \geqq 0$ that $f \bullet f=\int_{a}^{b} f^{2}(x) d x \geqq 0$. Clearly, if 0 denotes as usual the zero function $x \mapsto 0$ then $0 \bullet 0=0$. It remains to be shown that if $\int_{a}^{b} f^{2}(x) d x \geqq 0$ then $f=0$. This follows from prop.9.12.
b. Symmetry:

$$
f \bullet g=\int_{a}^{b} f(x) g(x) d x=\int_{a}^{b} g(x) f(x) d x=g \bullet f
$$

c. Additivity and homogeneity: This can be deduced from the well-known formulas

$$
\int_{a}^{b} f(x)+g(x) d x=\int_{a}^{b} f(x) d x+\int_{a}^{b} g(x) d x \text { and } \int_{a}^{b} \lambda g(x) d x=\lambda \int_{a}^{b} g(x) d x .
$$

as follows:

$$
\begin{align*}
& (f+g) \bullet h=\int_{a}^{b}(f(x)+g(x)) h(x) d x=\int_{a}^{b} f(x) h(x) d x+\int_{a}^{b} g(x) h(x) d x=f \bullet h+g \bullet h,  \tag{9.33}\\
& 4) \quad(\lambda f) \bullet g=\int_{a}^{b} \lambda f(x) g(x) d x=\lambda \int_{a}^{b} f(x) g(x) d x=\lambda(f \bullet g) . \tag{9.34}
\end{align*}
$$

According to def. 9.16 (norm for an inner product) and thm. 9.2 (inner products define norms) we now define the norm associated with $f \bullet g=\int_{a}^{b} f(x) g(x) d x$.

Definition 9.18 ( $L_{2}$-Norm for continuous functions). Let $a, b \in \mathbb{R}$ such that $a<b$. Let $f \bullet g$ be the the following inner product on the space $\mathscr{C}([a, b], \mathbb{R})$ of all continuous functions $[a, b] \rightarrow \mathbb{R}$ :

$$
\begin{equation*}
f \bullet g:=\int_{a}^{b} f(x) g(x) d x \tag{9.35}
\end{equation*}
$$

The associated norm

$$
\begin{equation*}
\|\cdot\|_{L^{2}}: f \mapsto\|f\|_{\bullet}=\sqrt{\int_{a}^{b} f^{2}(x) d x} \tag{9.36}
\end{equation*}
$$

is called the $L^{2}-$ norm. of $f$.
We saw in def.9.15 that the Euclidean norm is the p-norm $\|\vec{x}\|_{p}=\left(\sum_{j=1}^{n} x_{j}{ }^{p}\right)^{1 / p}$ for the special case $p=2$. There is an analogue for the $L^{2}$ norm.

Definition 9.19 ( $L^{p}$-norms for $\mathbb{R}^{n}$ ). Let $a, b \in \mathbb{R}$ such that $a<b$ and $p \geqq 1$. It will be shown in prop. 9.15 (The $L^{p}$-norm is a norm) on p. 176 that

$$
\begin{equation*}
f \mapsto\|f\|_{L^{p}}:=\left(\int_{a}^{b}|f(x)|^{p}\right)^{1 / p} \tag{9.37}
\end{equation*}
$$

is a norm on $\mathscr{C}([a, b], \mathbb{R})$. This norm is called the $L^{p}$-norm of $f$. .

### 9.2.3 The Inequalities of Young, Hoelder, and Minkowski ( $\star$ )

Note that this chapter is starred, hence optional.
Proposition 9.14 (The $p$-norm in $\mathbb{R}^{n}$ is a norm). Let $p \in[1, \infty[$.
Then the $p$-norm $\vec{x} \mapsto\|\vec{x}\|_{p}=\left(\sum_{j=1}^{n}\left|x_{j}\right|^{p}\right)^{1 / p}$ is a norm in $\mathbb{R}^{n}$.

## PROOF:

a. Positive definiteness:

Clearly, $\sum_{j=1}^{n}\left|x_{j}\right|^{p} \geqq 0$ because each term $\left|x_{j}\right|^{p}$ is non-negative, hence $\|\vec{x}\|_{p}=\sqrt{\sum_{j=1}^{n}\left|x_{j}\right|^{p}} \geqq 0$. Note that $\|\vec{x}\|_{p}=0$ is only possible if $\left|x_{j}\right|^{p}=0$ for all indices $j$, because, if $x_{j_{0}} \neq 0$ for some $j_{0}$ then $\left|x_{j_{0}}\right|^{p}>0$, hence $\left(\|\vec{x}\|_{p}\right)^{1 / p} \geqq\left|x_{j_{0}}\right|^{p}>0$.
b. Absolute homogeneity:

If $\lambda \in \mathbb{R}$ then

$$
\|(\lambda \vec{x})\|_{p}=\left(\sum_{j=1}^{n}\left(|\lambda|\left|x_{j}\right|\right)^{p}\right)^{1 / p}=\left(|\lambda|^{p} \sum_{j=1}^{n}\left|x_{j}\right|^{p}\right)^{1 / p}=|\lambda|\left(\sum_{j=1}^{n}\left|x_{j}\right|^{p}\right)^{1 / p}=|\lambda|\|\vec{x}\|_{p}
$$

c. Triangle inequality for $p=1$ :

It follows from $\left|x_{j}+y_{j}\right| \leqq\left|x_{j}\right|+\left|y_{j}\right|$ for all $j$ that

$$
\|\vec{x}+\vec{y}\|_{1}=\sum_{j=1}^{n}\left|x_{j}+y_{j}\right| \leqq \sum_{j=1}^{n}\left|x_{j}\right|+\sum_{j=1}^{n}\left|y_{j}\right|=\|\vec{x}\|_{1}+\|\vec{y}\|_{1} .
$$

d. Triangle inequality for $p>1$ :

This is Minkowski's inequality for $\left(\mathbb{R}^{n},\|\cdot\|_{p}\right)$ (thm.9.6 below). That $\|\cdot\|_{2}$ satisfies the triangle inequality (i.e., $p=2$ ) also follows independently from cor.9.1 on p.172.

Proposition 9.15 (The $L^{p}$-norm is a norm). Let $p \in[1, \infty[$ and let $a, b \in \mathbb{R}$ such that $a<b$. Then the $L^{p}$-norm $f \mapsto\|f\|_{L^{p}}=\left(\int_{a}^{b}|f(x)|^{p}\right)^{1 / p}$ is a norm in $\mathscr{C}([a, b], \mathbb{R})$.

## PROOF:

a. Positive definiteness:

Follows from prop. 9.12 on $p .174$ and the fact that $x \mapsto|f(x)|^{p}$ is a non-negative and continuous function.
b. Absolute homogeneity:

If $\lambda \in \mathbb{R}$ then

$$
\begin{aligned}
\|(\lambda f)\|_{L^{p}} & =\left(\int_{a}^{b}(|\lambda||f(x)|)^{p} d x\right)^{1 / p} \\
& =\left(|\lambda|^{p} \int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}=|\lambda|\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}=|\lambda|\|f\|_{L^{p}}
\end{aligned}
$$

c. Triangle inequality for $p=1$ :

It follows from $|f(x)+g(x)| \leqq|f(x)|+|g(x)|$ for all $x$ that

$$
\begin{aligned}
\|f+g\|_{L^{1}} & =\int_{a}^{b}|f(x)+g(x)| d x \leqq \int_{a}^{b}(|f(x)|+|g(x)|) d x \\
& =\int_{a}^{b}|f(x)| d x+\int_{a}^{b}|g(x)| d x=\|\vec{x}\|_{1}+\|\vec{y}\|_{1}
\end{aligned}
$$

d. Triangle inequality for $p>1$ :

This is Minkowski's inequality for $L^{p}-n o r m s$ (thm. 9.4 below). That $\|\cdot\|_{L^{2}}$ satisfies the triangle inequality (i.e., $p=2$ ) also follows independently from cor.9.1 on p.172.

We were referring to Minkowksi's inequalities for $\left(\mathbb{R}^{n},\|\cdot\|_{p}\right)$ and $L^{p}$-norms when proving the triangle inequality for those norms. We now build the machinery that will allows us to prove those inequalities.
Proposition 9.16 (Young's Inequality). Let $a, b>0$ and let $p, q>1$ be conjugate indices, i.e.,

$$
\begin{equation*}
\frac{1}{p}+\frac{1}{q}=1 \tag{9.38}
\end{equation*}
$$

Then the following holds Young's inequality:

$$
\begin{equation*}
a b \leqq \frac{a^{p}}{p}+\frac{b^{q}}{q} . \tag{9.39}
\end{equation*}
$$

PROOF:
Step 1: We show that $q-1=\frac{1}{p-1}$ :

$$
\begin{align*}
\frac{1}{p}+\frac{1}{q}=1 & \Rightarrow q+p=p q \Rightarrow q(1-p)=-p  \tag{9.40}\\
& \Rightarrow q=\frac{p}{p-1} \Rightarrow q-1=\frac{p-(p-1)}{p-1}=\frac{1}{p-1} .
\end{align*}
$$

Step 2: The functions

$$
\varphi:\left[0, \infty\left[\rightarrow \left[0, \infty\left[; \quad x \mapsto x^{p-1} \quad \text { and } \quad \psi:\left[0, \infty\left[\rightarrow \left[0, \infty\left[; \quad y \mapsto y^{q-1}\right.\right.\right.\right.\right.\right.\right.\right.
$$

are inverse to each other because we have

$$
\psi(\varphi(x))=\psi\left(x^{p-1}\right)=\left(x^{p-1}\right)^{q-1} \stackrel{(\star)}{=}\left(x^{p-1}\right)^{1 /(p-1)}=x
$$

(( $\star$ ) follows from step 1). We further have

$$
\varphi(\psi(y))=\varphi\left(y^{q-1}\right)=\left(y^{q-1}\right)^{p-1} \stackrel{(\star \star)}{=}\left(y^{q-1}\right)^{1 /(q-1)}=y
$$

(( $* *$ ) again follows from step 1). Note that those two functions are continuous (actually, differentiable) and strictly increasing because $\varphi^{\prime}(t)=(p-1) t^{(p-2)}>0$ and $\psi^{\prime}(t)=(q-1) t^{(q-2)}>0$ for all $t \geqq 0$. We further have $\varphi(0)=0=\psi(0)$.
Step 3: Let $f:[0, \infty[\rightarrow[0, \infty[$ be a continuous and strictly increasing (hence invertible) function such that $f(0)=0$. Then the following is true for any two real numbers $a, b>0$ :

$$
\begin{equation*}
a b \leqq \int_{0}^{a} f(x) d x+\int_{0}^{b} f^{-1}(y) d y \tag{9.41}
\end{equation*}
$$

To prove this, we distinguish three cases. Either $b<f(a)$ or $b>f(a)$ or $b=f(a)$.

The picture to the right shows what happens if $b<f(a)$ : The rec angle $a b$ is covered by the areas determined by the two integrals, $b$ not all of the area of $\int_{0}^{a} f(x) d x$ is covered by the rectangle.
Source: https://brilliant.org/wiki/youngs-inequality/


If $b>f(a)$ then the situation is similar, except that now not all of the area of $\int_{0}^{b} f^{-1}(y) d y$ is covered by the rectangle $a b$. Finally, if $b=f(a)$, the area covered by the two integrals matches the rectangle.

Step 4: We now apply the above to the function $y=f(x)=x^{p-1}$.
The inverse function is $x=f^{-1}(y)=y^{1 /(p-1)}=y^{q-1}$ (see (9.40)). We integrate and obtain

$$
\int_{0}^{a} f(x) d x=\int_{0}^{a} x^{p-1}=\left.\frac{x^{p}}{p}\right|_{0} ^{a}=\frac{a^{p}}{p}, \quad \int_{0}^{a} f-1(y) d y=\int_{0}^{a} y^{q-1}=\left.\frac{y^{q}}{q}\right|_{0} ^{b}=\frac{b^{q}}{q} .
$$

Young's inequality (9.39) now follows from (9.41).
Theorem 9.3 (Hoelder's inequality for $L^{p}$-norms). Let $a, b \in \mathbb{R}$ such that $a<b$. Let $p, q>1$ be conjugate indices, i.e.,

$$
\begin{equation*}
\frac{1}{p}+\frac{1}{q}=1 \tag{9.42}
\end{equation*}
$$

Then Hoelder's inequality is true:

$$
\begin{equation*}
\|f g\|_{L^{1}} \leqq\|f\|_{L^{p}}\|g\|_{L^{q}}, \text { i.e., } \int_{a}^{b}|f(x) g(x)| d x \leqq\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}\left(\int_{a}^{b}|g(x)|^{q} d x\right)^{1 / q} \tag{9.43}
\end{equation*}
$$

PROOF: We note that the composite function $x \mapsto|f(x)|^{p}$ is continuous, hence integrable, as the composite of the three continuous functions $x \mapsto f(x), y \mapsto|y|$, and $z \mapsto z^{p}$.

Note that $\|f\|_{L^{p}}=0$ is only possible if $|f(x)|^{p}=0$, i.e., $f(x)=0$ for all $x$ (see prop.9.12 on $p .174$ ). Likewise, $\|g\|_{L^{q}}=0$ implies $g(x)=0$ for all $x$. In either case, $\int_{a}^{b} f(x) g(x)=0$ and (9.43) is trivially satisfied. So we may assume that both $\|f\|_{L^{p}}>0$ and $\|g\|_{L^{q}}>0$ For some fixed $x \in[a, b]$ let

$$
A:=\|f\|_{L^{p}}, \quad a_{x}:=\frac{|f(x)|}{A}, \quad B:=\|g\|_{L^{q}}, \quad b_{x}:=\frac{|g(x)|}{B} .
$$

It follows from Young's inequality (9.39) that $a_{x} b_{x} \leqq \frac{a_{x}^{p}}{p}+\frac{b_{x}^{q}}{q}$. We integrate both sides of that inequality $\int_{a}^{b} \cdots d x$ and obtain from the monotonicity of the integral (see example 4.19 on $p .78$ ) that

$$
\begin{equation*}
\int_{a}^{b} a_{x} b_{x} d x \leqq \int_{a}^{b}\left(\frac{a_{x}^{p}}{p}+\frac{b_{x}^{q}}{q}\right) d x \tag{9.44}
\end{equation*}
$$

i.e.,

$$
\begin{align*}
\frac{1}{A B} \int_{a}^{b}|f(x) g(x)| d x & \leqq \int_{a}^{b}\left(\frac{|f(x)|^{p}}{p A^{p}}+\frac{|g(x)|^{q}}{q B^{q}}\right) d x  \tag{9.45}\\
& =\frac{1}{p A^{p}} \int_{a}^{b}|f(x)|^{p} d x+\frac{1}{q B^{q}} \int_{a}^{b}|g(x)|^{q} d x .
\end{align*}
$$

We use

$$
\begin{equation*}
\int_{a}^{b}|f(x)|^{p} d x=\left(\|f\|_{L^{p}}\right)^{p}, \quad \int_{a}^{b}|g(x)|^{q} d x=\left(\|g\|_{L^{q}}\right)^{q} \tag{9.46}
\end{equation*}
$$

in (9.45) and obtain

$$
\frac{1}{A B} \int_{a}^{b}|f(x) g(x)| d x \leqq \frac{A^{p}}{p A^{p}}+\frac{B^{q}}{q B^{q}}=\frac{1}{p}+\frac{1}{q}=1 .
$$

It follows from the definition of $A$ and $B$ that

$$
\int_{a}^{b}|f(x) g(x)| d x \leqq A B=\|f\|_{L^{p}}\|g\|_{L^{q}}
$$

Theorem 9.4 (Minkowski's inequality for $L^{p}$-norms). Let $a, b \in \mathbb{R}$ such that $a<b$ and let $p \in[1, \infty[$. Then Minkowski's inequality is true:

$$
\begin{align*}
\|f+g\|_{L^{p}} & \leqq\|f\|_{L^{p}}+\|g\|_{L^{p}, \text { i.e., }}  \tag{9.47}\\
\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / p} & \leqq\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}+\left(\int_{a}^{b}|g(x)|^{p} d x\right)^{1 / p} \tag{9.48}
\end{align*}
$$

PROOF: This follows for $p=1$ from part $\boldsymbol{c}$ of the proof of prop.9.15. We may assume that $p>1$. Let $q$ be the conjugate index to $p$, i.e.,

$$
\begin{equation*}
\frac{1}{p}+\frac{1}{q}=1, \text { hence }(p-1) q=p \tag{9.49}
\end{equation*}
$$

(see (9.40)). Let $a \leqq x \leqq b$. Then
$|f(x)+g(x)|^{p}=|f(x)+g(x)||f(x)+g(x)|^{p-1} \leqq|f(x)||f(x)+g(x)|^{p-1}+|g(x)||f(x)+g(x)|^{p-1}$.
The last inequality follows from $|f(x)+g(x)| \leqq|f(x)|+|g(x)|$ and $|f(x)+g(x)|^{p-1} \geqq 0$. We integrate and obtain

$$
\begin{equation*}
\int_{a}^{b}|f(x)+g(x)|^{p} d x \leqq \int_{a}^{b}|f(x)||f(x)+g(x)|^{p-1} d x+\int_{a}^{b}|g(x)||f(x)+g(x)|^{p-1} d x . \tag{9.50}
\end{equation*}
$$

We apply Hoelder's inequality to the first of the two integrals on the right hand side of (9.50) and obtain

$$
\begin{align*}
\int_{a}^{b}(|f(x)|)\left(|f(x)+g(x)|^{p-1}\right) d x & \leqq\left(\int_{a}^{b}(|f(x)|)^{p} d x\right)^{1 / p}\left(\int_{a}^{b}\left(|f(x)+g(x)|^{p-1}\right)^{q} d x\right)^{1 / q} \\
& =\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}\left(\int_{a}^{b}|f(x)+g(x)|^{(p-1) q} d x\right)^{1 / q}  \tag{9.51}\\
& =\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / q}
\end{align*}
$$

The last equality results from $(p-1) q=p$ (see (9.49)). Similarly, we obtain from the second integral on the right hand side of (9.49) the following:

$$
\begin{equation*}
\int_{a}^{b}(|g(x)|)\left(|f(x)+g(x)|^{p-1}\right) d x \leqq\left(\int_{a}^{b}|g(x)|^{p} d x\right)^{1 / p}\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / q} \tag{9.52}
\end{equation*}
$$

We apply (9.51) and (9.52) to (9.50) and obtain

$$
\begin{align*}
\int_{a}^{b}|f(x)+g(x)|^{p} d x & \leqq\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / q}  \tag{9.53}\\
& +\left(\int_{a}^{b}|g(x)|^{p} d x\right)^{1 / p}\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / q}
\end{align*}
$$

Minkowski's inequality (9.48) is trivially satisfied if $\int_{a}^{b}|f(x)+g(x)|^{p} d x=0$, so we may assume that $\int_{a}^{b}|f(x)+g(x)|^{p} d x>0$. This allows us to divide each term in (9.53) by $\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / q}$. We obtain

$$
\begin{equation*}
\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1-1 / q} \leqq\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}+\left(\int_{a}^{b}|g(x)|^{p} d x\right)^{1 / p} \tag{9.54}
\end{equation*}
$$

Note that $1-\frac{1}{q}=\frac{1}{p}$ because $\frac{1}{q}+\frac{1}{p}=1$, and (9.54) reads

$$
\left(\int_{a}^{b}|f(x)+g(x)|^{p} d x\right)^{1 / p} \leqq\left(\int_{a}^{b}|f(x)|^{p} d x\right)^{1 / p}+\left(\int_{a}^{b}|g(x)|^{p} d x\right)^{1 / p}
$$

Theorem 9.5 (Hoelder's inequality for the $p$-norms). Let $n \in \mathbb{N}$ and $\vec{x}=\left(x_{1}, \ldots x_{N}\right), \vec{y}=\left(y_{1}, \ldots y_{N}\right) \in \mathbb{R}^{n}$. Let $p, q>1$ be conjugate indices, i.e.,

$$
\begin{equation*}
\frac{1}{p}+\frac{1}{q}=1 \tag{9.55}
\end{equation*}
$$

Then Hoelder's inequality in $\mathbb{R}^{n}$ is true:

$$
\begin{equation*}
\sum_{j=1}^{n}\left|x_{j} y_{j}\right| \leqq\|\vec{x}\|_{p}\|\vec{y}\|_{q} \text {, i.e., } \sum_{j=1}^{n}\left|x_{j} y_{j}\right| \leqq\left(\sum_{j=1}^{n}\left|x_{j}\right|^{p}\right)^{1 / p}\left(\sum_{j=1}^{n}\left|y_{j}\right|^{q}\right)^{1 / q} \tag{9.56}
\end{equation*}
$$

PROOF: Let $\vec{x}, \vec{y} \in \mathbb{R}^{n}$. If $\vec{x}=0$ or $\vec{y}=0$ then $\sum_{j=1}^{n}\left|x_{j} y_{j}\right|=0$ and (9.56) is trivially satisfied. We hence may assume that both $\vec{x} \neq 0$ and $\vec{y} \neq 0$.

It follows from part $\boldsymbol{a}$ of the proof of prop.9.14 (positive definiteness of $\|\cdot\|_{p}$ for all $p$ ) on $p .175$ that $\|\vec{x}\|_{p}>0$ and $\|\vec{y}\|_{q}>0$ For some fixed index $1 \leqq j \leqq n$ let

$$
A:=\|\vec{x}\|_{p}, \quad a_{j}:=\frac{\left|x_{j}\right|}{A}, \quad B:=\|\vec{y}\|_{q}, \quad b_{x}:=\frac{\left|y_{j}\right|}{B} .
$$

It follows from Young's inequality (9.39) that

$$
a_{j} b_{j} \leqq \frac{a_{j}^{p}}{p}+\frac{b_{j}^{q}}{q} .
$$

We take sums $\sum_{j=1}^{n} \cdots$ of both sides of that inequality and obtain from the monotonicity of summation

$$
\begin{equation*}
\sum_{j=1}^{n} a_{j} b_{j} \leqq \sum_{j=1}^{n}\left(\frac{\left(a_{j}\right)^{p}}{p}+\frac{\left(b_{j}\right)^{q}}{q}\right) \tag{9.57}
\end{equation*}
$$

i.e.,

$$
\begin{equation*}
\frac{1}{A B} \sum_{j=1}^{n}\left|x_{j} y_{j}\right| \leqq \sum_{j=1}^{n}\left(\frac{\left|a_{j}\right|^{p}}{p A^{p}}+\frac{\left|b_{j}\right|^{q}}{q B^{q}}\right)=\frac{1}{p A^{p}} \sum_{j=1}^{n}\left|a_{j}\right|^{p}+\frac{1}{q B^{q}} \sum_{j=1}^{n}\left|b_{j}\right|^{q} \tag{9.58}
\end{equation*}
$$

But

$$
\begin{equation*}
\sum_{j=1}^{n}\left|x_{j}\right|^{p}=\left(\|f\|_{p}\right)^{p}, \quad \sum_{j=1}^{n}\left|y_{j}\right|^{q}=\left(\|g\|_{q}\right)^{q} . \tag{9.59}
\end{equation*}
$$

It follows from (9.58) that

$$
\frac{1}{A B} \sum_{j=1}^{n}\left|x_{j} y_{j}\right| \leqq \frac{A^{p}}{p A^{p}}+\frac{B^{q}}{q B^{q}}=\frac{1}{p}+\frac{1}{q}=1
$$

and we deduce from the definition of $A$ and $B$ that

$$
\sum_{j=1}^{n}\left|x_{j} y_{j}\right| \leqq A B=\|\vec{x}\|_{p}\|\vec{x}\|_{q}
$$

Theorem 9.6 (Minkowski's inequality for $\left(\mathbb{R}^{n},\|\cdot\|_{p}\right)$ ). Let $n \in \mathbb{N}$ and $\vec{x}=\left(x_{1}, \ldots x_{N}\right)$, $\vec{y}=\left(y_{1}, \ldots y_{N}\right) \in \mathbb{R}^{n}$. Let $p \in\left[1, \infty\left[\right.\right.$. Then Minkowski's inequality for $\left(\mathbb{R}^{n},\|\cdot\|_{p}\right)$ is true:

$$
\begin{align*}
\|\vec{x}+\vec{y}\|_{p} & \leqq\|\vec{x}\|_{p}+\|\vec{y}\|_{p}, \text { i.e., }  \tag{9.60}\\
\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / p} & \leqq\left(\sum_{j}|f(x)|^{p}\right)^{1 / p}+\left(\sum_{j}|g(x)|^{p}\right)^{1 / p} . \tag{9.61}
\end{align*}
$$

PROOF: This follows for $p=1$ from part $c$ of the proof of prop.9.15. We hence may assume that $p>1$. Let $q$ be the conjugate index to $p$, i.e.,

$$
\begin{equation*}
\frac{1}{p}+\frac{1}{q}=1, \text { hence }(p-1) q=p \tag{9.62}
\end{equation*}
$$

(see (9.40)). Let $a \leqq x \leqq b$. Then

$$
\left|x_{j}+y_{j}\right|^{p}=\left|x_{j}+y_{j}\right|\left|x_{j}+y_{j}\right|^{p-1} \leqq\left|x_{j}\right|\left|x_{j}+y_{j}\right|^{p-1}+\left|y_{j}\right|\left|x_{j}+y_{j}\right|^{p-1}
$$

The last inequality follows from $\left|x_{j}+y_{j}\right| \leqq\left|x_{j}\right|+\left|y_{j}\right|$ and $\left|x_{j}+y_{j}\right|^{p-1} \geqq 0$. We sum and obtain

$$
\begin{equation*}
\sum_{j}\left|x_{j}+y_{j}\right|^{p} \leqq \sum_{j}\left|x_{j}\right|\left|x_{j}+y_{j}\right|^{p-1}+\sum_{j}\left|y_{j}\right|\left|x_{j}+y_{j}\right|^{p-1} . \tag{9.63}
\end{equation*}
$$

Hoelder's inequality applied to the first of the two integrals on the right hand side of (9.63) yields

$$
\begin{align*}
\sum_{j}\left(\left|x_{j}\right|\right)\left(\left|x_{j}+y_{j}\right|^{p-1}\right) & \leqq\left(\sum_{j}\left(\left|x_{j}\right|\right)^{p}\right)^{1 / p}\left(\sum_{j}\left(\left|x_{j}+y_{j}\right|^{p-1}\right)^{q}\right)^{1 / q} \\
& =\left(\sum_{j}\left|x_{j}\right|^{p}\right)^{1 / p}\left(\sum_{j}\left|x_{j}+y_{j}\right|^{(p-1) q}\right)^{1 / q}  \tag{9.64}\\
& =\left(\sum_{j}\left|x_{j}\right|^{p}\right)^{1 / p}\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / q}
\end{align*}
$$

The last equality results from $(p-1) q=p$ (see (9.62)). Similarly, we obtain from the second integral on the right hand side of (9.63) the following:

$$
\begin{equation*}
\sum_{j}\left(\left|y_{j}\right|\right)\left(\left|x_{j}+y_{j}\right|^{p-1}\right) \leqq\left(\sum_{j}\left|y_{j}\right|^{p}\right)^{1 / p}\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / q} \tag{9.65}
\end{equation*}
$$

We apply (9.64) and (9.65) to (9.63) and obtain

$$
\begin{align*}
\sum_{j}\left|x_{j}+y_{j}\right|^{p} & \leqq\left(\sum_{j}\left|x_{j}\right|^{p}\right)^{1 / p}\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / q} \\
& +\left(\sum_{j}\left|y_{j}\right|^{p}\right)^{1 / p}\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / q} \tag{9.66}
\end{align*}
$$

Minkowski's inequality (9.61) is trivially satisfied if $\sum_{j}\left|x_{j}+y_{j}\right|^{p}=0$, so we may assume that $\sum_{j}\left|x_{j}+y_{j}\right|^{p}>0$. We divide each term in (9.66) by $\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / q}$ and obtain

$$
\begin{equation*}
\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1-1 / q} \leqq\left(\sum_{j}\left|x_{j}\right|^{p}\right)^{1 / p}+\left(\sum_{j}\left|y_{j}\right|^{p}\right)^{1 / p} \tag{9.67}
\end{equation*}
$$

Note that $1-\frac{1}{q}=\frac{1}{p}$ because $\frac{1}{q}+\frac{1}{p}=1$, and (9.67) reads

$$
\left(\sum_{j}\left|x_{j}+y_{j}\right|^{p}\right)^{1 / p} \leqq\left(\sum_{j}\left|x_{j}\right|^{p}\right)^{1 / p}+\left(\sum_{j}\left|y_{j}\right|^{p}\right)^{1 / p}
$$

### 9.3 Exercises for Ch. 9

Exercise 9.1. Prove prop.9.10 on p.170: Let $X$ be an arbitrary, nonempty set. Then the function $\|\cdot\|_{\infty}: \mathscr{B}(X, \mathbb{R}) \rightarrow \mathbb{R}_{+}, \quad h \rightarrow\|h\|_{\infty}=\sup \{|h(x)|: x \in X\}$ defines a norm.

Exercise 9.2. Prove prop.9.7 (Algebraic properties of the inner product) on p.169:
Let $V$ be a vector space with inner product $\bullet(\cdot, \cdot)$. Let $a, b, x, y \in V$. Then
a. $(a+b) \bullet(x+y)=a \bullet x+b \bullet x+a \bullet y+b \bullet y$
b. $(x+y) \bullet(x+y)=x \bullet x+2(x \bullet y)+y \bullet y$
c. $(x-y) \bullet(x-y)=x \bullet x-2(x \bullet y)+y \bullet y$

Exercise 9.3. Prove prop.9.11 on p.171:
Let $(V,\|\cdot\|)$ be a normed vector space and let $\gamma>0$. Let $p: V \rightarrow \mathbb{R}$ be defined as $p(x):=\gamma\|x\|$. Then $p$ also is a norm.

Exercise 9.4. Prove that the $p$-norm (see def.9.15 on p .171 ) is a norm on $\mathbb{R}^{n}$ for the special case $p=1$ :

$$
\|\vec{x}\|_{1}=\sum_{j=1}^{n}\left|x_{j}\right|
$$

In chapter 10 on the topology of real numbers ( $p$. 183) you will learn about metric spaces as a concept that generalizes the measurement of distance (or closeness, if you prefer) for the elements of a nonempty set.

## 10 Metric Spaces

There is a branch of Mathematics, called topology, which deals with the concept of closeness. The concept of limits of a sequence $\left(x_{n}\right)_{n}$ is based on closeness: The points of the sequence must get "arbitrarily close" to its limit as $n \rightarrow \infty$. Continuity of functions also can be phrased in terms of closeness: They map arbitrarily close elements of the domain to arbitrarily close elements of the codomain. In the most general setting Topology is about neighborhoods of a point without having the concept of measuring the distance of two points. We mostly won't deal with such a level of generality in this document. Instead we'll we'll focus on metric spaces $(X, d)$ : sets $X$ that are equipped with a distance function $(x, y) \mapsto d(x, y)$. Even this limited context will significantly generalize the material of ch.8.2 (Convergence and Continuity in $\mathbb{R}$ ) and ch.9.2.2 (Normed Vector Spaces).

### 10.1 The Topology of Metric Spaces

A metric is a real valued function of two arguments which associates with any two points $x, y \in X$ their "distance" $d(x, y)$.

It is clear how you measure the distance (or closeness, depending on your point of view) of two numbers $x$ and $y$ : you plot them on an $x$-axis where the distance between two consecutive integers is exactly one inch, grab a ruler and see what you get. Alternate approach: you compute the difference. For example, the distance between $x=12.3$ and $y=15$ is $x-y=12.3-15=-2.7$. Actually, we have a problem: There are situations where direction matters and a negative distance is one that goes into the opposite direction of a positive distance, but we do not want that in this context and understand the distance to be always non-negative, i.e.,

$$
\operatorname{dist}(x, y)=|y-x|=|x-y|
$$

More importantly, you must forget what you learned in your in your science classes: "Never ever talk about a measure (such as distance or speed or volume) without clarifying its dimension". Is the speed measured in miles per hour our inches per second? Is the distance measured in inches or miles or micrometers? In the context of metric spaces we measure distance simply as a number, without any dimension attached to it. For the above example, you get

$$
\operatorname{dist}(12.3,15)=|12.3-15|=2.7
$$

In section 9.1.3 on $p .157$ it is shown in great detail that the distance between two two-dimensional vectors $\vec{v}=\left(v_{1}, v_{2}\right)$ and $\vec{w}=\left(w_{1}, w_{2}\right)$ is

$$
\operatorname{dist}(\vec{v}, \vec{w})=\sqrt{\left(w_{1}-v_{1}\right)^{2}+\left(w_{2}-v_{2}\right)^{2}}
$$

and the distance between two three-dimensional vectors $\vec{v}=\left(v_{1}, v_{2}, v_{3}\right)$ and $\vec{w}=\left(w_{1}, w_{2}, w_{3}\right)$ is

$$
\operatorname{dist}(\vec{v}, \vec{w})=\sqrt{\left(w_{1}-v_{1}\right)^{2}+\left(w_{2}-v_{2}\right)^{2}+\left(w_{3}-v_{3}\right)^{2}}
$$

In the next chapter we will generalize the concept of distance to more general objects.

### 10.1.1 Definition and Examples of Metric Spaces

Definition 10.1 (Metric spaces). Let $X$ be an arbitrary, nonempty set.
A metric on $X$ is a real valued function of two arguments

$$
d(\cdot, \cdot): X \times X \rightarrow \mathbb{R}, \quad(x, y) \mapsto d(x, y)
$$

with the following three properties: ${ }^{87}$

```
(10.1a) \(d(x, y) \geqq 0 \quad \forall x, y \in X \quad\) and \(\quad d(x, y)=0 \quad \Leftrightarrow x=y\)
positive definiteness
(10.1b) \(d(x, y)=d(y, x) \quad \forall x, y \in X\)
(10.1c) \(d(x, z) \leqq d(x, y)+d(y, z) \quad \forall x, y, z \in X \quad\) triangle inequality
symmetry
```

Let $x, y \in X$. We say that $x$ and $y$ are $\varepsilon$-close if $d(x, y)<\varepsilon$.
The pair $(X, d(\cdot, \cdot)$, usually just written as $(X, d)$, is called a metric space. We'll write $X$ for short if it is clear which metric we are talking about.

To appreciate that last sentence, you must understand that there can be more than one metric on $X$. See the examples below.

Remark 10.1 (Metric properties). Let us quickly examine what those properties mean.
"Positive definite": The distance is never negative and two items $x$ and $y$ have distance zero if and only if they are equal.
"symmetry": the distance from $x$ to $y$ is no different to that from $y$ to $x$. That may come as a surprise to you if you have learned in Physics about the distance from point $a$ to point $b$ being the vector $\vec{v}$ that starts in $a$ and ends in $b$ and which is the opposite of the vector $\vec{w}$ that starts in $b$ and ends in $a$, i.e., $\vec{v}=-\vec{w}$. We only care about size and not about direction.
"Triangle inequality": If you directly drive from $x$ to $z$ then this will take less fuel than if you make a stopover at an intermediary $y$.

Remark 10.2. Do not make the mistake and think of $X$ as a set of numbers or vectors! For example, we might deal with

$$
X:=\{\text { all students who are currently taking this class }\} .
$$

We can define the distance of any two students $s_{1}$ and $s_{2}$ as

$$
d\left(s_{1}, s_{2}\right)= \begin{cases}0 & \text { for } s_{1}=s_{2} \\ 1 & \text { for } s_{1} \neq s_{2}\end{cases}
$$

We will learn later in this subchapter that the above function is called the discrete metric on $X$ and satisfied indeed the definition of a metric. ${ }^{88}$

The triangle inequality generalizes to more than two terms.
Proposition 10.1. Let $(X, d)$ be a metric space. Let $n \in \mathbb{N}$ and $x_{1}, x_{2}, \ldots, x_{n} \in X$. Then

$$
\begin{equation*}
d\left(x_{1}, x_{n}\right) \leqq \sum_{j=1}^{n-1} d\left(x_{j}, x_{j+1}\right)=d\left(x_{1}, x_{2}\right)+d\left(x_{2}, x_{3}\right)+d\left(x_{n-1}, x_{n}\right) . \tag{10.2}
\end{equation*}
$$

[^66]The proof is left as exercise 10.1.
Before we give some examples of metric spaces, here is a theorem that tells you that a vector space with a norm (see def.9.14 on p.171), becomes a metric space as follows:

Theorem 10.1 (Norms define metric spaces). Let $(V,\|\cdot\|)$ be a normed vector space. Then the function

$$
\begin{equation*}
d_{\|\cdot\|}(\cdot, \cdot): V \times V \rightarrow \mathbb{R}_{\geqq 0} ; \quad(x, y) \mapsto d_{\|\cdot\|}(x, y):=\|y-x\| \tag{10.3}
\end{equation*}
$$

defines a metric space ( $V, d_{\|\cdot\|}$ ).
The proof is left as exercise 10.2.
Definition 10.2 (Metric induced by a norm). We say that the metric $d_{\|\cdot\|}(\cdot, \cdot)$ defined by (10.3) is induced by the norm $\|\cdot\|$. We also say that $d_{\|\cdot\|}(\cdot, \cdot)$ is derived from the norm $\|\cdot\|$ or that $d_{\|\cdot\|}(\cdot, \cdot)$ is associated with the norm $\|\cdot\| .{ }^{89}$

Here are some examples of metric spaces.
Example $10.1\left(\left(\mathbb{R}\right.\right.$ with $\left.d_{|\cdot|}(a, b)=|b-a|\right)$. According to thm.10.1 $\left(\mathbb{R}, d_{|\cdot|}\right)$ is a metric space because the Euclidean norm $|\cdot|$ is a norm on $\mathbb{R}=\mathbb{R}^{1}$.

Here is a direct proof; It is obvious that if $x, y$ are real numbers then the difference $x-y$, and hence its absolute value, is zero if and only if $x=y$ and that proves positive definiteness.

Symmetry follows from $d_{|\cdot|}(x, y)=|x-y|=|-(y-x)|=|y-x|=d_{|\cdot|}(y, x)$.
The triangle inequality for a metric follows from $|a+b| \leqq|a|+|b|$ (see prop.2.2 on p.19):

$$
\begin{aligned}
d_{|\cdot|}(x, z) & =|x-z|=|(x-y)-(z-y)| \\
& \leqq|x-y|+|z-y|=d_{|\cdot|}(x, y)+d_{|\cdot|}(z, y)=d_{|\cdot|}(x, y)+d_{|\cdot|}(y, z)
\end{aligned}
$$

Example 10.2 (bounded real valued functions with $\left.d_{\|\cdot\|_{\infty}} f, g\right)=\sup -$ norm of $g(\cdot)-f(\cdot)$ ).

$$
\begin{equation*}
d_{\|\cdot\|_{\infty}}(f, g)=\|g-f\|_{\infty}=\sup \{|g(x)-f(x)|: x \in X\} \tag{10.4}
\end{equation*}
$$

is a metric on the set $\mathscr{B}(X, \mathbb{R})$ of all bounded real valued functions on $X$. This follows from thm.10.1 and prop. 9.10 on p. 170, according to which $\left(\mathscr{B}(X, \mathbb{R}),\|\cdot\|_{\infty}\right)$ is a normed vector space.
Example 10.3 (continuous real valued functions on [a,b]with $d_{\|\cdot\|_{L^{2}}}(f, g)=\|g-f\|_{L^{2}}$ ). We will see in ch.10.1.2 on p. 186 that $\|g-f\|_{\infty}$ is a good measure for the difference of the functions $f$ and $g$ and that an often even better measure is that of the area difference between their graphs which is given by the netric

$$
\begin{equation*}
d_{\|\cdot\|_{L^{2}}}(f, g)=\|g-f\|_{L^{2}}=\sqrt{\int_{a}^{b}(g(x)-f(x))^{2} d x} \tag{10.5}
\end{equation*}
$$

(see def.9.18 on p.175).

[^67]Example 10.4 ( $\mathbb{R}^{N}$ with the Euclidean metric).

$$
d_{\|\cdot\|_{2}}(\vec{x}, \vec{y})=\sqrt{\left(y_{1}-x_{1}\right)^{2}+\left(y_{2}-x_{2}\right)^{2}+\ldots+\left(y_{N}-x_{N}\right)^{2}}=\sqrt{\sum_{j=1}^{N}\left(y_{j}-x_{j}\right)^{2}}
$$

This follows from the fact that the Euclidean norm is a norm on the vector space $\mathbb{R}^{N}$ (see cor.9.1 on p .172 .)

Just in case you think that all metrics are derived from norms, here is a counterexample.
Definition 10.3 (Discrete metric). Let $X$ be nonempty. Then the function

$$
d(x, y)= \begin{cases}0 & \text { for } x=y \\ 1 & \text { for } x \neq y\end{cases}
$$

on $X \times X$ is called the discrete metric on X .
The above definition makes sense because of the following proposition.
Proposition 10.2. The discrete metric satisfies the properties of a metric.
PROOF: Obviously the function is non-negative and it is zero if and only if $x=y$. Symmetry is obvious too.

The triangle inequality $d(x, z) \leqq d(x, y)+d(y, z)$ is clear in the special case $x=z$. (Why?)
So let us assume $x \neq z$. But then $x \neq y$ or $y \neq z$ or both must be true. (Why?) That means that

$$
d(x, z)=1 \leqq d(x, y)+d(y, z)
$$

and this proves the triangle inequality.

### 10.1.2 Measuring the Distance of Real Valued Functions

How do we compare two functions? Let us make our lives easier: How do we compare two real valued functions $f(\cdot)$ and $g(\cdot)$ ? One answer is to look at a picture with the graphs of $f(\cdot)$ and $g(\cdot)$ and look at the shortest distance $|f(x)-g(x)|$ as you run through all $x$. That means that the distance between the functions $f(x)=x$ and $g(x)=x^{2}$ is zero because $f(1)=g(1)=1$. The distance between $f(x)=x+1$ and $g(x)=0$ (the $x$-axis) is also zero because $f(-1)=g(-1)=0$.

Do you really think this is a good way to measure closeness? You really do not want two items to have zero distance unless they coincide. It's a lot better to look for an argument $x$ where the value $|f(x)-g(x)|$ is largest rather than smallest. Now we are ready for a proper definition.

Definition 10.4 (Maximal displacement distance between real valued functions). Let $X$ be an arbitrary, nonempty set and let $f(\cdot), g(\cdot): X \rightarrow \mathbb{R}$ be two real valued functions on $X$. We define the maximal displacement distance, also called the sup-norm distance or $\|\cdot\|_{\infty}$ distance, between $f(\cdot)$ and $g(\cdot)$ as

$$
\begin{equation*}
d_{\infty}(f, g):=\|f(\cdot)-g(\cdot)\|_{\infty}=\sup \{|f(x)-g(x)|: x \in X\} \tag{10.6}
\end{equation*}
$$

i.e., as the metric induced by the sup-norm on the set $\mathscr{B}(X, \mathbb{R})$ of all bounded real valued function on $X$.

Remark 10.3. We have previously encountered formula (10.6) in example 10.2 on p.185. We will see in prop. 10.27 on p. 226 of ch.10.3.1 on convergence of function sequences that the sup-norm induced metric is suitable to measure what will be called "uniform convergence" of real valued functions. As a metric, the distance measure of two functions $f, g$ satisfies positive definiteness, symmetry and the triangle inequality. We have seen in other contexts what those properties mean.
"Positive definite": The distance is never negative and two functions $f(\cdot)$ and $g(\cdot)$ have distance zero if and only if they are equal, i.e., if and only if $f(x)=g(x)$ for each argument $x \in X$.
"Symmetry": the distance from $f(\cdot)$ to $g(\cdot)$ is no different than that from $g(\cdot)$ to $f(\cdot)$. Symmetry implies that you do not obtain a negative distance if you walk in the opposite direction.
"Triangle inequality": If you directly compare the maximum deviation between two functions $f(\cdot)$ and $h(\cdot)$ then this will never be more than than using an intermediary function $g(\cdot)$ and adding the distance between $f(\cdot)$ and $g(\cdot)$ to that between $g(\cdot)$ and $h(\cdot)$.

Remark 10.4. The following picture illustrates the last definition. Plot the graphs of $f$ and $g$ as usual and find the the spot $x_{0}$ on the $x$-axis for which the difference $\left|f\left(x_{0}\right)-g\left(x_{0}\right)\right|$ (the length of the vertical line that connects the two points with coordinates $\left(x_{0}, f\left(x_{0}\right)\right)$ and $\left(x_{0}, g\left(x_{0}\right)\right)$ ) has the largest possible value. The domain of $f$ and $g$ is the subset of $\mathbb{R}$ that corresponds to the thick portion of the $x$-axis.


Figure 10.1: Distance of two real valued functions.

This figure allows you to visualize for a given $\delta>0$ and $f: X \rightarrow \mathbb{R}$ the " $\delta$-neighborhood" of $f(\cdot)$ defined as

$$
\begin{equation*}
N_{\delta}(f):=\{g: X \rightarrow \mathbb{R}: d(\infty f, g)<\delta\}=\left\{g(\cdot): X \rightarrow \mathbb{R}: \sup _{x \in X}|f(x)-g(x)|<\delta\right\}, \tag{10.7}
\end{equation*}
$$

i.e., the set of all functions $g(\cdot)$ with distance less than $\delta$ from $f(\cdot)$.

You draw the graph of $f(\cdot)+\delta$ (the graph of $f(\cdot)$ shifted up by the amount of $\delta$ ) and the graph of $f(\cdot)-\delta$ (the graph of $f(\cdot)$ shifted down by the amount of $\delta$ ). Any function $g(\cdot)$ which stays completely inside this band, without actually touching it, belongs to the $\delta$-neighborhood of $f(\cdot)$.


Figure 10.2: $\delta$-neighborhood of a real valued function.

In other words, assuming that the domain $A$ is a single, connected chunk and not a collection of several separate intervals, the $\delta$-neighborhood of $f(\cdot)$ is a "band" whose contours are made up on the left and right by two vertical lines and on the top and bottom by two lines that look like the graph of $f(\cdot)$ itself but have been shifted up and down by the amount of $\delta$.
Definition 10.5 (Mean distances between real valued functions). Let $a, b \in \mathbb{R}$ such that $a<b$ and let $f(\cdot), g(\cdot): X \rightarrow \mathbb{R}$ be two continuous real valued functions on $X$. We define the mean square distance between $f(\cdot)$ and $g(\cdot)$ on $[a, b]$ as

$$
\begin{equation*}
d_{L^{2}}(f, g):=d_{\|\cdot\|_{L^{2}}(f, g)}=\|g-f\|_{L^{2}}=\int_{a}^{b}\left(g(x)-f(x)^{2} d x\right. \tag{10.8}
\end{equation*}
$$

i.e., as the metric induced by the $L^{2}$-norm on the set $\mathscr{C}_{\mathscr{B}}([a, b], \mathbb{R})$ of all continuous and bounded real valued function on $[a, b]$.

We further define the mean distance between $f(\cdot)$ and $g(\cdot)$ on $[a, b]$ as

$$
\begin{equation*}
d_{L^{1}}(f, g):=d_{\|\cdot\|_{L^{1}}(f, g)}=\|g-f\|_{L^{1}}=\int_{a}^{b} \mid g(x)-f\left(\left.x\right|^{2} d x\right. \tag{10.9}
\end{equation*}
$$

i.e., as the metric induced by the $L^{1}$-norm on the set $\mathscr{C}_{\mathscr{B}}([\mathrm{a}, \mathrm{b}], \mathbb{R})$.

Remark 10.5. We saw in def.9.17, example 9.21 and def. 9.18 on pp. 173 that both

$$
\begin{align*}
& d_{L^{1}}(f, g):=d_{\|\cdot\|_{L^{1}}(f, g)}=\|g-f\|_{L^{1}}=\int_{a}^{b}|g(x)-f(x)| d x,  \tag{10.10}\\
& d_{L^{2}}(f, g):=d_{\|\cdot\|_{L^{2}}(f, g)}=\|g-f\|_{L^{2}}=\int_{a}^{b}\left(g(x)-f(x)^{2} d x,\right. \tag{10.11}
\end{align*}
$$

are often better suitable than the distance derived from the sup-norm to measure the distance of two functions. One of the drawbacks from a teaching perspective is that there is no picture like figure 10.2 to visualize the set of all functions with an $L^{1}$-distance or $L^{2}$-distance from a given function.

### 10.1.3 Neighborhoods and Open Sets

A. Given a point $x_{0} \in \mathbb{R}$ (a real number) and $\varepsilon>0$, we can look at

$$
\begin{align*}
N_{\varepsilon}\left(x_{0}\right)=\left(x_{0}-\varepsilon, x_{0}+\varepsilon\right) & =\left\{x \in \mathbb{R}: x_{0}-\varepsilon<x<x_{0}+\varepsilon\right\}  \tag{10.12}\\
& =\left\{x \in \mathbb{R}: d\left(x, x_{0}\right)=\left|x-x_{0}\right|<\varepsilon\right\}
\end{align*}
$$

which is the set of all real numbers $x$ with a distance to $x_{0}$ of strictly less than a number $\varepsilon$ (the open interval with end points $x_{0}-\varepsilon$ and $x_{0}+\varepsilon$ ). (see example (10.1) on p.185).
B. Given a point $\vec{x}_{0}=\left(x_{0}, y_{0}\right) \in \mathbb{R}^{2}$ (a point in the $x y$-plane), we can look at

$$
\begin{align*}
N_{\varepsilon}\left(\vec{x}_{0}\right) & =\left\{(x, y) \in \mathbb{R}^{2}:\left(x-x_{0}\right)^{2}+\left(y-y_{0}\right)^{2}<\varepsilon^{2}\right\}  \tag{10.13}\\
& =\left\{\vec{x} \in \mathbb{R}^{2}: d_{\|\cdot\|_{2}}(\vec{x}, \vec{y})=\left\|\vec{x}-\vec{x}_{0}\right\|_{2}<\varepsilon\right\}
\end{align*}
$$

which is the set of all points in the plane with a distance to $\vec{x}_{0}$ of strictly less than a number $\varepsilon$ (the open disc around $\vec{x}_{0}$ with radius $\varepsilon$ from which the points on the boundary (those with distance equal to $\varepsilon$ ) are excluded).
C. Given a point $\vec{x}_{0}=\left(x_{0}, y_{0}, z_{0}\right) \in \mathbb{R}^{3}$ (a point in the 3 -dimensional space), we can look at

$$
\begin{align*}
N_{\varepsilon}\left(\vec{x}_{0}\right) & =\left\{(x, y, z) \in \mathbb{R}^{3}:\left(\vec{x}-\vec{x}_{0}\right)^{2}+\left(\vec{y}-\vec{y}_{0}\right)^{2}+\left(\vec{z}-\vec{z}_{0}\right)^{2}<\varepsilon^{2}\right\} \\
& =\left\{\vec{x} \in \mathbb{R}^{3}: d_{\|\cdot\|_{2}}(\vec{x}, \vec{y})=\left\|\vec{x}-\vec{x}_{0}\right\|_{2}<\varepsilon\right\} \tag{10.14}
\end{align*}
$$

which is the set of all points in space with a distance to $\vec{x}_{0}$ of strictly less than a number $\varepsilon$ (the open ball around $\vec{x}_{0}$ with radius $\varepsilon$ from which the points on the boundary (those with distance equal to $\varepsilon$ ) are excluded).
D. Given a normed vector space $(V,\|\cdot\|)$ and a vector $x_{0} \in V$, we can look at

$$
\begin{equation*}
N_{\varepsilon}\left(x_{0}\right)=\left\{x \in V:\left\|x-x_{0}\right\|<\varepsilon\right\} \tag{10.15}
\end{equation*}
$$

which is the set of all vectors in $V$ with a distance to $x_{0}$ of strictly less than a number $\varepsilon$ (the open set around $x_{0}$ with "radius" $\varepsilon$ from which the points on the boundary (those with distance equal to $\varepsilon$ ) are excluded).
E. Given a bounded real valued function $f \in \mathscr{B}(X, \mathbb{R})$, we can look at the sets $N_{\varepsilon}(f)(\varepsilon>0)$ defined in (10.7) on p.187, i.e., the set of all functions $g(\cdot)$ with distance less than $\varepsilon$ from $f(\cdot)$.
F. Given is a closed interval $[a, b](a, b \in \mathbb{R})$. For a continuous (hence bounded) real valued function $f \in \mathscr{B}([a, b], \mathbb{R})$, we can look at the sets

$$
\begin{equation*}
N_{\varepsilon}(f)=\left\{g \in \mathscr{B}([a, b], \mathbb{R}):\|g-f\|_{L^{2}}<\varepsilon\right\}, \tag{10.16}
\end{equation*}
$$

i.e., the set of all functions $g(\cdot)$ such that $\sqrt{\int_{a}^{b}(g(x)-f(x))^{2} d x}<\varepsilon$ (see def.9.18 on $p .175$ )

There is one more item more general than neighborhoods of elements belonging to normed vector spaces, and that would be neighborhoods in metric spaces. We have arrived at the final definition:

Definition 10.6 ( $\varepsilon$-Neighborhood). Given a metric space ( $X, d$ ) and an element $x_{0} \in X$, let

$$
\begin{equation*}
N_{\varepsilon}\left(x_{0}\right)=\left\{x \in X: d\left(x, x_{0}\right)<\varepsilon\right\} \tag{10.17}
\end{equation*}
$$

be the set of all elements of $X$ with a distance to $x_{0}$ of strictly less than the number $\varepsilon$ (the open set around $x_{0}$ with "radius" $\varepsilon$ from which the points on the boundary (those with distance equal to $\varepsilon$ ) are excluded). We call $N_{\varepsilon}\left(x_{0}\right)$ the $\varepsilon$-neighborhood of $x_{0}$.

The following should be intuitively clear: Look at any point $a \in N_{\varepsilon}\left(x_{0}\right)$. You can find $\delta>0$ such that the entire $\delta$-neighborhood $N_{\delta}(a)$ of a is contained inside $N_{\varepsilon}\left(x_{0}\right)$. Just in case you do not trust your intuition, this is shown in prop. 10.4 just a little bit further down.

It then follows that any $a \in N_{\varepsilon}\left(x_{0}\right)$ is an interior point of $N_{\varepsilon}\left(x_{0}\right)$ in the following sense:
Definition 10.7 (Interior points in metric spaces). Given is a metric space ( $X, d$ ).
An element $a \in A \subseteq X$ is called an interior point of $A$ if we can find some $\varepsilon>0$, however small it may be, so that $N_{\varepsilon}(a) \subseteq A$.

Definition 10.8 (Open sets in metric spaces). Given is a metric space ( $X, d$ ).
A set all of whose members are interior points is called an open set.
Proposition 10.3. Let $(X, d)$ be a metric space. Let $x, y \in X$ and $\varepsilon>0$ such that $y \in N_{\varepsilon}(x)$.

$$
\text { If } \delta>0 \text { Then } N_{\delta}(y) \subseteq N_{\delta+\varepsilon}(x)
$$

PROOF: Let $z \in N_{\delta}(y)$. Then

$$
d(z, x) \leqq d(z, y)+d(y, x)<\delta+\varepsilon
$$

In other words, each element $z$ of $N_{\delta}(y)$ is $\delta+\varepsilon$-close to $x$. Hence $N_{\delta}(y) \subseteq N_{\delta+\varepsilon}(x)$.
Proposition 10.4. $N_{\varepsilon}\left(x_{0}\right)$ is an open set
It is worth while to examine the following proof ${ }^{90}$ closely because you can see how the triangle inequality is put to work.
$a \in N_{\varepsilon}\left(x_{0}\right)$ means that $\varepsilon-d\left(a, x_{0}\right)>0$, say,

$$
\begin{equation*}
\varepsilon-d\left(a, x_{0}\right)=2 \delta \tag{10.18}
\end{equation*}
$$

where $\delta>0$. Let $b \in N_{\delta}(a)$. The claim is that any such $b$ is an element of $N_{\varepsilon}\left(x_{0}\right)$. How so?

$$
d\left(b, x_{0}\right) \leqq d(b, a)+d\left(a, x_{0}\right)<\delta+(\varepsilon-2 \delta)=\varepsilon-\delta<\varepsilon
$$

In the above chain, the first inequality is a consequence of the triangle inequality. The second one reflects the fact that $b \in N_{\delta}(a)$ and uses (10.18).

We have proved that for any $b \in N_{\delta}(a)$ it is true that $b \in N_{\varepsilon}\left(x_{0}\right)$ hence $N_{\delta}(a) \subseteq N_{\varepsilon}\left(x_{0}\right)$.
This proves that $a$ is an interior point of $N_{\varepsilon}\left(x_{0}\right)$. But $a$ is an arbitrary point in $N_{\varepsilon}\left(x_{0}\right)$. It follows that $N_{\varepsilon}\left(x_{0}\right)$ is open.

Placeholder for thm.10.32 on p.242:
Open intervals are open in $\left(\mathbb{R}, d_{|\cdot|}\right)$.
Definition 10.9 (Neighborhoods in Metric Spaces). Let $(X, d)$ be a metric space, $x_{0} \in X$. Any open set that contains $x_{0}$ is called an open neighborhood of $x_{0}$. Any superset of an open neighborhood of $x_{0}$ is called a neighborhood of $x_{0}$.

[^68]Remark 10.6 (Open neighborhoods are the important ones). You will see that the important neighborhoods are the small ones, not the big ones. The definition above says that for any neighborhood $A_{x}$ of a point $x \in X$ you can find an open neighborhood $U_{x}$ of $x$ such that $U_{x} \subseteq A_{x}$.

Because of this there are many propositions and theorems where you may assume that a neighborhood you deal with is open.
Definition 10.10 (boundary points). Let $A$ be a subset of the metric space ( $X, d$ ). $x \in X$ is called a boundary point of $A$ if any neighborhood of $x$ intersects both $A$ and $A^{\complement}$. We write $\partial A$ for the set of all boundary points of $A$ and call this set the boundary of $A$.
Theorem 10.2 (Metric spaces are topological spaces). The following is true about open sets of a metric space $(X, d)$ :

$$
\begin{equation*}
\text { An arbitrary union } \bigcup_{i \in I} U_{i} \text { of open sets } U_{i} \text { is open. } \tag{10.19a}
\end{equation*}
$$

A finite intersection $U_{1} \cap U_{2} \cap \ldots \cap U_{n}(n \in \mathbb{N})$ of open sets is open.
The entire set $X$ is open and the empty set $\emptyset$ is open.
PROOF of a: Let $U:=\bigcup_{i \in I} U_{i}$ and assume $x \in U$. We must show that $x$ is an interior point of $U$. An element belongs to a union if and only if it belongs to at least one of the participating sets of the union. So there exists an index $i_{0} \in I$ such that $x \in U_{i_{0}}$.

Because $U_{i_{0}}$ is open, $x$ is an interior point and we can find a suitable $\varepsilon>0$ such that $N_{\varepsilon}(x) \subseteq U_{i_{0}}$. But $U_{i_{0}} \subseteq U$, hence $N_{\varepsilon}(x) \subseteq U$. It follows that $x$ is interior point of $U$. But $x$ was an arbitrary point of $U=\bigcup_{i \in I} U_{i}$ which therefore is shown to be an open set.
PROOF of b: Let $x \in U:=U_{1} \cap U_{2} \cap \ldots \cap U_{n}$. Then $x \in U_{j}$ for all $1 \leqq j \leqq n$ according to the definition of an intersection and it is inner point of all of them because they all are open sets. Hence, for each $j$ there is a suitable $\varepsilon_{j}>0$ such that $N_{\varepsilon_{j}}(x) \subseteq U_{j}$ Now define

$$
\varepsilon:=\min \left\{\varepsilon_{1}, \varepsilon_{2}, \varepsilon_{3}, \ldots, \varepsilon_{n}\right\}
$$

Then $\varepsilon>0$ and ${ }^{91}$

$$
N_{\varepsilon}(x) \subseteq N_{\varepsilon_{j}}(x) \subseteq U_{j}(1 \leqq j \leqq n), \quad \text { hence } \quad N_{\varepsilon}(x) \subseteq \bigcap_{j=1}^{n} U_{j} .
$$

We have shown that an arbitrary $x \in U$ is interior point of $U$ and this proves part $b$.
PROOF of $c$ : First we deal with the set $X$. Choose any $x \in X$. No matter how small or big an $\varepsilon>0$ you choose, $N_{\varepsilon}(x)$ is a subset of $X$. But then $x$ is an inner point of $X$, so all members of $x$ are inner points and this proves that $X$ is open.
Now to the empty set $\emptyset$. You may have a hard time to accept the logic of this statement: All elements of $\emptyset$ are interior points. But remember, the premise "let $x \in X$ " is always false and you may conclude from it whatever you please (see ch. 3 (Logic).

This last theorem provides the underpinnings for the definition of abstract topological spaces, a subject which will be touched upon in ch.10.1.5 on p.193.

[^69]
### 10.1.4 Convergence

You have already encountered the precise definition of the convergence of sequences of real numbers in ch.8.2. It is only a small step to generalize this concept to all metric spaces and therefore also to all normed vector spaces.

Definition 10.11 (convergence of sequences in metric spaces). Given is a metric space ( $X, d$ ).
We say that a sequence $\left(x_{n}\right)$ of elements of $X$ converges to $a \in X$ for $n \rightarrow \infty$ if almost all of the $x_{n}$ will come arbitrarily close to $a$ in the following sense:

Let $\delta$ be an arbitrarily small positive real number. Then there is a (possibly extremely large) integer $n_{0}$ such that all $x_{j}$ belong to $N_{\delta}(a)$ just as long as $j \geqq n_{0}$. To say this another way: Given any number $\delta>0$, however small, you can find an integer $n_{0}$ such that

$$
\begin{equation*}
d\left(a, x_{j}\right)<\delta \text { for all } j \geqq n_{0} \tag{10.20}
\end{equation*}
$$

We write either of

$$
\begin{equation*}
a=\lim _{n \rightarrow \infty} x_{n} \quad \text { or } \quad x_{n} \rightarrow a \tag{10.21}
\end{equation*}
$$

and we call $a$ the limit of the sequence $\left(x_{n}\right)$
There is an equivalent way of expressing convergence towards $a$ : No matter how small a neighborhood of $a$ you choose: at most finitely many of the $x_{n}$ will be located outside that neighborhood.

Theorem 10.3 (Limits in metric spaces are uniquely determined). Let $(X, d)$ be a metric space . Let $\left(x_{n}\right)_{n}$ be a convergent sequence in $X$ Then its limit is uniquely determined.

PROOF: Otherwise there would be two different points $L_{1}, L_{2} \in X$ such that both $\lim _{n \rightarrow \infty} x_{n}=L_{1}$ and $\lim _{n \rightarrow \infty} x_{n}=L_{2}$ Let $\varepsilon:=d\left(L_{1}, L_{2}\right) / 2$. There will be $N_{1}, N_{2} \in \mathbb{N}$ such that

$$
d\left(x_{n}, L_{1}\right)<\varepsilon \forall n \geqq N_{1} \text { and } d\left(x_{n}, L_{2}\right)<\varepsilon \forall n \geqq N_{2} .
$$

It follows that, for $n \geqq \max \left(N_{1}, N_{2}\right),{ }^{92}$

$$
d\left(L_{1}, L_{2}\right) \leqq d\left(L_{1}, x_{n}\right)+d\left(x_{n}, L_{2}\right)<2 \varepsilon=d\left(L_{1}, L_{2}\right)
$$

and we have reached a contradiction.
The following proposition shows that the limit behavior of a sequence is a property of its tail, i.e., it does not depend on the first finitely many indices.
Proposition 10.5. ${ }^{93}$ Let $x_{n}$ be a convergent sequence in a metric space $(X, d)$ with limit $L \in E$. Let $K \in \mathbb{N}$. For $n \in \mathbb{N}$ let $y_{n}:=x_{n+K}$. Then $\lim _{n \rightarrow \infty}\left(y_{n}\right)_{n}=L$.

## PROOF:

Left as an exercise.

[^70]Proposition 10.6. ${ }^{94}$ Let $x_{n}$, $y_{n}$ be two sequences in a metric space $(X, d)$. Assume there is $K \in \mathbb{N}$ such that $x_{n}=y_{n}$ for all $n \geqq K$. Let $L \in X$ Then

$$
\lim _{n \rightarrow \infty} x_{n}=L \Leftrightarrow \lim _{n \rightarrow \infty} y_{n}=L
$$

PROOF:
Left as an exercise.

## Remark 10.7.

The following allows you to prove convergence of $x_{n}$ to $L$ in $(X, d)$ by utilizing what you know about convergence in $\left(\mathbb{R}, d_{\mid} \cdot \mid\right)$.

$$
\lim _{n \rightarrow \infty} x_{n}=L \Leftrightarrow \lim _{n \rightarrow \infty} d\left(x_{n}, L\right)=0 .
$$

Remark 10.8 (Opposite of convergence). Given a metric space $(X, d)$, what is the opposite of $\lim _{k \rightarrow \infty} x_{k}=L$ ?
Beware! It is NOT the statement that $\lim _{k \rightarrow \infty} x_{k} \neq L$ because such a statement would mislead you to believe that such a limit exists, it just happens not to coincide with $L$

The correct answer: There exists some $\varepsilon>0$ such that for all $N \in \mathbb{N}$ there exists some natural number $j=j(N)$ such that $j \geq N$ and $d\left(x_{j}, L\right) \geq \varepsilon$.

It is easy to prove from the above remark the following:
Proposition 10.7 (Opposite of convergence). A sequence $\left(x_{k}\right)_{k}$ with values in $(X, d)$ does not have $L \in X$ as its limit if and only if there exists some $\varepsilon>0$ and $n_{1}<n_{2}<n_{3}<\cdots \in \mathbb{N}$ such that $d\left(x_{n_{j}}, L\right) \geq \varepsilon$ for all $j$. In other words, we can find a subsequence $\left(x_{n_{j}}\right)_{j}$ which completely stays out of some $\varepsilon$-neighborhood of $L$.

The proof is left as exercise 10.23.

### 10.1.5 Abstract Topological spaces

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Theorem 10.2 on p. 191 gives us a way of defining neighborhoods for sets which do not have a metric.
Definition 10.12 (Abstract topological spaces). Let $X$ be an arbitrary nonempty set and let $\mathfrak{U}$ be a set of subsets ${ }^{95}$ of $X$ whose members satisfy the properties $\mathrm{a}, \mathrm{b}$ and c of (10.19) on p.191:

$$
\begin{align*}
& \text { An arbitrary union } \bigcup_{i \in I} U_{i} \text { of sets } U_{i} \in \mathfrak{U} \text { belongs to } \mathfrak{U},  \tag{10.22a}\\
& U_{1}, U_{2}, \ldots, U_{n} \in \mathfrak{U}(n \in \mathbb{N}) \Rightarrow U_{1} \cap U_{2} \cap \ldots \cap U_{n} \in \mathfrak{U},  \tag{10.22b}\\
& X \in \mathfrak{U} \text { and } \emptyset \in \mathfrak{U} \text {. } \tag{10.22c}
\end{align*}
$$

[^71]Then $(X, \mathfrak{U})$ is called a topological space The members of $\mathfrak{U}$ are called open sets. of $(X, \mathfrak{U})$ and the collection $\mathfrak{U}$ of open sets is called the topology of $X$.

Definition 10.13 (Topology induced by a metric). Let $(X, d)$ be a metric space and let $\mathfrak{U}_{d}$ be the set of open subsets of $(X, d)$, i.e., all sets $U \subseteq X$ which consist of interior points only: for each $x \in U$ there exist $\varepsilon>0$ such that

$$
N_{\varepsilon}(x)=\{y \in X: d(x, y)<\varepsilon\} \subseteq U
$$

(see (10.7) on p.190). We have seen in theorem (10.2) that those open sets satisfy the conditions of the previous definition. In other words, $\left(X, \mathfrak{U}_{d}\right)$ defines a topological space. We say that its topology is induced by the metric $d(\cdot, \cdot)$ or that it is generated by the metric $d(\cdot, \cdot)$. If there is no confusion about which metric we are talking about, we also simply speak about the metric topology.

Let $V$ be a vector space with a norm $\|\cdot\|$. Remember that any norm defines a metric $d_{\|\cdot\|}(\cdot, \cdot)$ via $d_{\|\cdot\|}(x, y)=\|x-y\|$ (see (10.1) on p.185). Obviously, this norm defines open sets

$$
\mathfrak{U}_{\|\cdot\|}:=\mathfrak{U}_{d_{\|\cdot\|}}
$$

on $V$ by means of this metric. We say that this topology is induced by the norm $\|\cdot\|$ or that it is generated by the norm $\|\cdot\|$. If there is no confusion about which norm we are talking about, we also simply speak about the norm topology of $V$.

Example 10.5 (Discrete topology). Let $X$ be nonempty. Def. 10.3 on p. 186 gave the discrete metric as the function

$$
d(x, y)= \begin{cases}0 & \text { for } x=y \\ 1 & \text { for } x \neq y\end{cases}
$$

The associated topology is

$$
\mathfrak{U}_{d}=2^{X}=\{A: A \subseteq X\} .
$$

Note that the discrete metric defines the biggest possible topology on $X$, i.e., the biggest possible collection of subsets of $X$ whose members satisfy properties a, b, c of definition 10.12 on p.193. We call this topology the discrete topology of $X$.

Example 10.6 (Indiscrete topology). Here is an example of a topology which is not generated by a metric. Let $X$ be an arbitrary nonempty set and define $\mathfrak{U}:=\{\emptyset, X\}$. Then $(X, \mathfrak{U})$ is a topological space. This is trivial because any intersection of members of $\mathfrak{U}$ is either $\emptyset$ (if at least one member is $\emptyset$ ) or $X$ (if all members are $X$ ). Conversely, any union of members of $\mathfrak{U}$ is either $\emptyset$ (if all members are $\emptyset$ ) or $X$ (if at least one member is $X$ ).

The topology $\{\emptyset, X\}$ is called the indiscrete topology of $X$. It is the smallest possible topology on $X$.

Definition 10.14 (Base of the topology). Let $(X, \mathfrak{U})$ be a topological space.
A subset $\mathfrak{B} \subseteq \mathfrak{U}$ of open sets is called a base of the topology if any nonempty open set $U$ can be written as a union of elements of $\mathfrak{B}$ :

$$
\begin{equation*}
U=\bigcup_{i \in I} B_{i} \quad\left(B_{i} \in \mathfrak{B} \text { for all } i \in I\right) \tag{10.23}
\end{equation*}
$$

where I is a suitable index set which of course will in general depend on $U$.

We note that, because $X$ itself is open, (10.23) implies that $X=\bigcup[B: B \in \mathfrak{B}]$.
Definition 10.15 (Neighborhoods and interior points in topological spaces). Let $x \in X$ and $A \subseteq X$. It is not assumed that $A$ be open. $A$ is called a neighborhood of $x$ and $x$ is called an inner point or interior point of $A$ if there exists an open set $U$ such that

$$
\begin{equation*}
x \in U \subseteq A \tag{10.24}
\end{equation*}
$$

$x$ is called an exterior point of $A$ if there exists an open set $U^{\prime}$ such that

$$
\begin{equation*}
x \in U^{\prime} \subseteq A^{\complement} \tag{10.25}
\end{equation*}
$$

i.e., $x$ is an inner point of $A^{\complement}$.

We call $\operatorname{ext}(A):=\{$ all exterior points of A$\}$ the open exterior of $A .{ }^{96}$
Definition 10.16 (Interior of a set in a topological space). Let ( $X, \mathfrak{U}$ ) be a topological space and let $A \subseteq X$. The union of all open subsets of $A$

$$
\begin{equation*}
A^{o}:=\bigcup[U \in \mathfrak{U}: U \subseteq A] \tag{10.26}
\end{equation*}
$$

is called the interior of $A$. An alternate notation for $A^{o}$ is $\operatorname{int}(A)$.
It follows from def. 10.12 (abstract topological spaces) on $p .193$ that $A^{o}$ is an open set which is, as a union of subsets of $A$, also a subset of $A$. Because $A^{\circ}$ is the union of all such sets, we conclude that

The interior $A^{o}$ of $A$ is the largest of all open subsets of $A$.

Remark 10.9. Let $(X, \mathfrak{U})$ be a topological space and $A \subseteq X$. Note that the open exterior of $A$ is

$$
\begin{equation*}
\operatorname{ext}(A)=\overline{A^{C}} . \tag{10.27}
\end{equation*}
$$

Remark 10.10. Let ( $X, \mathfrak{U}$ ) be a topological space.

$$
\begin{equation*}
\text { If } A \subseteq B \subseteq X \text { then } A^{o} \subseteq B^{o} . \tag{10.28}
\end{equation*}
$$

Definition 10.17 (boundary points in topological spaces). Let ( $X, \mathfrak{U}$ ) be a topological space and $A \subseteq X$. Then $x \in X$ is called a boundary point of $A$ if any neighborhood of $x$ intersects both $A$ and $A^{\complement}$. We write $\partial A$ for the set of all boundary points of $A$ and call this set the boundary of $A$.

[^72]Figure 10.3: Inner points v.s. boundary points.
Source: https: / / en.wikipedia.org/wiki/Interior_(topology)


We note that this definition is exactly the same as that given for metric spaces (compare def.10.10 on p .191 ).

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> The following prop. belongs into the next (optional) chapter10.1.6 (Neighborhood Bases).

Proposition 10.8. Let $(X, d)$ be a metric space and let $\mathfrak{B}:=\left\{N_{1 / k}(x): x \in X, k \in \mathbb{N}\right\}$. Then $\mathfrak{B}$ is a base of the topology for the associated topological space $\left(X, \mathfrak{U}_{d}\right)$.

The proof is left as exercise 10.12.
Definition 10.18 (Second axiom of countability). Let ( $X, \mathfrak{U}$ ) be a topological space.
We say that $X$ satisfies the second axiom of countability or $X$ is second countable if we can find a countable base for $\mathfrak{U}$.

Theorem 10.4 (Euclidean space $\mathbb{R}^{n}$ is second countable). Let

$$
\begin{equation*}
\mathfrak{B}:=\left\{N_{1 / j}(\vec{q}): \vec{q} \in \mathbb{Q}^{n}, j \in \mathbb{N}\right\} . \tag{10.29}
\end{equation*}
$$

Here $\mathbb{Q}^{n}=\left\{\vec{q}=\left(q_{1}, \ldots, q_{n}\right): q_{j} \in \mathbb{Q}, 1 \leqq j \leqq n\right\}$ is the set of all points in $\mathbb{R}^{n}$ with rational coordinates.
Then $\mathfrak{B}$ is a countable base of $\mathbb{R}^{n}$.
PROOF (outline): Let $U \in \mathfrak{U}$ be an arbitrary open set in $X$. Any vector $\vec{x} \in U$ is interior point of $U$, hence we can find some $n_{\vec{x}} \in \mathbb{N}$ such that the entire $\frac{2}{n_{\vec{x}}}$-neighborhood $N_{2 /\left(n_{\vec{x}}\right)}(\vec{x})$ is contained within $U$.

As any vector can be approximated by vectors with rational coordinates, there exists $\vec{q}=\vec{q} \vec{x} \in \mathbb{Q}^{n}$ such that $d\left(\vec{x}, \vec{q}_{\vec{x}}\right)<\frac{1}{n_{\vec{x}}}$, hence $\vec{x} \in N_{1 / n_{\vec{x}}}\left(\vec{q}_{\vec{x}}\right)$.

It follows from $N_{2 /\left(n_{\vec{x}}\right)}(\vec{x}) \subseteq U$ and prop. 10.3 on $p .190$, applied to $\delta=\varepsilon=\frac{1}{n_{\vec{x}}}$, that

$$
\begin{aligned}
& N_{1 / n_{\vec{x}}}\left(\vec{q}_{\vec{x}}\right) \subseteq N_{2 / n_{\vec{x}}}(\vec{x}) \subseteq U \text { for all } \vec{x} \in U . \\
& \text { Hence } U=\bigcup_{\vec{x} \in U}\{\vec{x}\} \subseteq \bigcup\left[N_{1 / n_{\vec{x}}}\left(\vec{q}_{\vec{x}}\right): \vec{x} \in U\right] \subseteq U .
\end{aligned}
$$

We have managed to write the arbitrarily chosen open set $U$ as a union of the sets $N_{1 / n_{\vec{x}}}\left(\vec{q}_{\vec{x}}\right)$ which belong to $\mathfrak{B}$. This proves that $\mathfrak{B}$ is a basis of the topology.

We recall from cor. 7.6 on $p .122$ that $\mathbb{Q}^{n}$ is countable. For $j \in \mathbb{N}$ let $\mathfrak{B}_{j}:=\left\{N_{1 / j}(\vec{q}): \vec{q} \in \mathbb{Q}^{n}\right\}$. Then each $\mathfrak{B}_{j}$ is countable because $\vec{q} \mapsto N_{1 / j}(\vec{q})$ is a surjection from the countable set $\mathbb{Q}^{n}$ onto $\mathfrak{B}_{j}$. It follows that the base of the topology $\mathfrak{B}=\bigcup_{j} \mathfrak{B}_{j \in \mathbb{N}}$ is countable as the countable union of countable sets.

We'll conclude this chapter with a summary of what we have learned about the classification of sets with a concept of closeness of points.

Remark 10.11 (Hierarchy of topological spaces). We have seen the following:
a. $\quad \mathbb{R}^{N}$ is an inner product space (see (9.8) on p.169).
b. All inner product spaces are normed spaces (see (9.2) on p.172).
c. All normed spaces are metric spaces (see (10.1) on p.185).
d. All metric spaces are topological spaces. (see (10.12) on p. 193 and (10.13) on p.194).

### 10.1.6 Neighborhood Bases ( $\star$ )

## Note that this chapter is starred, hence optional.

Definition 10.19 (Neighborhood base). Let $(X, \mathfrak{U})$ be a topological space.
The set of subsets of $X$

$$
\begin{equation*}
\mathfrak{N}(x):=\{A \subseteq X: A \text { is a neighborhood of } x\} \tag{10.30}
\end{equation*}
$$

is called the neighborhood system of $x$
Given a point $x \in X$, any subset $\mathfrak{B}:=\mathfrak{B}(x) \subseteq \mathfrak{N}(x)$ of the neighborhood system of $x$ is called a neighborhood base of $x$ if it satisfies the following condition: For any $A \in \mathfrak{N}(x)$ you can find a $B \in \mathfrak{B}(x)$ such that $B \subseteq A$.

In many propositions where proving closeness to some element is the issue, It often suffices to show that something is true for all sets that belong to a neighborhood base of $x$ rather than having to show it for all neighborhoods of $x$. The reason is that often only the small neighborhoods matter and a neighborhood base has "enough" of those.

Definition 10.20 (First axiom of countability). Let ( $X, \mathfrak{U}$ ) be a topological space.
We say that $X$ satisfies the first axiom of countability or $X$ is first countable if we can find for each $x \in X$ a countable neighborhood base.

Proposition 10.9 ( $\varepsilon$-neighborhoods are a base of the topology). Let ( $X, d$ ) be a metric space. Then the set $\mathscr{B}_{1}:=\left\{N_{\varepsilon}(x): x \in X, \varepsilon>0\right\}$ is a base for the topology of $(X, d)$ (see 10.14 on $p .194$ ) and the same is true for the "thinner" set $\mathscr{B}_{2}:=\left\{N_{1 / n}(x): x \in X, n \in \mathbb{N}\right\}$.

PROOF: To show that $\mathscr{B}_{1}$ (resp., $\mathscr{B}_{2}$ ) is a base we must prove that any open subset of $X$ can be written as a union of (open) sets all of which belong to $\mathscr{B}_{1}$ (resp., $\mathscr{B}_{2}$ ). We prove this for $\mathscr{B}_{2}$.

Let $U \subseteq X$ be open. As any $x \in U$ is an interior point of $U$ we can find some $\varepsilon=\varepsilon(x)>0$ such that $N_{\varepsilon(x)}(x) \subseteq U$. We note that for any such $\varepsilon(x)$ there is $n(x) \in \mathbb{N}$ such that $1 / n(x) \leqq \varepsilon(x)$.

We observe that $U \subseteq \bigcup\left[N_{1 / n(x)}(x): x \in U\right] \subseteq U$.
The first inclusion follows from $\{x\} \subseteq N_{1 / n(x)}(x)$ for all $x \in U$ and the second inclusion follows from $N_{1 / n(x)}(x) \subseteq N_{\varepsilon(x)}(x) \subseteq U$ and the inclusion lemma (lemma 5.1 on $p .96$ ).

We obtain $U=\bigcup\left[N_{1 / n(x)}(x): x \in U\right]$ and we have managed to represent our open $U$ as a union of elements of $\mathscr{B}_{2}$. This proves that $\mathscr{B}_{2}$ is a base for the topology of $(X, d)$.

As $\mathscr{B}_{2} \subseteq \mathscr{B}_{1}$ it follows that $\mathscr{B}_{1}$ also is such a base.
Theorem 10.5 (Metric spaces are first countable). Let $(X, d)$ be a metric space. Then $X$ is first countable.
Proof (outline): For any $x \in X$ let

$$
\begin{equation*}
\mathfrak{B}(x):=\left\{N_{1 / n}(x): n \in \mathbb{N}\right\} . \tag{10.31}
\end{equation*}
$$

Then $\mathfrak{B}(x)$ is a neighborhood base of $x$ because, by def. 10.7 and def. 10.9 (Interior points and neighborhoods in metric spaces), any neighborhood of $x$ will contain one of the form $N_{\varepsilon}(x)$ and for any such $\varepsilon>0$ there exists $n \in \mathbb{N}$ such that $\frac{1}{n}<\varepsilon$.

### 10.1.7 Metric and Topological Subspaces

Definition 10.21 (Metric subspaces). Given is a metric space ( $X, d$ ) and a nonempty $A \subseteq(X, d)$. Let $\left.d\right|_{A \times A}: A \times A \rightarrow \mathbb{R}_{\geq 0}$ be the restriction $\left.d\right|_{A \times A}(x, y):=d(x, y)(x, y \in A)$ of the metric $d$ to $A \times A$ (see def.4.11 on p.83). It is trivial to verify that $\left(A,\left.d\right|_{A \times A}\right)$ is a metric space in the sense of def.10.1 on p.183. We call $\left(A,\left.d\right|_{A \times A}\right)$ a metric subspace of $(X, d)$ and we call $\left.d\right|_{A \times A}$ the metric induced by $d$ or the metric inherited from $(X, d)$.

## Remark 10.12.

Metric subspaces come with their own collections of open and closed sets,
 neighborhoods, $\varepsilon$-neighborhoods, convergent sequences, ...
You must watch out when looking at statements and their proofs whether those concepts refer to the entire space $(X, d)$ or to the subspace $\left(A,\left.d\right|_{A x A}\right)$.

## Notations 10.1.

a) Because the only difference between $d$ and $d_{A \times A}$ is the domain, it is customary to write $d$ instead of $\left.d\right|_{A \times A}$ to make formulas look simpler if doing so does not give rise to confusion.
b) We often shorten "open in $\left(A,\left.d\right|_{A \times A}\right)$ " to "open in $A$ ", "closed in $\left(A,\left.d\right|_{A \times A}\right)$ " to "closed in $A$ ", "convergent in $\left(A,\left.d\right|_{A \times A}\right)$ " to "convergent in $A^{\prime}, \ldots .$.

Definition 10.22 (Traces of sets in a metric subspace). Let $(X, d)$ be a metric space and $A \subseteq X$ a nonempty subset of $X$, viewed as a metric subspace $\left(A,\left.d\right|_{A \times A}\right)$ of $(X, d)$. Let $Q \subseteq X$. We call $Q \cap A$ the trace of $Q$ in $A$.

For $\varepsilon>0$ and $a \in A$ let $N_{\varepsilon}(a)$ be the $\varepsilon$-neighborhood of $a$ (in $(X, d)$ ). We define

$$
N_{\varepsilon}^{A}(a):=N_{\varepsilon}(a) \cap A
$$

i.e., $N_{\varepsilon}^{A}(a)$ is defined as the trace of $N_{\varepsilon}(a)$ in $A$.

Proposition 10.10 (Open sets in metric subspaces as traces of open sets in $X$ ). Let $(X, d)$ be a metric space and $A \subseteq X$ a nonempty subset of $X$.
a. Let $\varepsilon>0$ and $a \in A$. Then

$$
\begin{equation*}
N_{\varepsilon}^{A}(a)=\left\{x \in A:\left.d\right|_{A \times A}(x, a)<\varepsilon\right\} \tag{10.32}
\end{equation*}
$$

i.e., $N_{\varepsilon}^{A}(a)$ is the $\varepsilon-$ Neighborhood of a as defined in def. 10.6 on p.189.

It then follows from

$$
\begin{equation*}
N_{\varepsilon}^{A}(a)=N_{\varepsilon}(a) \cap A \tag{10.33}
\end{equation*}
$$

that each $\varepsilon$-neighborhood in the subspace $A$ is the trace of an $\varepsilon$-neighborhood in $X$.
b. More generally, a set $U \subseteq A$ is open in $A$ if and only if there is an open $V \subseteq$ in $(X, d)$ such that

$$
\begin{equation*}
U=V \cap A \tag{10.34}
\end{equation*}
$$

i.e., $U$ is the trace of a set $V$ which is open in $X$.

PROOF of $a$ : First we prove (10.32). As $\left.d\right|_{A \times A}$ is the restriction of $d$ to $A \times A$ it follows that

$$
\begin{aligned}
N_{\varepsilon}^{A}(a) & =N_{\varepsilon}(a) \cap A=\{x \in X: d(x, a)<\varepsilon\} \cap A \\
& =\{x \in A: d(x, a)<\varepsilon\}=\left\{x \in A:\left.d\right|_{A \times A}(x, a)<\varepsilon\right\} .
\end{aligned}
$$

This finishes the proof of $\boldsymbol{a}$.
PROOF of $\boldsymbol{b}$ : First we show that if $V$ is open in $X$ then $U:=V \cap A$ is open in the subspace $A$.
Let $a \in U$. We must prove that $a$ is an interior point of $U$ with respect to $\left(A,\left.d\right|_{A \times A}\right)$.
Because $a \in V$ and $V$ is open in $X$, there is $\varepsilon>0$ such that $N_{\varepsilon}(a) \subseteq V$. It follows that
$N_{\varepsilon}^{A}(a)=N_{\varepsilon}(a) \cap A \subseteq V \cap A=U$ and $N_{\varepsilon}^{A}(a)$ is open in $A$, hence $a$ is interior point of $U$ with respect to the subspace $\left(A,\left.d\right|_{A \times A}\right)$.

Finally we prove that if $U \subseteq A$ is open in $A$ there is $V \subseteq X$ open in $X$ such that $U=V \cap A$ :
We can write $U=\bigcup\left[N_{\varepsilon(a)}^{A}(a): a \in U\right]$ for suitable $\varepsilon(a)>0$ (see the proof of prop.10.9 on p.197).
Let $V:=\bigcup\left[N_{\varepsilon(a)}(a): a \in U\right]$. Then $V$ is open in $\left.(X, d)\right)$ as union of the open sets $N_{\varepsilon(a)}(a)$, hence

$$
\begin{aligned}
V \cap A & =A \cap \bigcup\left[N_{\varepsilon(x)}(x): x \in U\right]=\bigcup\left[N_{\varepsilon(x)}(x) \cap A: x \in U\right] \\
& =\bigcup\left[N_{\varepsilon(x)}^{A}(x): x \in U\right]=U
\end{aligned}
$$

(the second equalitity follows from prop.5.2 on p.98). This finishes the proof.
The last proposition justifies to define subspaces of abstract topological spaces as follows.

Definition 10.23 (Topological subspaces). Let $(X, \mathfrak{U})$ be a topological space and $A \subseteq X$. We say that $V \subseteq A$ is open in $\mathbf{A}$ if $V$ is the trace of an open set in $X$, i.e., if there is some $U \in \mathfrak{U}$ such that $V=U \cap A$. We denote the collection of all open sets in $A$ as $\mathfrak{U}_{A}$, i.e.,

$$
\mathfrak{U}_{A}=\{V \cap A: V \in \mathfrak{U}\} .
$$

We call $\left(A, \mathfrak{U}_{A}\right)$ a topological subspace or also just a subspace of $(X, \mathfrak{U})$ and
we call $\mathfrak{U}_{A}$ the topology induced by the metric $d$ or the topology inherited from ( $X, d$ ).
Proposition 10.11 (Topological subspaces are topological spaces). Let ( $X, \mathfrak{L}$ ) be a topological space, $A \subseteq X$, and let $\mathfrak{U}_{A}$ be the collection of all open sets in $A$. Then $\left(A, \mathfrak{U}_{A}\right)$ is a topological space, i.e., it satisfies the definition def. 10.12 on $p .193$ of an abstract topological space.

PROOF: Left as an exercise.
Remark 10.13 (Convergence does not extend to subspaces). Let $A \subseteq(X, d)$ and $a_{n} \in A$.
Let $a_{n} \in A$ for all $n \in \mathbb{N}$. Note that convergence of the sequence $\left(a_{n}\right)$ in the space $(X, d)$ (i.e., there exists $x \in X$ such that $x=\lim _{n \rightarrow \infty} a_{n}$ ), does NOT imply convergence of the sequence in the subspace $\left(A,\left.d\right|_{A \times A}\right)$. Rather, we have the following dichotomy:
a. $\quad x \in A$ : Then $a_{n}$ converges in the subspace $\left(A,\left.d\right|_{A \times A}\right)$ (and also in $(X, d)$ ).
b. $\quad x \in A^{\complement}$ : Then $a_{n}$ converges in $(X, d)$ but not in $\left(A,\left.d\right|_{A \times A}\right)$.

### 10.1.8 Bounded Sets and Bounded Functions

Definition 10.24 (bounded sets). Given is a subset $A$ of a metric space ( $X, d$ ). The diameter of $A$ is defined as

$$
\begin{equation*}
\operatorname{diam}(\emptyset):=0, \quad \operatorname{diam}(A):=\sup \{d(x, y): x, y \in A\} \text { if } A \neq \emptyset \tag{10.35}
\end{equation*}
$$

We call $A$ a bounded set if $\operatorname{diam}(A)<\infty$.
Proposition 10.12. Given is a metric space $(X, d)$ and a nonempty subset $A$. The following are equivalent:

$$
\begin{equation*}
\text { a. } \operatorname{diam}(A)<\infty, \quad \text { i.e., } A \text { is bounded. } \tag{10.36}
\end{equation*}
$$

b. There exists $\gamma>0$ and $x_{0} \in X$ such that $A \subseteq N_{\gamma}\left(x_{0}\right)$.
c. For all $x \in X$ there exists $\gamma>0$ such that $A \subseteq N_{\gamma}(x)$.

PROOF of " $\boldsymbol{b} \Rightarrow \boldsymbol{a}^{\prime \prime}$ : For any $x, y \in A$ we have

$$
d(x, y) \leqq d\left(x, x_{0}\right)+d\left(x_{0}, y\right) \leqq 2 \gamma
$$

and it follows that $\operatorname{diam}(A) \leqq 2 \gamma$.
PROOF of " $\boldsymbol{a} \Rightarrow \boldsymbol{b}^{\prime}$ ": Pick an arbitrary $x_{0} \in A$ and let $\gamma:=\operatorname{diam}(A)$. Then for all $a \in A$

$$
d\left(x_{0}, a\right) \leqq \sup _{x \in A} d(x, a) \leqq \sup _{x, z \in A} d(x, z)=\operatorname{diam}(A)=\gamma
$$

It follows that $A \subseteq N_{\gamma}\left(x_{0}\right)$.

PROOF of " $\boldsymbol{c} \Rightarrow \boldsymbol{a}$ ": We pick an arbitrary $x_{0} \in A$ which is possible as $A$ is not empty. Then there is $\gamma=\gamma\left(x_{0}\right)$ such that $A \subseteq N_{\gamma}\left(x_{0}\right)$. For any $y, z \in A$ we then have

$$
d(y, z) \leqq d\left(y, x_{0}\right)+d\left(x_{0}, z\right) \leqq 2 \gamma
$$

and it follows that $\operatorname{diam}(A) \leqq 2 \gamma<\infty$.
PROOF of " $\boldsymbol{a} \Rightarrow \boldsymbol{c}$ ": Given $x \in X$, pick an arbitrary $x_{0} \in A$ and let $\gamma:=d\left(x, x_{0}\right)+\operatorname{diam}(A)$. Then

$$
\begin{aligned}
y \in A \quad \rightarrow \quad d(x, y) & \leqq d\left(x, x_{0}\right)+d\left(x_{0}, y\right)
\end{aligned} \begin{aligned}
& \leqq\left(x, x_{0}\right)+\sup _{u \in A} d(u, y) \\
& \leqq d\left(x, x_{0}\right)+\sup _{u, z \in A} d(u, z)=d\left(x, x_{0}\right)+\operatorname{diam}(A)=\gamma
\end{aligned}
$$

It follows that $A \subseteq N_{\gamma}(x)$.

### 10.1.9 Contact Points and Closed Sets

If you look at any closed interval $[a, b]=\{y \in \mathbb{R}: a \leqq y \leqq b\}$ of real numbers, then all of its points are interior points, except for the end points $a$ and $b$. Moreover $a$ and $b$ are contact points according to the following definition which makes sense for any abstract topological space.
Definition 10.25 (Contact points). Given is a topological space ( $X, \mathfrak{U}$ ).
Let $A \subseteq X$ and $x \in X$ ( $x$ may or may not to belong to $A$ ). $x$ is called a contact point ${ }^{97}$ of $A$ if

$$
\begin{equation*}
A \cap N \neq \emptyset \text { for any neighborhood } N \text { of } x \tag{10.39}
\end{equation*}
$$

The following theorem shows that we can characterize contact points of subsets of metric spaces by means of sequences.

Theorem 10.6 (Sequence criterion for contact points in metric spaces). Given is a metric space ( $X, d$ ). Let $A \subseteq X$ and $x \in X$. Then $x$ is a contact point of $A$ if and only if there exists a sequence $x_{1}, x_{2}, x_{3}, \ldots$ of members of $A$ which converges to $x$.

PROOF of " $\Rightarrow$ ": Let $x \in X$ be such that $N \cap A \neq \emptyset$ for any neighborhood $N$ of $x$. Let $x_{n} \in N_{1 / n}(x) \cap A$. Such $x_{n}$ exists because the neighborhood $N_{1 / n}(x)$ has nonempty intersection with $A$.

Given $\varepsilon>0$, let $N \in \mathbb{N}$ be chosen such that $\frac{1}{N}<\varepsilon$. This is possible because $\mathbb{N}$ is not bounded (above) in $\mathbb{R}$.

For any $j \geqq N$ we obtain $d\left(x_{j}, x\right)<1 / j \leqq 1 / N<\varepsilon$. This proves convergence $x_{n} \rightarrow x$.
PROOF of " $\Leftarrow$ " Let $x \in X$ and assume there is $\left(x_{n}\right)_{n \in \mathbb{N}}$ such that $x_{n} \in A$ for all $n$ and $x_{n} \rightarrow x$.
We must show that if $U_{x}$ is a (open) neighborhood of $x$ then $U_{x} \cap A \neq \emptyset$. Let $\varepsilon>0$ such that $N_{\varepsilon}(x) \subseteq U_{x}$.
It follows from $x_{n} \rightarrow x$ that there is $N=N(\varepsilon)$ such that $x_{n} \in N_{\varepsilon}(x)$ for all $n \geqq N$, especially, $x_{N} \in N_{\varepsilon}(x)$. By assumption, $x_{N} \in A$, hence $x_{N} \in N_{\varepsilon}(x) \cap A \subseteq U_{x} \cap A$, hence $U_{x} \cap A \neq \emptyset$.
Note 10.1. Note that any $a \in A$ is a contact point of $A$ but not necessarily the other way around:
a. Let $a \in A$. Then any neighborhood $U_{a}$ of $a$ contains $a$, hence $U_{A} \cap A$ is not empty, hence $a$ is a contact point of $A$. This proves that any $a \in A$ is a contact point of $A$.
b. Here is a counterexample which shows that the converse need not be true. Let $(X, d):=\mathbb{R}$ with the standard Euclidean metric and let $A$ be the subset $] 0,1[$. We show now that 0 is a contact point of $A$.

[^73]Any neighborhood $A_{0}$ of 0 contains for some small enough $\delta>0$ the entire interval $]-\delta, \delta[$ Let $x:=\min (\delta / 2,1 / 2)$.
Clearly, $x \in]-\delta, \delta\left[\subseteq A_{0}\right.$ and $\left.x \in\right] 0,1[=A$.
It follows that $x \in A \cap A_{0}$. As $A_{0}$ was an arbitrary neighborhood of 0 , we have proved that 0 is a contact point of $A$, even though $0 \notin A$.
c. The above counterexample can be proven much faster if the criterion for contact points in metric spaces is employed: Let $x_{n}:=1 / n(n \geqq 2)$ Then $\left.x_{n} \in\right] 0,1[$ for all $n$ and the sequence converges to 0 . It follows that 0 is a contact point of $] 0,1[$.

Note $\mathbf{1 0 . 2}$ (Contact points vs Limit points). Besides contact points there also is the concept of a limit point. Here is the definition (see [9] Munkres, a standard book on topology):

Given is a metric space $(X, d)$. Let $A \subseteq X$ and $x_{0} \in X . x_{0}$ is called a limit point or cluster point or point of accumulation of $A$ if any neighborhood $U$ of $x_{0}$ intersects $A$ in at least one point other than $x_{0}$. This definition excludes "isolated points" 98 of $A$ from being limit points of $A$.

Definition 10.26 (Closed sets). Given is a metric space ( $X, d$ ) and a subset $A \subseteq X$. We call

$$
\begin{equation*}
\bar{A}:=\{x \in X: x \text { is a contact point of } A\} \tag{10.40}
\end{equation*}
$$

the closure of $A$. A set that contains all its contact points is called a closed set.
Remark 10.14. It follows from note 10.1.a that $A \subseteq \bar{A}$.
Proposition 10.13. The complement of an open set is closed.
PROOF of 10.13: Let $A$ be an open set. Each point $a \in A$ is an interior point which can be surrounded by a $\delta$-neighborhood $N_{\delta}(a)$ which, for small enough $\delta$, will be entirely contained within $A$.

Assume $x \in X$ is a contact point of $A^{\complement}$. We want to prove that $A^{\complement}$ is a closed set, so we must show that $x \in A^{\complement}$.

We assume to the contrary that $x \in A$. $A$ is open, so $x$ is an interior point of $A$. It follows that there is a neighborhood $N_{\delta}(x)$ surrounding $x$ which is entirely contained in $A$, hence $N_{\delta}(x) \cap A^{\complement}=\emptyset$.

We also assumed that $x$ is a contact point of $A^{\complement}$. This implies that $N_{\delta}(x) \cap A^{\complement} \neq \emptyset$. We have reached a contradiction.

Proposition 10.14. The complement of a closed set is open.
PROOF: We will give two proofs of the above. The first one makes use of the criterion for contact points (theorem 10.6) and works with sequences. The second proof is based on the definition of contact points and works with neighborhoods and interior points.
a. First proof of prop.10.14:

Let $A \subset(X, d)$ be closed. If $A^{\complement}$ is not open then there must some be $b \in A^{\complement}$ which is not an interior point of $A^{\complement}$.

We show that this assumption leads to a contradiction.
Because $b$ is not an interior point of $A^{\complement}$, there is no $\delta$-neighborhood, for whatever small $\delta$, that entirely belongs to $A^{\complement}$. So, for each $j \in \mathbb{N}$, there is an $x_{j} \in N_{1 / j}(b)$ which does not belong to $A^{\complement}$, i.e., $x_{j} \in A$.

[^74]We have constructed a sequence $x_{j}$ which is entirely contained in $A$ and which converges to $b$. The latter is true because, for any $j$, all but finitely many members are contained in $N_{1 / j}(b)$.

The closed set A contains all its contact points and it follows from the criterion for contact points that $b \in A$.

But we had assumed at the outset that $b \in A^{\complement}$ and we have a contradiction.
b. Alternate proof of prop.10.14, based on neighborhoods and interior points:

Let $A$ be closed set and $b \in A^{\complement}$.
The closed set $A$ contains all its contact points, so $b \notin A$ implies that $b$ is not a contact point of $A$ : according to def. 10.25 there exists some neighborhood $V$ of $b$ such that $V \cap A=\emptyset$, i.e., $V \subseteq A^{\complement}$.

We have proved that an arbitrary $b \in A^{\complement}$ is an interior point of $A^{\complement}$, i.e., the complement of the closed set $A$ is open.
Theorem 10.7 (Open iff complement is closed). Let $(X, d)$ be a metric space and $A \subseteq X$. Then $A$ is open if and only if $A^{\complement}$ is closed.

PROOF: Immediate from prop.10.13 and prop.10.14
Remark 10.15. Note the following.
a. We have seen that def.10.25 for contact points and hence def.10.26 for closed sets are entirely based on the concept of neighborhood which itself is entirely based on that of open sets. It follows that those two definitions make perfect sense not only in metric spaces but, more generally, in abstract topological spaces ( $X, \mathfrak{U}$ ) which are characterized by the set $\mathfrak{U}$ of all open subsets of $X$ (see def.10.12 on p.193).
b. Moreover the proof for prop. 10.13 (complements of open sets are closed) and the first proof for prop.10.14 (complements of closed sets are open) are based on those definitions and do not employ specific properties of metric spaces either; theorem 10.7 also works for abstract topological spaces.
c. Matter of fact, many books define closed sets as the complements of open sets and only afterwards define contact points as we did. No surprise then that our definition of closed sets becomes their theorem: It is proved from those definitions that closed sets are exactly those that contain all their contact points.
Definition $\mathbf{1 0 . 2 7}$ (Contact points, closed sets and boundary points in topological spaces). ${ }^{99}$ Given is an abstract topological space $(X, \mathfrak{U})$.

Let $A \subseteq X$ and $x \in X$ ( $x$ may or may not to belong to $A$ ). $x$ is called a contact point of $A$ if

$$
A \cap N \neq \emptyset \text { for any neighborhood } N \text { of } x .
$$

We call

$$
\bar{A}:=\{x \in X: x \text { is a contact point of } A\}
$$

the closure of $A$. A set that contains all its contact points is called a closed set.
A point $x \in X$ is called a boundary point of $A$ if any neighborhood of $x$ intersects both $A$ and $A^{\complement}$. As for metric spaces, we write $\partial A$ for the set of all boundary points of $A$ and call this set the boundary of $A$.

[^75]Remark 10.16. We note that $A \subseteq \bar{A}$ : Let $a \in A$ and let $V_{a}$ be a neighborhood of $a$. Because $a \in V_{a}$, we obtain $a \in V_{a} \cap A$, hence $V_{a} \cap A \neq \emptyset$, hence $a \in \bar{A}$.

It follows that $A$ is closed if and only if $A=\bar{A}$ (which justifies the name "closure of $A$ " for $\bar{A}$.)

Proposition 10.15. Let $(X, \mathfrak{U})$ be a topological space and $A \subseteq B \subseteq X$. Then $\bar{A} \subseteq \bar{B}$.
The proof is left as exercise 10.16.
Proposition 10.16. Let $(X, \mathfrak{U})$ be a topological space.
The closed sets of $X$ satisfy the following property:
a. An arbitrary intersection of closed sets is closed.
b. A finite union of closed sets is closed.
c. The entire set $X$ is closed and $\emptyset$ is closed.

The proofs of $\boldsymbol{a}$ and $\boldsymbol{b}$ follow easily from De Morgan's law (the duality principle for sets: see (5.1) on p.97). Observe that $X$ plays the role of a universal set because all members $U$ of $\mathfrak{U}$ and their complements $U^{\complement}$ are subsets of $X$.
PROOF of $\boldsymbol{a}$ : Let $\left(C_{\alpha}\right)$ be an arbitrary familiy of closed sets. Then $U_{\alpha}:=C_{\alpha}^{\complement}$ is an open set for each $\alpha$. Observe that $C_{\alpha}^{\complement}=U_{\alpha}$ because the complement of the complement of any set gives you back that set. Let $C:=\bigcap_{\alpha} C_{\alpha}$. Then

$$
C^{\complement}=\left(\bigcap_{\alpha} C_{\alpha}\right)^{\complement}=\bigcup_{\alpha} C_{\alpha}^{\complement}=\bigcup_{\alpha} U_{\alpha} .
$$

In other words $C^{\complement}$ is an arbitrary union of open sets which is open by the very definition of open sets of a topological space. We have proved a.
PROOF of $\boldsymbol{b}$ : Let $C_{1}, C_{2}, \ldots C_{n}$ be closed sets. Then $U_{j}:=C_{j}^{\complement}$ is an open set for each $j$. Let $C:=\bigcup_{1 \leqq j \leqq n} C_{j}$.
Then

$$
C^{\complement}=\left(\bigcup_{j} C_{j}\right)^{\complement}=\bigcap_{j} C_{j}^{\complement} .=\bigcap_{j} U_{j}
$$

Hence, $C^{\complement}$ is the intersection of finitely many open sets. This shows that $C^{\complement}$ is open, i.e., $C$ is closed. We have proved $b$.
PROOF of $\boldsymbol{c}$ : Trivial because

$$
X^{\complement}=\emptyset \quad \text { and } \quad \emptyset^{\complement}=X .
$$

Proposition 10.17. Let $(X, \mathfrak{U})$ be a topological space and $A \subseteq X$. Then

$$
\begin{equation*}
\partial A=\bar{A} \cap \overline{A^{\complement}} \tag{10.42}
\end{equation*}
$$

i.e., $x \in X$ is a boundary point of $A$ if and only if $x$ is a contact point of both $A$ and $A^{\complement}$.

PROOF: Left as an exercise.

Proposition 10.18 (Minimality of the closure of a set). Let $(X, \mathfrak{U})$ be a topological space and $A \subseteq X$. Then

$$
\begin{equation*}
\bar{A}=\bigcap[C \supseteq A: C \text { is closed }] . \tag{10.43}
\end{equation*}
$$

The closure $\bar{A}$ of $A$ is the smallest of all closed supersets of $A$.

PROOF: Let $\mathfrak{C}:=\{C \supset A: C$ is closed $\}$ and let $F:=\bigcap \mathfrak{C}$. We need to show that $\bar{A}=F$.
It follows from prop.10.16.a that $F$ is closed, hence $F=\bar{F}$. It follows from $C \supseteq A$ for all $C \in \mathfrak{C}$ that $F \supseteq A$, hence $F=\bar{F} \supseteq \bar{A}$.

It remains to be shown that $F \subseteq \bar{A}$. It is true that $\bar{A} \in \mathfrak{C}$ because $\bar{A}$ is a closed set which contains $A$, hence $\bar{A} \supseteq \bigcap \mathfrak{C}=F$. (See prop. 10.15 on $p .204$ ).

Proposition 10.19 (Closure of a set as a hull operator ${ }^{100}$ ). Let ( $X, \mathfrak{U}$ ) be a topological space. We can think of the closure of sets as a function ${ }^{-}: 2^{X} \rightarrow 2^{X} ; \quad A \mapsto \bar{A}$. This function has the following properties:
a. $\bar{\emptyset}=\emptyset$,
b. $A \subseteq \bar{A}$,
c. $\bar{A}=\bar{A}$,
d. $\overline{A \cup B}=\bar{A} \cup \bar{B}$.

PROOF: $\boldsymbol{a}$ follows from (10.41).c and $\boldsymbol{b}$ follows from remark 10.14.
The proof of $\boldsymbol{c}$ and $\boldsymbol{d}$ is left as exercise 10.18.

### 10.1.10 Completeness in Metric Spaces

Often you are faced with a situation where you need to find a contact point a and all you have is a sequence which behaves like one converging to a contact point in the sense of inequality 10.20 (page 192)

Definition 10.28 (Cauchy sequences). Given is a metric space ( $X, d$ ).
A sequence $\left(x_{n}\right)$ in $X$ is called a Cauchy sequence ${ }^{101}$ or, in short, it is Cauchy if it has the following property: Given any whatever small number $\varepsilon>0$, you can find a (possibly very large) number $n_{0}$ such that

$$
\begin{equation*}
d\left(x_{i}, x_{j}\right)<\varepsilon \quad \text { for all } i, j \geqq n_{0} \tag{10.44}
\end{equation*}
$$

This is called the Cauchy criterion for convergence of a sequence.
${ }^{100}$ This proposition states that the closure is a so-called closure operator which is defined to be a function

$$
c l: 2^{X} \rightarrow 2^{X} ; \quad A \mapsto c l(A):=\bar{A}
$$

on some abstract, nonempty set $X$ (which need not be a topological space) such that the following are satisfied:
a. $(\emptyset)=\emptyset$,
b. $A \subseteq c l(A)$,
c. $c l(c l(A))=c l(A)$,
d. $c l(A \cup B)=c l(A) \cup c l(B)$.

It can be shown that if we define

$$
\mathfrak{U}:=\left\{A^{\complement}: c l(A)=A\right\}
$$

then $(X, \mathfrak{U})$ satisfies the properties of a topological space.
${ }^{101}$ so named after the great french mathematician Augustin-Louis Cauchy (1789-1857) who contributed massively to the most fundamental ideas of Calculus.

Example 10.7 (Cauchy criterion for real numbers). In $\mathbb{R}$ we have $d(x, y)=|x-y|$ and the Cauchy criterion requires for any given $\varepsilon>0$ the existence of $n_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\left|x_{i}-x_{j}\right|<\varepsilon \quad \text { for all } i, j \geqq n_{0} . \tag{10.45}
\end{equation*}
$$

Proposition 10.20. A Cauchy sequence in a metric space is bounded.
PROOF: Let $\left(x_{n}\right)_{n}$ be a Cauchy sequence in a metric space $(X, d)$. There is $N=N(1 / 2)$ such that $d\left(x_{i}, x_{j}\right)<1 / 2$ for all $i, j \geqq N$. In particular, $d\left(x_{i}, x_{N}\right)<1 / 2$.

Let $M:=\max \left\{d\left(x_{j}, x_{N}\right): j<N\right\}$. We obtain for any two indices $i, j \in \mathbb{N}$ that

$$
d\left(x_{i}, x_{j}\right) \leqq d\left(x_{i}, x_{N}\right)+d\left(x_{N}, x_{j}\right) .
$$

$d\left(x_{i}, x_{N}\right)$ is bounded by $M$ in case that $i<N$ and by $1 / 2$ if $i \geqq N$; hence $d\left(x_{i}, x_{N}\right)<1 / 2+M$. We use the same reasoning to conclude that $d\left(x_{N}, x_{j}\right)<1 / 2+M$ and obtain $d\left(x_{i}, x_{j}\right)<1+2 M$. This proves the boundedness of $\left(x_{n}\right)_{n}$.

The following theorem of the completeness of the set of all real numbers ${ }^{102}$ states that any Cauchy sequence converges to a real number. This is a big deal: To show that a sequence has a finite limit you need not provide the actual value of that limit. All you must show is that this sequence satisfies the Cauchy criterion. One can say that this preoccupation with proving existence rather than computing the actual value is one of the major points which distinguish Mathematics from applied Physics and the engineering disciplines.

Here is the formal definition of a complete set in a metric space.
Definition 10.29 (Completeness in metric spaces). Given is a metric space ( $X, d$ ). A subset $A \subseteq X$ is called complete if any Cauchy sequence $\left(x_{n}\right)$ with elements in $A$ converges to an element of $A$.

## Remark 10.17.

a. In particular, $X$ itself is complete iff any Cauchy sequence in $X$ converges.
b. $A$ is complete as a subset of $(X, d)$ iff the subspace $\left(\left(A,\left.d\right|_{A \times A}\right)\right.$ is complete "in itself".

Theorem 10.8 (Completeness of the real numbers). The following is true for the real numbers with the metric $d(a, b)=|b-a|$ but will in general be false for arbitrary metric spaces: Let $\left(x_{n}\right)$ be a Cauchy sequence in $\mathbb{R}$. then there exists a real number $L$ such that $L=\lim _{n \rightarrow \infty} x_{n}$.

PROOF: It follows from prop.10.20 that $x_{n}$ is bounded, hence $\left(x_{n}\right)_{n}$ possesses finite liminf and limsup. ${ }^{103}$ We now show that $\liminf _{n \rightarrow \infty} x_{n}=\underset{n \rightarrow \infty}{\limsup } x_{n}$.

Let $\varepsilon>0$ and $N \in \mathbb{N}$ such that $\left|x_{i}-x_{j}\right| \leqq \varepsilon$ for all $i, j \geqq N$.
Let $T_{n}:=\left\{x_{j}: j \geqq n\right\}$ be the tail set of the sequence $\left(x_{n}\right)_{n}$. Let $\alpha_{n}:=\inf T_{n}, \beta_{n}:=\sup T_{n}$.
There is some $i \geqq N$ such that $\left|x_{i}-\alpha_{N}\right|=x_{i}-\alpha_{N} \leqq \varepsilon$ and there is some $j \geqq N$ such that $\left|\beta_{N}-x_{j}\right|=$ $\beta_{N}-x_{j} \leqq \varepsilon$. It follows that

$$
0 \leqq \beta_{N}-\alpha_{N}=\left|\beta_{N}-\alpha_{N}\right| \leqq\left|\beta_{N}-x_{j}\right|+\left|x_{j}-x_{i}\right|+\left|x_{i}-\alpha_{N}\right| \leqq 3 \varepsilon .
$$

[^76]Further, if $k \geqq N$ then $T_{k} \subseteq T_{N}$, hence $\alpha_{k} \geqq \alpha_{N}$ and $\beta_{k} \leqq \beta_{N}$. It follows that

$$
\inf \beta_{k}-\sup \alpha_{k} \leqq \beta_{N}-\alpha_{N} \leqq 3 \varepsilon
$$

But then

$$
0 \leqq \limsup _{k \rightarrow \infty} x_{k}-\liminf _{k \rightarrow \infty} x_{k}=\inf _{k} \beta_{k}-\sup _{k} \alpha_{k} \leqq 3 \varepsilon .
$$

$\varepsilon>0$ was arbitrary, hence $\limsup _{k \rightarrow \infty} x_{k}=\liminf _{k \rightarrow \infty} x_{k}$.
Figure 10.4: liminf = limsup for Cauchy sequences


Part 3: It follows from theorem 8.3 on $p .147$ that the sequence $\left(x_{n}\right)_{n}$ converges to $L:=\limsup _{k \rightarrow \infty} x_{k}$ and the proof is finished.

Now that the completeness of $\mathbb{R}$ has been established, it is not very difficult to see that $N$-dimensional space $\mathbb{R}^{N}$ also is complete.

Theorem 10.9 (Completeness of $\mathbb{R}^{N}$ ). Let $\left(\vec{x}_{n}\right)$ be a Cauchy sequence in $\mathbb{R}^{N}$.
Then there exists a vector $\vec{a} \in \mathbb{R}^{N}$ such that $\vec{a}=\lim _{n \rightarrow \infty} \vec{x}_{n}$.
Proof (outline): Let $\vec{x}_{j}=\left(x_{j, 1}, x_{j, 2}, \ldots, x_{j, N}\right)$ be Cauchy in $\mathbb{R}^{N}$. For fixed $k$, each coordinate sequence $\left(x_{j, k}\right)_{j}$ is Cauchy because, if $\varepsilon>0$, there exists $K \in \mathbb{N}$ such that if $i, j \geqq K$ then $\left\|\vec{x}_{i}-\vec{x}_{j}\right\|_{2}<\varepsilon$. Hence

$$
\left|x_{i, k}-x_{j, k}\right|=\sqrt{\left|x_{i, k}-x_{j, k}\right|^{2}} \leqq \sqrt{\sum_{k=1}^{N}\left|x_{i, k}-x_{j, k}\right|^{2}}=\left\|\vec{x}_{i}-\vec{x}_{j}\right\|_{2}<\varepsilon
$$

It follows from the completeness of $\mathbb{R}$ as a metric space that there exist real numbers

$$
a_{1}, a_{2}, a_{3}, \ldots, a_{N} \quad \text { such that } a_{k}=\lim _{n \rightarrow \infty} x_{n, k}(1 \leqq k \leqq N) .
$$

For a given number $\varepsilon$ we can find natural numbers $n_{0,1}, n_{0,2}, \ldots, n_{0, N}$ such that

$$
\left|x_{n, k}-a_{k}\right|<\frac{\varepsilon}{N} \quad \text { for all } n \geqq n_{0, j} \text { and for all } 1 \leqq k \leqq N
$$

Let $n^{\star}:=\max \left(n_{0,1}, n_{0,2}, \ldots, n_{0, N}\right)$ and $\vec{a}:=\left(a_{1}, a_{2}, \ldots, a_{n}\right)$. It follows that

$$
d\left(\vec{x}_{n}-\vec{a}\right)=\sqrt{\sum_{k=1}^{N}\left|x_{n, k}-a_{k}\right|^{2}} \leqq \sqrt{N \cdot\left(\frac{\varepsilon}{N}\right)^{2}}=\frac{\varepsilon}{\sqrt{N}} \leqq \varepsilon \quad \text { for all } n \geqq n^{\star}
$$

Example 10.8 (Approximation of decimals). The following illustrates Cauchy sequences and completeness in $\mathbb{R}$. Take any real number $x \geqq 0$ and write it as a decimal:

$$
x=m+\sum_{j=1}^{\infty} d_{j} \cdot 10^{-j} \quad\left(m \in \mathbb{Z}_{>0}, d_{j} \in\{0,1,2, \ldots, 9\}\right)
$$

As was explained in (2.9) on (p.18), anything that can be written as a decimal number is a real number. Let's say, $x$ starts out on the left as

$$
x=258.1408926584207531 \ldots
$$

We define as $x_{k}$ the leftmost part of $x$, truncated $k$ digits after the decimal points:

$$
x_{1}=258.1, \quad x_{2}=258.14, \quad x_{3}=258.140, \quad x_{4}=258.1408, \quad x_{5}=258.14089, \quad \ldots
$$

We further define $y_{k}$ the leftmost part of $x$, truncated $k$ digits after the decimal points, but the rightmost digit incremented by 1 (there might be a carry-over to the left when you add 1 to 9 )

$$
y_{1}=258.2, \quad y_{2}=258.15, \quad y_{3}=258.141, \quad y_{4}=258.1409, \quad y_{5}=258.14090, \quad \ldots
$$

The sequence $\left(x_{n}\right)$ is non-decreasing: $x_{n+1} \geqq x_{n}$ for all $n$ and the sequence $\left(y_{n}\right)$ is non-increasing: $y_{n+1} \leqq y_{n}$ for all $n$. We have the sandwiching property: $x_{n} \leqq x \leqq y_{n}$ for all $n$. Both sequences are Cauchy because

$$
\begin{aligned}
& d\left(x_{n+i}, x_{n+j}\right)=\left|x_{n+i}-x_{n+j}\right| \leqq 10^{-n} \rightarrow 0 \quad(n \rightarrow \infty) \\
& d\left(y_{n+i}, y_{n+j}\right)=\left|y_{n+i}-y_{n+j}\right| \leqq 10^{-n} \rightarrow 0 \quad(n \rightarrow \infty)
\end{aligned}
$$

It follows that both sequences have limits. It is obvious that $x=\lim _{n \rightarrow \infty} x_{n}=\lim _{m \rightarrow \infty} y_{m}$.
What just has been illustrated is that there a natural way to construct for a given $x \in \mathbb{R}$ Cauchy sequences that converge towards $x$. The completeness principle states that the reverse is true: For any Cauchy sequence there is an element $x$ towards which the sequence converges.

We won't really talk about completeness in general until the chapter on compact spaces. Just to mention one of the simplest facts about completeness:

Theorem 10.10 (Complete sets are closed). Any complete subset of a metric space is closed.

Figure 10.5: complete $\Rightarrow$ closed


PROOF: Let $(X, d)$ be a metric space and $A \subseteq X$. Let $a \in X$ be a contact point of $A$. The theorem is proved if we can show that $a \in A$.
a) We employ def.10.25 on p.201: A point $x \in X$ is a contact point of $A$ if and only if $A \cap V \neq \emptyset$ for any neighborhood $V$ of $x$.

Let $m \in \mathbb{N}$. Then $N_{1 / m}(a)$ is a neighborhood of the contact point a, hence hence $A \bigcap N_{1 / m}(a) \neq \emptyset$ and we can pick a point from this intersection which we name $x_{m}$.
b) We prove next that $\left(x_{m}\right)_{m}$ is Cauchy. Let $\varepsilon>0$ and let $N \in \mathbb{N}$ be such that $N>\frac{1}{2 \varepsilon}$. if $j \in \mathbb{N}$ and $k \in \mathbb{N}$ both exceed $N$ then

$$
d\left(x_{j}, x_{k}\right) \leqq d\left(x_{j}, a\right)+d\left(a, x_{k}\right) \leqq \frac{1}{j}+\frac{1}{k} \leqq \frac{1}{N}+\frac{1}{N}<\varepsilon
$$

It follows that the sequence $\left(x_{j}\right)$ is Cauchy.
c) Because $A$ is complete, this sequence must converge to some $b \in A$. But $b$ cannot be different from $a$. Otherwise we could "separate" a and b by two disjoint neighborhoods: choose the open $\rho$-balls $N_{\rho}(a)$ and $N_{\rho}(b)$ where $\rho$ is one half the distance between a and $b$ (see the proof of thm.10.3 on $p .192$ ).

Only finitely many of the $x_{n}$ are allowed to be outside $N_{\rho}(a)$ and the same is true for $N_{\rho}(b)$. That is a contradiction and it follows that $b=a$, i.e., $a \in A$.
d) We summarize: if $a$ is a contact point of $A$ then $a \in A$. It follows that $A$ is closed.

The following is the reverse of thm.10.10.
Theorem 10.11 (Closed subsets of a complete space are complete). Let $(X, d)$ be a complete metric space and let $A \subseteq X$ be closed. Then $A$ is complete, i.e., the metric subspace $\left(A,\left.d\right|_{A \times A}\right)$ is complete.

PROOF: Let $\left(x_{n}\right)_{n}$ be Cauchy in $A$. We must show that there is $a \in A$ such that $x_{n} \rightarrow a$. Note that $\left(x_{n}\right)$ also is Cauchy in $X$ because the Cauchy criterion is entirely specified in terms of members of the sequence $\left(x_{n}\right)$.

Because $X$ is complete there exists $x \in X$ such that $x_{n} \rightarrow x$. All $x_{n}$ belong to $A$. According to thm.10.6 (Sequence criterion for contact points in metric spaces), $x$ is a contact point of $A$.

As the set $A$ is assumed to be closed, it contains all its contact points. It follows that $x \in A$, i.e., the arbitrary Cauchy sequence $\left(x_{n}\right)$ in $A$ converges to an element of $A$. We conclude that $A$ is complete.

Theorem 10.12 (Convergent sequences are Cauchy). Let $\left(x_{n}\right)_{n}$ be a convergent sequence in a metric space $(X, d)$. Then $\left(x_{n}\right)_{n}$ is Cauchy.

PROOF: Let $L \in X$ and $x_{n} \rightarrow L$. Let $\varepsilon>0$. There exists $N \in \mathbb{N}$ such that

$$
\begin{equation*}
k \geqq N \Rightarrow d\left(x_{k}, L\right)<\varepsilon / 2 . \tag{10.46}
\end{equation*}
$$

It follows from (10.46) that, for any $i, j \geqq N$,

$$
d\left(x_{i}, x_{j}\right) \leqq d\left(x_{i}, L\right)+d\left(L, x_{j}\right)<\varepsilon / 2+\varepsilon / 2=\varepsilon .
$$

It follows that the sequence satisfies (10.44) of the definition of a Cauchy sequence (def. 10.28 on p.205).
Proposition 10.21. Let $\left(x_{n}\right)_{n}$ be a Cauchy sequence in a metric space $(X, d)$ such that some subsequence $x_{n_{j}}$ converges to a limit $L$. Then ANY subsequence of $\left(x_{n}\right)_{n}$ converges to $L$. This is true in particular for the full sequence $x_{1}, x_{2}, \ldots, i . e .,\left(x_{n}\right)_{n}$ is a convergent sequence.

PROOF: Let $n_{1}<n_{2}<n_{3} \ldots$ be such that $x_{n_{j}}$ converges to $L$. For $k \in \mathbb{N}$ let $y_{k}:=x_{n_{k}}$.
Let $\varepsilon>0$. Convergence $y_{j} \rightarrow L$ implies that there is $N \in \mathbb{N}$ such that

$$
\begin{equation*}
d\left(y_{j}, L\right)<\varepsilon / 2 \text { for all } j \geqq N . \tag{10.47}
\end{equation*}
$$

Because $\left(x_{j}\right)$ is Cauchy there also exists $N^{\prime} \in \mathbb{N}$ such that

$$
\begin{equation*}
d\left(x_{i}, x_{j}\right)<\varepsilon / 2 \text { for all } i, j \geqq N^{\prime} . \tag{10.48}
\end{equation*}
$$

Let $K:=\max \left(n_{N}, N^{\prime}\right)$ and $j \geqq K$. Then

$$
d\left(x_{j}, L\right) \leqq d\left(x_{j}, y_{K}\right)+d\left(y_{K}, L\right)
$$

It follows from $n_{K} \geqq K$ and $j \geqq K$ and (10.48) that $d\left(x_{j}, y_{K}\right)=d\left(x_{j}, x_{n_{K}}\right)<\varepsilon / 2$ and it follows from (10.47) that $d\left(y_{K}, L\right)<\varepsilon / 2$. We conclude that $d\left(x_{j}, L\right)<\varepsilon$ for all $j \geqq K$ and this proves convergence $x_{j} \rightarrow L$.

### 10.2 Continuity (Study this!)

### 10.2.1 Definition and Characterizations of Continuous Functions

We have briefly discussed in ch. 8.2 on $p .133$. the continuity of functions with arguments and values in $\mathbb{R}$. We now extend this definition to functions that map from metric spaces to metric spaces.

Definition 10.30 (Sequence continuity). Given are two metric spaces ( $X, d_{1}$ ) and ( $Y, d_{2}$ ). Let $A \subseteq X$, $x_{0} \in A$ and let $f: A \rightarrow Y$ be a mapping from $A$ to $Y$. We say that $f$ is sequence continuous at $x_{0}$ and we write

$$
\begin{equation*}
\lim _{x \rightarrow x_{0}} f(x)=f\left(x_{0}\right) \tag{10.49}
\end{equation*}
$$

if the following is true for any sequence $\left(x_{n}\right)$ with values in $A$ :

$$
\begin{equation*}
\text { if } x_{n} \rightarrow x_{0} \text { then } f\left(x_{n}\right) \rightarrow f\left(x_{0}\right) . \tag{10.50}
\end{equation*}
$$

In other words, the following must be true for any sequence $\left(x_{n}\right)$ in $A$ and $x_{0} \in A$ :

$$
\begin{equation*}
\lim _{n \rightarrow \infty} x_{n}=x_{0} \Rightarrow \lim _{n \rightarrow \infty} f\left(x_{n}\right)=f\left(\lim _{n \rightarrow \infty} x_{n}\right)=f\left(x_{0}\right) . \tag{10.51}
\end{equation*}
$$

We say that $f$ is sequence continuous if $f$ is sequence continuous at $x_{0}$ for all $x_{0} \in A$.
Remark 10.18. Important points to notice:
a) It is not enough for the above to be true for some sequences that converge to $x_{0}$. Rather, it must be true for all such sequences!
b) We restrict our universe to the domain $A$ of $f: x_{0}$ and the entire sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ must belong to $A$ because there must be function values for all $x$-values. In other words, $f$ is continuous at $x_{0} \in A$ if and only if $f$ is continuous at $x_{0}$ in the metric subspace $\left(A,\left.d\right|_{A \times A}\right)$.

Definition 10.31 ( $\varepsilon-\delta$ continuity). Given are two metric spaces $\left(X, d_{1}\right)$ and $\left(Y, d_{2}\right)$. Let $A \subseteq X$, $x_{0} \in A$ and let $f(\cdot): A \rightarrow Y$ be a mapping from $A$ to $Y$. We say that $f(\cdot)$ is $\varepsilon-\delta$ continuous at $x_{0}$ if the following is true: For any (whatever small) $\varepsilon>0$ there exists $\delta>0$ such that either one of the following equivalent statements is satisfied:

$$
\begin{equation*}
f\left(N_{\delta}\left(x_{0}\right) \cap A\right) \subseteq N_{\varepsilon}\left(f\left(x_{0}\right)\right) \tag{10.52}
\end{equation*}
$$

$$
\begin{equation*}
d_{1}\left(x, x_{0}\right)<\delta \Rightarrow d_{2}\left(f(x), f\left(x_{0}\right)\right)<\varepsilon \text { for all } x \in A \tag{10.53}
\end{equation*}
$$

We say that $f(\cdot)$ is $\varepsilon-\delta$ continuous if $f(\cdot)$ is $\varepsilon-\delta$ continuous at $a$ for all $a \in A$.

Figure 10.6: $\varepsilon-\delta$ continuity


Remark 10.19. We recall from thm. 10.32 on p. 199 that

$$
\begin{equation*}
N_{\delta} \cap A=N_{\delta}^{A}(a)=\left\{x \in A:\left.d\right|_{A \times A}(x, a)<\delta\right\} . \tag{10.54}
\end{equation*}
$$

Hence, (10.52) states that $f$ is $\varepsilon-\delta$ continuous at $x_{0}$ if and only if $f\left(N_{\delta}^{A}\left(x_{0}\right)\right) \subseteq N_{\varepsilon}\left(f\left(x_{0}\right)\right)$.
Theorem 10.13 (Continuity criterion). Let $\left(X, d_{1}\right)$ and $\left(Y, d_{2}\right)$ be two metric spaces. Let $A \subseteq X, x_{0} \in A$ and let $f(\cdot): A \rightarrow Y$ be a mapping from $A$ to $Y$. Then $f$ is sequence continuous at $x_{0}$ if and only if $f$ is $\varepsilon-\delta$ continuous at $x_{0}$.

In particular $f$ is sequence continuous if and only if $f$ is $\varepsilon-\delta$ continuous.
a) $\Rightarrow$ : Proof that sequence continuity implies $\varepsilon$ - $\delta$-continuity:

We assume to the contrary that there exists some function $f$ which is sequence continuous but not $\varepsilon-\delta$ continuous at $x_{0}$, i.e., there exists some $\varepsilon>0$ such that neither (10.52) nor the equivalent (10.53) is true for any $\delta>0$.
a.1. In other words, No matter how small $\delta$ is choosen, there is at least one $x=x(\delta) \in A$ such that $d_{1}\left(x, x_{0}\right)<\delta$ but $d_{2}\left(f(x), f\left(x_{0}\right)\right) \geqq \varepsilon$. In particular we obtain for $\delta:=1 / m(m \in \mathbb{N})$ that

$$
\begin{equation*}
\text { there exists some } x_{m} \in N_{1 / m}\left(x_{0}\right) \cap A \text { such that } d_{2}\left(f\left(x_{m}\right), f\left(x_{0}\right)\right) \geqq \varepsilon \text {. } \tag{10.55}
\end{equation*}
$$

a.2. We now show that the sequence $\left(x_{m}\right)_{m \in \mathbb{N}}$ converges to $x_{0}$ : Let $\gamma>0$. There exists $N=N(\gamma) \in \mathbb{N}$ so big that $N>1 / \gamma$, i.e., $1 / N<\gamma$. As $x_{m} \in N_{1 / m}\left(x_{0}\right)$, we obtain for all $m \geqq N$ that

$$
d_{1}\left(x_{m}, x_{0}\right)<1 / m \leqq 1 / N<\gamma .
$$

This proves that $x_{m} \rightarrow x_{0}$.
a.3. Clearly, the sequence $\left(f\left(x_{m}\right)\right)_{m \in \mathbb{N}}$ does not converge to $f\left(x_{0}\right)$ as that requires $d_{2}\left(f\left(x_{m}\right), f\left(x_{0}\right)\right)<\varepsilon$ for all sufficiently big $m$, contrary to (10.55) which implies that there is not even one such $m$. In other words, the function $f$ is not sequence continuous, contrary to our assumption. We have our contradiction.
b) $\Leftarrow$ : Proof that $\varepsilon$ - $\delta$-continuity implies sequence continuity:

Let $x_{n} \rightarrow x_{0}$. Let $y_{n}:=f\left(x_{n}\right)$ and $y_{0}:=f\left(x_{0}\right)$. We must prove that $y_{n} \rightarrow y_{0}$ as $n \rightarrow \infty$.
b.1. Let $\varepsilon>0$. We can find $\delta>0$ such that (10.52) and hence (10.53) is satisfied. We assumed that $x_{n} \rightarrow x_{0}$. Hence there exists $N:=N(\delta) \in \mathbb{N}$ such that $d_{1}\left(x_{n}, x_{0}\right)<\delta$ for all $n \geqq N$.
b.2. It follows from (10.53) that $d_{2}\left(y_{n}, y\right)=d_{2}\left(f\left(x_{n}\right), f\left(x_{0}\right)\right)<\varepsilon$ for all $n \geqq N$. In other words, $y_{n} \rightarrow y_{0}$ as $n \rightarrow \infty$ and the proof of " $\Leftarrow$ " is finished.

Finally, the equivalence $f$ is $\varepsilon-\delta$ continuous $\Leftrightarrow f$ is $\varepsilon-\delta$ continuous.

## Definition (Continuity in metric spaces):

From now on we can use the terms " $\varepsilon-\delta$ continuous at $x_{0}$ " and "sequence continuous at $x_{0}$ " interchangeably for functions between metric spaces and we will simply speak about continuity of $f$ at $x_{0}$.

We saw in the $\varepsilon-\delta$ continuity definition of a function with metric spaces for both domain and codomain and the subsequent remark 10.19 that continuity of $f:\left(A,\left.d_{1}\right|_{A \times A}\right) \rightarrow\left(Y, d_{2}\right)$ in $x_{0} \in A$ was equivalent to demanding that for any $\varepsilon$-neighborhood of $f\left(x_{0}\right)$ there is a $\delta$-neighborhood of $x_{0}$ such that

$$
f\left(N_{\delta}^{A}\left(x_{0}\right)\right) \subseteq N_{\varepsilon}\left(f\left(x_{0}\right)\right) .
$$

The fact that any neighborhood of a point $z$ in a metric space contains a $\gamma$-neighborhood of $z$ for suitably small $\gamma$, is at the basis of the following theorem.

Theorem 10.14 (Neighborhood characterization of continuity). Let $\left(X, d_{1}\right)$ and $\left(Y, d_{2}\right)$ be two metric spaces. Let $A \subseteq X, x_{0} \in A$, and let $f(\cdot): A \rightarrow Y$ be a mapping from $A$ to $Y$. Then $f$ is continuous at $x_{0}$ if and only if for any neighborhood $V_{f\left(x_{0}\right)}$ of $f\left(x_{0}\right)$ there exists a neighborhood $U_{x_{0}}$ of $x_{0}$ in the metric space $\left(X, d_{1}\right)$ such that

$$
\begin{equation*}
f\left(U_{x_{0}} \cap A\right) \subseteq V_{f\left(x_{0}\right)} . \tag{10.56}
\end{equation*}
$$

Equivalently, this can be stated in terms of the subspace $\left(A,\left.d_{1}\right|_{A \times A}\right)$ as follows. for any neighborhood $V_{f\left(x_{0}\right)}$ of $f\left(x_{0}\right)$ there exists a neighborhood $U_{x_{0}}^{A}$ of $x_{0}$ in the metric space $\left(A,\left.d_{1}\right|_{A \times A}\right)$ such that

$$
\begin{equation*}
f\left(U_{x_{0}}^{A}\right) \subseteq V_{f\left(x_{0}\right)} . \tag{10.57}
\end{equation*}
$$

$\boldsymbol{a}) \Rightarrow):$ Assume that $f$ is continuous, i.e., $\varepsilon-\delta$ continuous at $a$. Let $V_{f\left(x_{0}\right)}$ be a neighborhood of $f\left(x_{0}\right)$.
Then $f\left(x_{0}\right)$ is interior point of $V_{f\left(x_{0}\right)}$ and we can find suitable $\varepsilon>0$ such that $N_{\varepsilon}\left(f\left(x_{0}\right)\right) \subseteq V_{f\left(x_{0}\right)} . \varepsilon-\delta$ continuity at a implies the existence of $\delta>0$ such that $f\left(N_{\delta}\left(x_{0}\right) \cap A\right) \subseteq V_{f\left(x_{0}\right)}$.

This proves both (10.56) (choose $U_{x_{0}}:=N_{\delta}\left(x_{0}\right)$ ) and (10.57) (choose $U_{x_{0}}^{A}:=N_{\delta}\left(x_{0}\right) \cap A$ ).
$b) \Leftarrow)$ : Assume that (10.56) is satisfied for any arbitrary neighborhood $V_{f\left(x_{0}\right)}$ of $f\left(x_{0}\right)$.
Let $\varepsilon>0$. We need to show that there exists $\delta>0$ such that

$$
\begin{equation*}
f\left(N_{\delta}\left(x_{0}\right) \cap A\right) \subseteq N_{\varepsilon}\left(f\left(x_{0}\right)\right) . \tag{10.58}
\end{equation*}
$$

$N_{\varepsilon}\left(f\left(x_{0}\right)\right)$ is a neighborhood of $f\left(x_{0}\right)$. It follows from (10.56) that there exists a neighborhood $U_{x_{0}}$ of $x_{0}$ such that

$$
\begin{equation*}
f\left(U_{x_{0}} \cap A\right) \subseteq N_{\varepsilon}\left(f\left(x_{0}\right)\right) \tag{10.59}
\end{equation*}
$$

$x_{0}$ is interior point of any of its neighborhoods. In particular, it is interior to $U_{x_{0}}$.

Accordingly, there exists $\delta>0$ such that $N_{\delta}\left(x_{0}\right) \subseteq U_{x_{0}}$, hence $N_{\delta}\left(x_{0}\right) \cap A \subseteq U_{x_{0}} \cap A$. It follows from the monotonicity of the direct image $\Gamma \mapsto f(\Gamma)$ that

$$
\begin{equation*}
f\left(N_{\delta}\left(x_{0}\right) \cap A\right) \subseteq f\left(U_{x_{0}} \cap A\right) \subseteq N_{\varepsilon}\left(f\left(x_{0}\right)\right) \tag{10.60}
\end{equation*}
$$

The second inclusion relation follows from (10.59). We have proved the existence of $\delta>0$ such that (10.58) is satisfied. This finishes the proof of " $\Leftarrow$ )".

The last theorem allows us to define continuity for functions between abstract topological spaces.
Definition 10.32 (Continuity for topological spaces). Given are two topological spaces ( $X, \mathfrak{L}_{1}$ ) and $\left(Y, \mathfrak{L}_{2}\right)$. Let $A \subseteq X, x_{0} \in A$ and let $f(\cdot): A \rightarrow Y$ be a mapping from $A$ to $Y$. We say that $f(\cdot)$ is continuous at $x_{0}$ if the following is true:

For any neighborhood $V_{f\left(x_{0}\right)}$ of $f\left(x_{0}\right)$ there exists a neighborhood $U_{x_{0}}$ of $x_{0}$ in the topological space ( $X, \mathfrak{U}_{1}$ ) such that

$$
\begin{equation*}
f\left(U_{x_{0}} \cap A\right) \subseteq V_{f\left(x_{0}\right)} \tag{10.61}
\end{equation*}
$$

Equivalently, this can be stated in terms of the subspace $\left(A, \mathfrak{U}_{1_{A}}\right)$ as follows. For any neighborhood $V_{f\left(x_{0}\right)}$ of $f\left(x_{0}\right)$ there exists a neighborhood $U_{x_{0}}^{A}$ of $x_{0}$ in $\left(A, \mathfrak{U}_{1_{A}}\right)$ such that

$$
\begin{equation*}
f\left(U_{x_{0}}^{A}\right) \subseteq V_{f\left(x_{0}\right)} \tag{10.62}
\end{equation*}
$$

We say that $f(\cdot)$ is continuous if $f(\cdot)$ is continuous at $a$ for all $a \in A$.
[1] B/G: Art of Proof defines in appendix A, p.136, continuity of a function $f$ as follows: " $f^{-1}$ (open) $=$ open". The following proposition proves that their definition coincides with the one given here: the validity of (10.52) for all $x_{0} \in X$.

Proposition 10.22 (" $f^{-1}$ (open) = open" continuity). Let ( $X, \mathfrak{U}$ ) and $(Y, \mathfrak{V})$ be two topological spaces and let $f(\cdot): X \rightarrow Y$ be a mapping from $X$ to $Y$. Then $f(\cdot)$ is continuous if and only if the following is true: If $V$ is open in $Y$ then the inverse image $f^{-1}(V)$ is open in $X$. ${ }^{104}$

PROOF of " $\Rightarrow$ ": Let $V$ be an open set in $Y$. Let $U:=f^{-1}(V), a \in U$ and $b:=f(a)$. Then $b \in V$ by the definition of inverse images. $b$ is inner point of $V$ because $V$ is open. According to def.10.32 there exists a neighborhood $U_{a}$ of a such that $f\left(U_{a}\right) \subseteq V$.

We conclude from the monotonicity of direct and inverse images and prop. 6.1 on $p .107$ that

$$
U_{a} \subseteq f^{-1}\left(f\left(U_{a}\right)\right) \subseteq f^{-1}(V)=U
$$

It follows that the arbitrarily chosen $a \in U$ is an interior point of $U$ and this proves that $U$ is open. PROOF of " $\Leftarrow$ ": We now assume that all inverse images of open sets in $Y$ are open in $X$.

Let $a \in X, b=f(a)$, and let $V_{b}$ be a neighborhood of $b$. Any neighborhood of $b$ contains an open neighborhood of $b$, hence we may assume that $V_{b}$ is open. We are done if we can find an open neighborhood $U_{a}$ of a such that

$$
\begin{equation*}
f\left(U_{a}\right) \subseteq V_{b} \tag{10.63}
\end{equation*}
$$

${ }^{104}$ This is easily extended to $f: A \rightarrow Y(\emptyset \neq A \subseteq X)$ by demanding that $f^{-1}(V)$ is open in $\left(A, \mathfrak{U}_{A}\right)$.

Let $U_{a}:=f^{-1}\left(V_{b}\right)$. Then $U_{a}$ is open as the inverse image of the open set $V_{b}$ It follows from the monotonicity of direct and inverse images and prop. 6.6 on $p .107$ that

$$
f(U)=f\left(f^{-1}\left(V_{b}\right)\right)=V_{b} \cap f(X) \subseteq V_{b}
$$

We have proved (10.63)

## Remark:

Note the following for the previous proposition 10.22.
a. In the interest of simplicity $f$ was defined on all of $X$ and not more generally on some arbitrary subset $A$ of $X$.
Note that the general case of $f: A \rightarrow Y$ is covered by replacing $(X, \mathfrak{U})$ with $\left(A, \mathfrak{U}_{A}\right)$, i.e., we deal with $f:\left(A, \mathfrak{U}_{A}\right) \rightarrow(Y, \mathfrak{V})$.
b. Also note that the above proposition only addresses "global" continuity of $f$ for all $x \in X$ and there is no local version which handles continuity at a specific $x_{0}$.

Remark 10.20 (continuity for real valued functions of real numbers). Let $\left(X, d_{1}\right)=\left(Y, d_{2}\right)=\mathbb{R}$. In this case equation (10.53) on p. 212 becomes

$$
\left|x-x_{0}\right|<\delta \Rightarrow\left|f(x)-f\left(x_{0}\right)\right|<\varepsilon
$$

See thm.8.1 on p.140.
Proposition 10.23 (continuity of the identity mapping). Let $X$ be nonempty and let

$$
i d_{X}: X \rightarrow X ; \quad x \mapsto x
$$

be the identity function on $X$. Then $i d_{X}$ is continuous in the following two cases:
a. $X$ is a metric space $(X, d)$,
b. $\quad X$ is a topological space $(X, \mathfrak{U})$.

PROOF of $\boldsymbol{b}$ : It suffices to show that inverse images of open sets are open. So let $V \in \mathfrak{U}$. Then $i d_{X}^{-1}(V)=V$, hence $i d_{X}^{-1}(V)$ is open. This finishes the proof.
PROOF of $\boldsymbol{a}$ : We associate with $d(\cdot, \cdot)$ its metric topology $\mathfrak{U}_{d}=\{U \subseteq X: U$ open in $(X, d)\}$ (see def. 10.13 (Topology induced by a metric) on p.194). We apply part $\boldsymbol{b}$ which we have already proved to ( $X, \mathfrak{U}_{d}$ ) and conclude that $i d_{X}$ is continuous.
Alternate proof of $\boldsymbol{a}$ : Here is a direct proof of $\boldsymbol{a}$ which works with the metric $d$.
Let $x \in X$ and $\varepsilon>0$. let $\delta:=\varepsilon$. If $x^{\prime} \in X$ such that $d\left(x, x^{\prime}\right)<\delta$, then

$$
d\left(i d_{X}(x), i d_{X}\left(x^{\prime}\right)\right)=d\left(x, x^{\prime}\right)<\delta=\varepsilon
$$

We have verified condition (10.53) of the $\varepsilon-\delta$ characterization of continuity and it follows that $i_{X}$ is continuous at $x .{ }^{105} x$ was an arbitrary point in $X$ and it follows that the identity is continuous.

Which one of the two proofs of $\boldsymbol{a}$ do you find more appealing?

[^77]
### 10.2.2 Continuity of Constants and Sums and Products

For all the following, unless stated differently, let $(X, d)$ be a metric space and $A \subseteq X, A \neq \emptyset$. Let

$$
f: A \rightarrow \mathbb{R}, \quad g: A \rightarrow \mathbb{R}
$$

be two real valued functions which both are continuous in a point $x_{0} \in A$. Moreover, let $a, b \in \mathbb{R}$. You can think of any fixed number a as a function

$$
a(\cdot): A \rightarrow \mathbb{R} ; x \mapsto a .
$$

In other words, the function $a(\cdot)$ assigns to each $x \in X$ one and the same value $a$. We called such a function a constant function (see (4.14) on p.84).

Proposition 10.24 (Rules of arithmetic for continuous real-valued functions). Given is a metric space ( $X, d$ ). Let the functions

$$
f(\cdot), g(\cdot), f_{1}(\cdot), f_{2}(\cdot), f_{3}(\cdot), \ldots, f_{n}(\cdot): A \longrightarrow \mathbb{R}
$$

all be continuous at $x_{0} \in A \subseteq X$. Then
a. Constant functions are continuous everywhere on $A$.
b. The product $f g(\cdot): x \mapsto f(x) g(x)$ is continuous at $x_{0}$. Specifically, af( $\left.\cdot\right): x \mapsto a \cdot f(x)$ is continuous at $x_{0}$ and, using -1 as a constant, $-f(\cdot): x \mapsto-f(x)$ is continuous at $x_{0}$.
c. The sum $f+g(\cdot): x \mapsto f(x)+g(x)$ is continuous at $x_{0}$.
d. Any linear combination ${ }^{106} \sum_{j=0}^{n} a_{j} f_{j}(\cdot): x \mapsto \sum_{j=0}^{n} a_{j} f_{j}(x)$ is continuous in $x_{0}$.

PROOF of $a$ : Let $\varepsilon>0$. We do not even have to look for a suitable $\delta$ to restrict the distance between two arguments $x$ and $x_{0}$ because it is always true that $\left|a(x)-a\left(x_{0}\right)\right|=|a-a|=0<\varepsilon$ This proves $\boldsymbol{a}$.
PROOF of $\boldsymbol{b}$ : Let $x_{n} \in A$ for all $n \in \mathbb{N}$ such that $x_{n} \rightarrow x_{0}$ as $n \rightarrow \infty$. All we need to show is $f\left(x_{n}\right) g\left(x_{n}\right) \rightarrow$ $f\left(x_{0}\right) g\left(x_{0}\right)$. This follows from prop.8.10 (Rules of arithmetic for limits) on p. 136 and we have proved $\boldsymbol{b}$.
PROOF of $\boldsymbol{c}$ : Let $x_{n} \in A$ for all $n \in \mathbb{N}$ such that $x_{n} \rightarrow x_{0}$ as $n \rightarrow \infty$. We must show $f\left(x_{n}\right)+g\left(x_{n}\right) \rightarrow$ $f\left(x_{0}\right)+g\left(x_{0}\right)$. This also follows from prop.8.10 and we have proved $\boldsymbol{c}$.
proof of $\boldsymbol{d}$ (outline): The proof is done by (strong) induction.
Base case: For $n=2$ the proof is obvious from parts $\boldsymbol{a}, \boldsymbol{b}$ and $\boldsymbol{c}$.
Induction step: Write

$$
\sum_{j=0}^{n+1} a_{j} f_{j}(x)=\left(\sum_{j=0}^{n} a_{j} f_{j}(x)\right)+a_{n+1} f_{n+1}(x)=I+I I .
$$

The left term " $I$ " is continuous by the induction assumption and the entire sum $I+I I$ then is continuous as the sum of two continuous functions (proved in $c$ ). This proves $d$.

Example 10.9 (Vector spaces of continuous real valued functions). Let ( $X, d$ ) be a metric space. Then

$$
\begin{equation*}
\mathscr{C}(X, \mathbb{R}):=\{f(\cdot): f(\cdot) \text { is a continuous real valued function on } X\} \tag{10.64}
\end{equation*}
$$

[^78]of all real continuous functions on $X$ is a vector space. Note that we have seen this before in example 9.11 (Vector spaces of real valued functions) on p .162 for the special case of $X \subseteq\left(\mathbb{R}, d_{|\cdot|}\right)$.

The sup-norm

$$
\|f(\cdot)\|_{\infty}=\sup \{|f(x)|: x \in X\}
$$

(see (9.13) on p.170) is not a real valued function on all of $\mathscr{C}(X, \mathbb{R})$ because $\|f(\cdot)\|_{\infty}=+\infty$ for any unbounded $f(\cdot) \in \mathscr{C}(X, \mathbb{R})$. To avoid complications from dealing with infinity, we often restrict the scope to the subspace

$$
\mathscr{C}_{\mathscr{B}}(X, \mathbb{R}):=\{h: h \text { is a bounded continuous real valued function on } X\}
$$

(see prop.9.10 on p .170 ) of the normed vector space $\mathscr{B}(X, \mathbb{R})$ of all bounded real valued functions on $X$. On this subspace the sup-norm truly is a real valued function in the sense that $\|f(\cdot)\|_{\infty}<\infty$.

### 10.2.3 Continuity of Polynomials (Understand this!)

Definition 10.3 (Polynomials). Anything that has to do with polynomials takes place in $\mathbb{R}$ and not on a metric space.

Let $A$ be subset of the real numbers and let $p(\cdot): A \rightarrow \mathbb{R}$ be a real valued function on $A$. p(.) is called a polynomial. if there is an integer $n \geqq 0$ and real numbers $a_{1}, a_{2}, \ldots, a_{n}$ which are constant (they do not depend on $x$ ) so that $p(\cdot)$ can be written as a sum

$$
\begin{equation*}
p(x)=a_{0}+a_{1} x+a_{2} x^{2}+\ldots+a_{n} x^{n}=\sum_{j=0}^{n} a_{j} x^{j} \tag{10.65}
\end{equation*}
$$

In other words, polynomials are linear combinations of the monomials $x \rightarrow x^{k} \quad\left(k \in(N)_{0}\right.$.
Proposition 10.25 (All polynomials are continuous).
The proof is left as exercise 10.24.
Proposition 10.26 (Vector space property of polynomials). Sums and scalar products of polynomials are polynomials.

## PROOF of $\boldsymbol{a}$. Additivity:

Let

$$
p_{1}(x)=a_{0}+a_{1} x+a_{2} x^{2}+\ldots+a_{n} x_{1}^{n}=\sum_{j=0}^{n_{1}} a_{j} x^{j}
$$

and

$$
p_{2}(x)=b_{0}+b_{1} x+b_{2} x^{2}+\ldots+b_{n} x_{2}^{n}=\sum_{j=0}^{n_{2}} b_{j} x^{j}
$$

be two polynomials. We may assume that $n_{1} \leqq n_{2}$. Let $a_{n_{1}+1}=a_{n_{1}+2}=\ldots=a_{n_{2}}=0$. Then $p_{1}(x)=$
$\sum_{j=0}^{n_{2}} a_{j} x^{j}$, hence

$$
\begin{aligned}
p_{1}(x)+p_{1}(x) & =\sum_{j=0}^{n_{2}} a_{j} x^{j}+\sum_{j=0}^{n_{2}} b_{j} x^{j} \\
& =\sum_{j=0}^{n_{2}}\left(a_{j}+b_{j}\right) x^{j} \\
& =\sum_{j=0}^{n_{2}} c_{j} x^{j} \quad\left(c_{j}:=a_{j}+b_{j}\right)
\end{aligned}
$$

This proves that the function $p_{1}(\cdot)+p_{2}(\cdot)$ is of the form (10.65) and we have shown that it is a polynomial. The proof for the sum of more than two polynomials is easily done by induction. See def.2.12 on p.19.
PROOF of b. Scalar product:
Let $p(x)=\sum_{j=0}^{n} a_{j} x^{j}$ be a polynomial. Let $\lambda \in \mathbb{R}$. Then

$$
(\lambda p)(x)=\lambda p(x)=\lambda \sum_{j=0}^{n} a_{j} x^{j}=\sum_{j=0}^{n} \lambda a_{j} x^{j}=\sum_{j=0}^{n} c_{j} x^{j} \quad\left(c_{j}:=\lambda a_{j}\right)
$$

This proves that the function $\lambda p(\cdot)$ is of the form (10.65).
Polnomials may not always be given in their normalized form (10.65) on p.218. Here is an example:

$$
\begin{aligned}
p(x) & =a_{0} x^{0}(1-x)^{n}+a_{1} x^{1}(1-x)^{n-1}+a_{2} x^{2}(1-x)^{n-2}+\ldots+a_{n-1} x^{n-1}(1-x)^{1}+a_{n} x^{n} \\
& =\sum_{k=0}^{n} a_{k} x^{k}(1-x)^{n-k}
\end{aligned}
$$

is a linear combination of monomials and hence a polynomial. All you need to do is "multiply out" the $x^{k}(1-x)^{n-k}$ terms and then regroup the resulting mess. The so called Bernstein polynomials ${ }^{107}$

$$
p(x)=\sum_{k=0}^{n}\binom{n}{k} f\left(\frac{k}{n}\right) x^{k}(1-x)^{n-k}
$$

are of that form.
Example 10.10 (Vector space of polynomials with sup-norm). Let $A \subseteq \mathbb{R}$. It follows from (10.26) and (10.25) that the set

$$
\{p(\cdot): p(\cdot) \text { is a polynomial on } A\}
$$

of all polynomials on an arbitrary nonempty subset $A$ of the real numbers is a subspace of the vector space $\mathscr{C}(A, \mathbb{R})$. (see example (10.9) on p.217. If $A$ is not bounded, the sup-norm

$$
\|p\|_{\infty}=\sup \{|p(x)|: x \in A\}
$$

is not a real valued function on the set of all polynomials on $A$ as its value may be $\infty$. Matter of fact, it can be shown that, if $A$ is not bounded, then the only polynomials for which $\|p(\cdot)\|_{\infty}<\infty$ are the constant functions on $A$.

[^79]
### 10.2.4 Uniform Continuity

It will be proved in theorem 11.12 (Uniform continuity on sequence compact spaces) on p. $257{ }^{108}$ that continuous real valued functions on the compact set $[0,1]$ are uniformly continuous in the sense of the following definition. ${ }^{109}$

Definition 10.34 (Uniform continuity of functions). Let $\left(X, d_{1}\right),\left(Y, d_{2}\right)$ be metric spaces and let $A$ be a subset of $X$. A function

$$
f(\cdot): A \rightarrow Y \quad \text { is called uniformly continuous }
$$

if for any $\varepsilon>0$ there exists a (possibly very small) $\delta>0$ such that

$$
\begin{equation*}
d_{2}(f(x)-f(y))<\varepsilon \quad \text { for any } x, y \in A \text { such that } d_{1}(x, y)<\delta \tag{10.66}
\end{equation*}
$$

Remark 10.21 (Uniform continuity vs. continuity). Note the following:
$\boldsymbol{a}$. Condition (10.66) for uniform continuity looks very close to the $\varepsilon-\delta$ characterization of ordinary continuity (10.53) on p.212. Can you spot the difference?

Uniform continuity is more demanding than plain continuity because, when dealing with the latter, you can ask for specific values of both $\varepsilon$ and $x_{0}$ according to which you must find a suitable $\delta$. In other words, for plain continuity

$$
\delta=\delta\left(\varepsilon, x_{0}\right) .
$$

In the case of uniform continuity all you get is $\varepsilon$. You must come up with a suitable $\delta$ regardless of what arguments are thrown at you. To write that one in functional notation,

$$
\delta=\delta(\varepsilon)
$$

b. It follows that uniform continuity implies continuity but the opposite need not be true.
c. Many concepts that are defined in metric spaces can be generalized to topological spaces Examples were neigborhoods, interior points and contact points, subspaces and continuity. Uniform continuity is not a concept that can be defined without a metric. ${ }^{110}$

Example 10.11 (Uniform continuity of the identity mapping). Let us have another look at proposition(10.23) where we proved the continuity of the identity mapping on a metric space. We chose $\delta=\varepsilon$ no matter what value of x we were dealing with and it follows that the identity mapping is always uniformly continuous.

Remark 10.22. Now that you have learned the definitions for both continuity and uniform continuity, have another look at example 3.30, p. 54 in ch.3.6.3 (Quantifiers for Statement Functions of more than Two Variables) where it was explained how you could obtain one definition from the other just by switching around $\mathrm{a} \forall$ quantifier and $\mathrm{a} \exists$ quantifier.

[^80]
### 10.2.5 Continuity of Linear Functions (Understand this!)

Lemma 10.1. Let $f:(V,\|\cdot\|) \rightarrow(W, \mathbf{I} \cdot \mathbf{I})$ be a linear function between two normed vector spaces. Let

$$
\begin{aligned}
a & :=\sup \{|f(x)|: x \in V,\|x\|=1\}, \\
b & :=\sup \{|f(x)|: x \in V,\|x\| \leqq 1\}, \\
c & :=\sup \left\{\frac{|f(x)|}{\|x\|}: x \in V, x \neq 0\right\} .
\end{aligned}
$$

Then $a=b=c$.
PROOF: We introduce the following three sets for this proof:

$$
\begin{aligned}
& A:=\{|f(x)|: x \in V,\|x\|=1\}, \\
& B:=\{|f(x)|: x \in V,\|x\| \leqq 1\}, \\
& C:=\left\{\frac{|f(x)|}{\|x\|}: x \in V, x \neq 0\right\} .
\end{aligned}
$$

Proof that $a=b$ :
It follows from $A \subseteq B$ that $a \leqq b$. On the other hand let $x \in B$ such that $x \neq 0$ (if $x=0$ then $f(x)=0$ certainly could not exceed $a$ ). Let $y:=\|x\|^{-1} x$. Then $y \in A$ and $\|x\|^{-1} \geqq 1$, hence

$$
|f(y)|=|f(x /\|x\|)|=(1 /\|x\|)|f(x)| \geqq|f(x)| .
$$

We conclude that the sup over the bigger set $B$ does not exceed the sup over $A$, hence $a=b$.
Proof that $a=c$ :
Let $x \in C$ and $y:=\|x\|^{-1} x$. Then $y \in A$ and

$$
|f(x)| /\|x\|=|f(x) /\|x\||=|f(x /\|x\|)|=|f(y)| .
$$

It follows that the sup over the bigger set $C$ does not exceed the sup over $A$, hence $c=b$.
Definition 10.35 (norm of linear functions). Let $f:(V,\|\cdot\|) \rightarrow(W, \mathbf{I} \cdot \mathbf{I})$ be a linear function between two normed vector spaces. We denote the quantity $a=b=c$ from lemma 10.1 by $\|f\|$, i.e.,

$$
\begin{align*}
\|f\| & =\sup \{|f(x)|: x \in V,\|x\|=1\} \\
& =\sup \{|f(x)|: x \in V,\|x\| \leqq 1\}  \tag{10.67}\\
& =\sup \left\{\frac{|f(x)|}{\|x\|}: x \in V, x \neq 0\right\} .
\end{align*}
$$

$\|f\|$ is called the norm of $f .{ }^{111}$
We note that $\|f\|$ need not be finite.

[^81]Theorem $\mathbf{1 0 . 1 5}$ (Continuity criterion for linear functions). Let $f:(V,\|\cdot\|) \rightarrow(W, \mathbf{I} \cdot \mathbf{I})$ be a linear function between two normed vector spaces. Then the following are equivalent.
A. $f$ is continuous at $x=0$,
B. $f$ is continuous in all points of $V$,
C. $f$ is uniformly continuous on $V$,
D. $\|f\|<\infty$.

Moreover, we then have

$$
\begin{equation*}
|f(x)| \leqq\|f\| \cdot\|x\| \quad \text { for all } x \in V \tag{10.68}
\end{equation*}
$$

PROOF: Clearly we have $\boldsymbol{C} \Rightarrow \boldsymbol{B} \Rightarrow \boldsymbol{A}$. We now show $\boldsymbol{A} \Rightarrow \boldsymbol{D}$.
It follows from the continuity of $f$ at 0 that there exists $\delta>0$ such that

$$
\begin{equation*}
\text { if } z \in V \text { and }\|z\|<\delta \text { then }|f(z)|=|f(z)-f(0)|<1 \tag{10.69}
\end{equation*}
$$

Let $x \in V$ such that $\|x\| \leqq 1$. Then $\|\delta / 2 \cdot x\| \leqq \delta / 2<\delta$, hence, according to (10.69),

$$
\delta / 2 \cdot|f(x)|=|f(\delta / 2 \cdot x)|<1, \text { hence }|f(x)|<2 / \delta
$$

Because this last inequality is true for all $x \in V$ with norm bounded by 1 , it follows that

$$
\|f\|=\sup \{|f(x)|: x \in V,\|x\| \leqq 1\}<2 / \delta<\infty
$$

We have proved that $\boldsymbol{A} \Rightarrow \boldsymbol{D}$.
We finally show $\boldsymbol{D} \Rightarrow \boldsymbol{C}$ and we do this in two steps.
First we show $\boldsymbol{D} \Rightarrow$ (10.68). The inequality trivially holds for $x=0$ because linearity of $f$ implies $f(0)=0$. If $x \neq 0$ then $\|x\|>0$ (norms are positive definite) and the inequality follows from the last characterization of $\|f\|$ in (10.67).

Second step: Let $\varepsilon>0$ and $\delta:=\varepsilon /\|f\|$. Let $x, y \in V$ such that $\|x-y\|<\delta$. If we can prove that this implies $|f(x)-f(y)|<\varepsilon$, then $f$ is indeed uniformly continuous and the proof is done. We show this as follows.

$$
|f(x)-f(y)|=|f(x-y)| \stackrel{(10.68)}{\leqq}\|f\| \cdot\|x-y\|<\|f\| \cdot \delta=\|f\| \cdot \varepsilon /\|f\|=\varepsilon
$$

Note 10.3 ( $\|f\|$ is a norm). Let
(10.70) $\mathscr{C}_{\text {lin }}(V, W):=\mathscr{C}_{\text {lin }}((V,\|\cdot\|),(W, \mathbf{I} \cdot \mathbf{I})):=\{f: V \rightarrow W: f$ is linear and continuous $\}$.

Then $\mathscr{C}_{\text {lin }}(V, W)$ is a vector space and

$$
\begin{equation*}
f \mapsto\|f\|=\sup \{|f(x)|:\|x\|=1\} \tag{10.71}
\end{equation*}
$$

defines a norm on $\mathscr{C}_{\text {lin }}(V, W)$.
In all of the proof that that is given now let $A:=x \in V:\|x\|=1$.
A. Proof that $\mathscr{C}_{\text {lin }}(V, W)$ is a vector space.

Let $f, g, \in \mathscr{C}_{\text {lin }}(V, W)$. We need to show that $f+g \in \mathscr{C}_{\text {lin }}(V, W)$, i.e., $f+g$ is both linear and continuous. Linearity is immediate. We now show continuity.

Let $x \in A$. Then

$$
\begin{equation*}
|f(x)+g(x)| \leqq|f(x)|+|f(x)| \leqq\|f\|+\|g\|<\infty . \tag{10.72}
\end{equation*}
$$

The first inequality holds because the norm $|f(x)|$ satisfies the triangle inequality for norms. The second follows from (10.67) on p.221, and the finiteness of $\|f\|+\|g\|$ is, according to the continuity criterion for linear functions (thm. 10.15 on $p .222$ ), equivalent to the continuity of both $f$ and $g$.

We still must show that if $f \in \mathscr{C}_{\text {lin }}(V, W)$ and $\lambda \in \mathbb{R}$ then $\lambda f: x \mapsto \lambda f(x) \in \mathscr{C}_{\text {lin }}(V, W)$, i.e., we must show that this function is linear and continuous. Again, linearity is immediate. To show continuity we proceed as follows.

Let $x \in$ A. I•I is absolutely homogeneous. Hence

$$
\begin{equation*}
|\lambda f(x)|=||\lambda| f(x)|=|\lambda| \cdot|f(x)| \tag{10.73}
\end{equation*}
$$

It follows from prop.8.4 (positive homogeneity of the sup) on p. 131 that

$$
\begin{align*}
\|\lambda f\| & =\sup \{|\lambda f(x)|:\|x\|=1\}=\sup \{| | \lambda|f(x)|:\|x\|=1\}  \tag{10.74}\\
& =|\lambda| \cdot \sup \{|f(x)|:\|x\|=1\}=|\lambda| \cdot\|f\|<\infty \tag{10.75}
\end{align*}
$$

This proves that $\lambda f$ is continuous.
B. Proof that $\|f\|$ is a norm on $\mathscr{C}_{\text {lin }}(V, W)$.

Because (10.72) is valid for all $x \in A$, we obtain

$$
\begin{equation*}
\|f+g\|=\sup \{|f(x)+g(x)|: x \in A\} \leqq\|f\|+\|g\| \tag{10.76}
\end{equation*}
$$

This proves the triangle inequality.
Likewise, we obtain from the validity of (10.74) for all $x \in A$,

$$
\begin{equation*}
\|\lambda f\|=\sup \{| | \lambda|f(x)|: x \in A\}=|\lambda| \sup \{|f(x)|: x \in A\}=|\lambda|\|f\| . \tag{10.77}
\end{equation*}
$$

This proves absolute homogeneity.
Finally we show positive definiteness. Clearly $\|f\|$ is non-negative as the sup of non-negative numbers $|f(x)|$. Assume that $\|f\|>0$. Then $\delta:=\frac{1}{2}\|f\|>0$ and there exists $x_{0} \in A$ such that
(10.78) $\sup \{|f(x)|: x \in A\}-\left|f\left(x_{0}\right)\right|<\delta$, i.e., $\|f\|-\left|f\left(x_{0}\right)\right|<\delta$, hence $\left|f\left(x_{0}\right)\right|>\delta$.

Positive definiteness of $\mathbf{I} \cdot \mathbf{I}$ implies that $f\left(x_{0}\right) \neq 0$ and hence $f \neq 0$. We have proved positive definiteness of $\|\cdot\|$.

### 10.3 Function Sequences and Infinite Series

### 10.3.1 Convergence of Function Sequences (Study this!)

## Notation Alert:

This chapter makes heary use of the notation $f(\cdot)$ instead of $f$ for a function $X \rightarrow \mathbb{R}$ to emphasize when sequences of functions $f_{n}(\cdot)$ are used and when function values (real numbers) $f_{n}(x)$ are used.

Vectors are more complicated than numbers because an $n$-dimensional vector $v \in \mathbb{R}^{n}$ represents a grouping of a finite number $n$ of real numbers. Matter of fact, any such vector $\left(x_{1}, x_{2}, x_{3}, \cdots, x_{n}\right)$ can be interpreted as a real valued function (remember: a real valued function is one which maps it arguments into $\mathbb{R}$ )

$$
\begin{equation*}
f(\cdot):\{1,2,3, \cdots, n\} \rightarrow \mathbb{R} \quad j \mapsto x_{j} \tag{10.79}
\end{equation*}
$$

(see (9.4) on p.156).
Next come sequences $\left(x_{j}\right)_{j \in \mathbb{N}}$ which can be interpreted as real valued functions

$$
\begin{equation*}
g(\cdot): \mathbb{N} \rightarrow \mathbb{R} \quad j \mapsto x_{j} . \tag{10.80}
\end{equation*}
$$

Finally we deal with any kind of real valued function

$$
\begin{equation*}
h(\cdot): X \rightarrow \mathbb{R} \quad x \mapsto h(x) \tag{10.81}
\end{equation*}
$$

as the most general case.
Now we add more complexity by not just dealing with one or two or three real valued functions but with an entire sequence of functions

$$
\begin{equation*}
f_{n}(\cdot): X \rightarrow \mathbb{R} \quad x \mapsto f_{n}(x) \tag{10.82}
\end{equation*}
$$

For any fixed argument $x_{0}$ we have a sequence $f_{1}\left(x_{0}\right), f_{2}\left(x_{0}\right), f_{3}\left(x_{0}\right), \cdots$ of real numbers which we can examine for convergence. This sequence may converge for some or all arguments $x_{0} \in X$ to some $L=$ $L\left(X_{0}\right) \in \mathbb{R}$. It's time now for some definitions.

Definition 10.36 (Pointwise convergence of function sequences). Let X be a nonempty set, $(Y, d)$ a metric space and let $f_{n}(\cdot): X \rightarrow Y$ and $f(\cdot): X \rightarrow Y$ be functions on $X(n \in \mathbb{N})$. Let $A \subseteq X$ be a subset of $X$. We say that $f_{n}(\cdot)$ converges pointwise or, simply, converges to $f(\cdot)$ on $A$ and we write $f_{n}(\cdot) \rightarrow f(\cdot)$ on $A$ if

$$
\begin{equation*}
f_{n}(x) \rightarrow f(x) \quad \text { for all } x \in A \tag{10.83}
\end{equation*}
$$

We omit the phrase "on $A$ " if it is clear how $A$ is defined, in particular if $A=X$.
Definition 10.37 (Uniform convergence of function sequences). Let X be a nonempty set, $(Y, d)$ a metric space and let $f_{n}(\cdot): X \rightarrow Y$ and $f(\cdot): X \rightarrow Y$ be functions on $X(n \in \mathbb{N})$. Let $A \subseteq X$ be a subset of $X$. We say that $f_{n}(\cdot)$ converges uniformly to $f(\cdot)$ on $A$ and we write ${ }^{112}$

$$
\begin{equation*}
f_{n}(\cdot) \xrightarrow{u c} f(\cdot) \text { on } A \tag{10.84}
\end{equation*}
$$

if the following is true: For each $\varepsilon>0$ (no matter how small) there exists a (possibly huge) number $n_{0}$ which can be chosen once and for all, independently of the specific argument $x$, such that

$$
\begin{equation*}
d\left(f_{n}(x), f(x)\right)<\varepsilon \quad \text { for all } x \in A \quad \text { and } n \geqq n_{0} . \tag{10.85}
\end{equation*}
$$

We omit the phrase "on $A$ " if it is clear how $A$ is defined, in particular if $A=X$.

[^82]Remark 10.23 (Uniform convergence implies pointwise convergence). Take another look at definition def.10.11 (convergence of sequences in metric spaces) on p.192. Note that (10.85) implies, for any given $x \in A$, ordinary convergence $f(x)=\lim _{n \rightarrow \infty} f_{n}(x)$. The reason is that the number $n_{0}=n_{0}(\varepsilon)$ chosen in (10.85) will also satisfy (10.20) of that definition for $x_{n}=f_{n}(x)$ and $a=f(x)$.

In other words, unform convergence implies pointwise convergence. But what is the difference between pointwise and uniform convergence? The difference is that, for poinwise convergence, the number $n_{0}$ will depend on both $\varepsilon$ and $x: n_{0}=n_{0}(\varepsilon, x)$. In the case of uniform convergence, the number $n_{0}$ will still depend on $\varepsilon$ but can be chosen independently of the argument $x \in A$.

Example 10.12 (a. Constant sequence of functions). Let X be a set and let $f: X \rightarrow \mathbb{R}$ be a real valued function on $X$ which may or may not be continuous anywhere. Define a sequence of functions

$$
f_{n}: X \rightarrow \mathbb{R}(n \in \mathbb{N}) \quad \text { as } \quad f_{1}=f_{2}=\cdots=f
$$

i.e.,

$$
f_{1}(x)=f_{2}(x)=\cdots=f(x) \forall n \in \mathbb{N}, \forall x \in X
$$

In other words, we are looking at a constant sequence of functions (not to be confused with a sequence of constant functions - seriously!).

Then $f_{n}(\cdot) \xrightarrow{u c} f(\cdot)$.
PROOF of the example: This is trivial. No matter how small an $\varepsilon$ and $n_{0}$ we choose and no matter what argument $x \in X$ we are looking at, we have

$$
\left|f_{n}(x)-f(x)\right|=0<\varepsilon \quad \text { for all } x \in A \quad \text { and } n>n_{0}
$$

Example 10.13 (b. Pointwise but not uniformly convergent sequence of functions). Let $X=[0,1]$, i.e., $X$ is the closed unit interval $\{x \in \mathbb{R}: 0 \leqq x \leqq 1\}$. Let the functions $f_{n}$ be defined as follows on $X$ :

$$
f_{n}(x)= \begin{cases}n^{2} x & \text { for } 0 \leqq x \leqq \frac{1}{n} \\ \frac{1}{x} & \text { for } \frac{1}{n} \leqq x \leqq 1\end{cases}
$$

Note that both pieces fit together in the point $x=1 / n$ because the " $\frac{1}{x}$ definition" gives $f_{n}(a)=$ $\frac{1}{1 / n}=n$ and the " $n^{2} x$ definition" gives the same value $n=n^{2} \frac{1}{n}$. We do not give a formal proof that $f_{n}(\cdot)$ is continuous at every point of $[0,1]$ but you should draw a picture of $f_{n}$ to see that the two graphs flow into each other at the "splicing point" $1 / n$. Instead you are asked in exercise 10.26 on p. 242 to give such a proof.

Now we define the function $f(\cdot):[0,1] \rightarrow \mathbb{R}$ as

$$
f(x)= \begin{cases}\frac{1}{x} & \text { for } 0<x \leqq 1 \\ 0 & \text { for } x=0\end{cases}
$$

Then the functions $f_{n}(\cdot)$ converge pointwise but not uniformly to $f(\cdot)$ on the entire unit interval.

PROOF of pointwise convergence:
first we inspect the point $a=0$. We have $f(0)=0=n^{2} 0=f^{n}(0)$ and the constant sequence of zeroes certainly converges to zero. Now assume $a>0$. If $n>1 / a$ then $f_{n}(a)=\frac{1}{a}$ for all such $n$. We have $a$
constant sequence ( $\frac{1}{a}$ ) except for the first finitely many $n$ and this sequence converges to $\frac{1}{a}=f(a)$. See prop.8.6 on p.135. We have thus proved pointwise convergence.

Proof that there is no uniform convergence:
To prove that (10.85) is not satisfied, we must find $\varepsilon>0$ and points $x_{N}$ so that for no matter how big a natural number $N$ we choose, there will be at least one $j>N$ such that $\left|f_{j}\left(x_{N}\right)-f\left(x_{N}\right)\right| \geqq \varepsilon$. Let $N \in \mathbb{N}$ be any natural number and let $x_{N}:=\frac{1}{N^{2}}$. Then

$$
\begin{aligned}
f_{N}\left(x_{N}\right) & =\frac{N^{2}}{N^{2}}=1 \\
f_{2 N}\left(x_{N}\right) & =\frac{(2 N)^{2}}{N^{2}}=4
\end{aligned}
$$

Hence

$$
\left|f_{2 N}\left(x_{N}\right)-f_{N}\left(x_{N}\right)\right|=3
$$

To recap: We found $\varepsilon>0$ so that for each $N \in \mathbb{N}$ there is at least one $j \geqq N$ and $x_{N} \in[0,1]$ such that $\left|f_{j}\left(X_{N}\right)-f_{N}\left(x_{N}\right)\right|>\varepsilon$ : we chose

$$
\varepsilon=2, \quad j=2 N, \quad x_{N}=\frac{1}{N^{2}}
$$

We have proved that convergence is pointwise but not uniform.
Proposition 10.27 (Uniform convergence is $\|\cdot\|_{\infty}$ convergence). Let $X$ be a nonempty set and $\mathscr{B}(X, \mathbb{R})$ the set of all bounded real valued functions on $X$. We remember that this set is a vector space with the norm $\|f\|_{\infty}=\sup \{|g(x)-f(x)|: x \in X\}$ and it is a metric space with the corresponding metric

$$
d_{\|\cdot\|_{\infty}}(f, g)=\sup \{|g(x)-f(x)|: x \in X\}
$$

(see example 10.2 on $p .185$ ). The following is true:

$$
f_{n}(\cdot) \xrightarrow{u c} f(\cdot) \quad \text { if and only if } \quad f_{n}(\cdot) \xrightarrow{\|\cdot\|_{\infty}} f(\cdot)
$$

In other words, the sequence $f_{n}$ converges to $f$ uniformly if and only if it converges to $f$ in the metric space $\left(\mathscr{B}(X, \mathbb{R}), d_{\|\cdot\|_{\infty}}(\cdot, \cdot)\right)$.

PROOF of " $\Rightarrow$ ": Assume that $f_{n}(\cdot) \xrightarrow{u c} f(\cdot)$. Let $\varepsilon>0$. According to def. 10.37 (Uniform convergence of function sequences) on $p .224$, there exists an index $n_{0}=n_{0}(\varepsilon)$ (which does not depend on the function argument $x \in X$ ) such that

$$
d\left(f_{n}(x), f(x)\right)=\left|f_{n}(x)-f(x)\right|<\varepsilon / 2 \quad \text { for all } x \in X \quad \text { and } n \geqq n_{0} .
$$

Note that here the metric space $Y$ in def. 10.37 is $\mathbb{R}$, so $d\left(f_{n}(x), f(x)\right)$ becomes $\left|f_{n}(x)-f(x)\right|$. We obtain

$$
\left\|f_{n}-f\right\|_{\infty}=\sup \left\{\left|f_{n}(x)-f(x)\right|: x \in X\right\} \leqq \varepsilon / 2 \text { for all } n \geqq n_{0}
$$

i.e., $d_{\|\cdot\|_{\infty}}\left(f_{n}, f\right)<\varepsilon$ for all $n \geqq n_{0}$. It follows that $f_{n}(\cdot) \xrightarrow{\|\cdot\|_{\infty}} f(\cdot)$.

PROOF of " $\Leftarrow$ ": Assume that $f_{n} \xrightarrow{\|\cdot\|_{\infty}}$ f, i.e., $\lim _{n \rightarrow \infty} f_{n}=f$ in the metric space $\left(\mathscr{B}(X, \mathbb{R}), d_{\|\cdot\|_{\infty}}\right)$.

Let $\varepsilon>0$. There exists $n_{0} \in \mathbb{N}$ such that

$$
d_{\|\cdot\|_{\infty}}\left(f_{n}, f\right)=\left\|f_{n}-f\right\|_{\infty}=\sup \left\{\left|f_{n}(x)-f(x)\right|: x \in X\right\}<\varepsilon \text { for all } n \geqq n_{0}
$$

But then

$$
\left|f_{n}(x)-f(x)\right|<\varepsilon \quad \text { for all } x \in X \text { and all } n \geqq n_{0} .
$$

This proves $f_{n}(\cdot) \xrightarrow{u c} f(\cdot)$.
The last proposition justifies the next definition.
Definition 10.38 (Norm and metric of uniform convergence). We also call the sup-norm on $\mathscr{B}(X, \mathbb{R})$ the norm of uniform convergence on $X$ and its associated metric $d_{\|\cdot\|_{\infty}}(\cdot, \cdot)$ the metric of uniform convergence on $X$.

Theorem 10.16 (Uniform limits of continuous functions are continuous). Let ( $X, d_{1}$ ) and $\left(Y, d_{2}\right)$ be metric spaces and let $f_{n}(\cdot): X \rightarrow Y$ and $f(\cdot): X \rightarrow Y$ be functions on $X(n \in \mathbb{N})$. Let $x_{0} \in X$ and let $V \subseteq X$ be a neighborhood of $x_{0}$. Assume $\boldsymbol{a}$ ) that the functions $f_{n}(\cdot)$ are continuous at $x_{0}$ for all $n$ and $\boldsymbol{b}$ ) that $f_{n}(\cdot) \xrightarrow{u c} f(\cdot)$ on $V$. Then $f$ is continuous at $x_{0}$

PROOF: Let $\varepsilon>0$.
A. Uniform convergence $f_{n}(\cdot) \xrightarrow{u c} f(\cdot)$ on $V$ guarantees the existence of some $N=N(\varepsilon)$ such that

$$
d_{2}\left(f_{n}(x), f(x)\right)<\frac{\varepsilon}{3} \text { for all } x \in V \text { and } n \geqq \mathbb{N} .
$$

In particular, for $n=N$,

$$
\begin{equation*}
d_{2}\left(f_{N}(x), f(x)\right)<\frac{\varepsilon}{3} \text { for all } x \in V \tag{10.86}
\end{equation*}
$$

B. All functions $f_{n}$ and in particular $f_{N}$ are continuous in $V$. There is $\tilde{\delta}>0$ such that

$$
\begin{equation*}
d_{2}\left(f_{N}(x), f_{N}\left(x_{0}\right)\right)<\frac{\varepsilon}{3} \text { for all } x \in N_{\tilde{\delta}}\left(x_{0}\right) . \tag{10.87}
\end{equation*}
$$

C. As $x_{0}$ is an interior point of $V$, there exists $\hat{\delta}>0$ such that $N_{\hat{\delta}}\left(x_{0}\right) \subseteq V$. Let $\delta$ be the smaller of $\hat{\delta}$ and $\tilde{\delta}$.

Then (10.86) and (10.87) both hold for any $x \in N_{\delta}\left(x_{0}\right)$. Because $x_{0} \in N_{\delta}\left(x_{0}\right)$ we obtain

$$
d\left(f(x), f\left(x_{0}\right)\right) \leqq d\left(f(x), f_{N}(x)\right)+d\left(f_{N}(x), f_{N}\left(x_{0}\right)\right)+d\left(f_{N}\left(x_{0}\right), f\left(x_{0}\right)\right)<\frac{\varepsilon}{3}+\frac{\varepsilon}{3}+\frac{\varepsilon}{3}=\varepsilon
$$

The proof is finished.

### 10.3.2 Infinite Series

We start by repeating the definition of a sequence given in section 4.2 on p.70: A sequence $\left(x_{j}\right)$ is nothing but a family of things $x_{j}$ which are indexed by integers, usually the natural numbers or the non-negative integers. We make throughout this entire document the following

Assumption 10.1 (indices of sequences).
Unless explicitly stated otherwise, sequences are always indexed $1,2,3, \ldots$, i.e., the first index is 1 and, given any index, you obtain the next one by adding 1 to it.

The simplest things that a mathematician deals with are numbers. One nice thing that is always possible with numbers, is that you can add them. Here is a very simple definition:

Definition 10.39 (Numeric Sequences and Series). A sequence ( $a_{j}$ ) of real numbers is called a numeric sequence if each $a_{j}$ is a real number. For any such sequence, we can build another sequence $\left(s_{n}\right)$ as follows:

$$
\begin{equation*}
s_{1}:=a_{1} ; \quad s_{2}:=a_{1}+a_{2} ; \quad s_{3}:=a_{1}+a_{2}+a_{3} ; \cdots \quad s_{n}:=\sum_{k=1}^{n} a_{k} \tag{10.88}
\end{equation*}
$$

We write this more compactly as

$$
\begin{equation*}
a_{1}+a_{2}+a_{3}+\cdots=\sum a_{k} \tag{10.89}
\end{equation*}
$$

and we call any such object, which represents a sequence of partial sums, a series. Loosely speaking, a series is a sum of infinitely many terms. We call $\left(s_{n}\right)$ the sequence of partial sums associated with the series $\sum a_{k}$.

Let $s \in \mathbb{R}$. We say that the series converges to $s$ and we write

$$
\begin{equation*}
\sum_{k=1}^{\infty} a_{k}=s \tag{10.90}
\end{equation*}
$$

if this is true for the associated sequence of partial sums (10.88). We then also say that the series has limit $s$.

We say that the series has limit $\pm \infty$ if this true for the associated partial sums. In this case we write

$$
\begin{equation*}
\sum_{k=1}^{\infty} a_{k}= \pm \infty \tag{10.91}
\end{equation*}
$$

Proposition 10.28 (Convergence criteria for series). A series $s:=\sum a_{k}$ of real numbers converges if and only if for all $\varepsilon>0$ there exists $n_{o} \in \mathbb{N}$ such that one of the following is true:

$$
\begin{array}{ll}
\left|\sum_{k=n}^{\infty} a_{k}\right|<\varepsilon & \text { for all } n \geqq n_{0} \\
\left|\sum_{k=n}^{m} a_{k}\right|<\varepsilon & \text { for all } m, n \geqq n_{0} \tag{10.92b}
\end{array}
$$

PROOF: Write

$$
\begin{equation*}
s=\sum_{k=1}^{\infty} a_{k}=\sum_{k=1}^{n} a_{k}+\sum_{k=n+1}^{\infty} a_{k}=s_{n}+\sum_{k=n+1}^{\infty} a_{k} \tag{10.93}
\end{equation*}
$$

Remember the convergence criteria for numeric sequences. Convergence of a sequence ( $s_{n}$ ) to a real number $s$ means that, for any $\varepsilon>0$, all but finitely many members $s_{n}$ will be inside the $\varepsilon$-neighborhood $N_{\varepsilon}(s)$ of $s$. Expressed in terms of the distance to sthis means there exists a suitable $n_{0} \in \mathbb{N}$ such that

$$
\left|s-s_{n}\right|<\varepsilon \quad \text { for all } n \geqq n_{0}
$$

(see (10.11) on p.192). According to (10.93) we can write that as

$$
\left|\sum_{k=n+1}^{\infty} a_{k}\right|<\varepsilon \quad \text { for all } n \geqq n_{0},
$$

which is the same as (10.92.a) because it does not matter whether we look at the sum of all terms bigger than $n$ or $n+1$.

Alternatively, there was the Cauchy criterion

$$
\left|s_{i}-s_{j}\right|<\delta \text { for all } i, j \geqq n_{0}
$$

(see (10.28) on p.205) which ensures convergence to some number s without specifying what it might actually be. Again we use (10.93) and obtain, assuming without loss of generality that $i<j$,

$$
\left|\sum_{k=i+1}^{j} a_{k}\right|<\delta \text { for all } j>i \geqq n_{0}
$$

Corollary 10.1. If a series $\sum a_{j}$ converges then $\lim _{n \rightarrow \infty} a_{n}=0$.
PROOF: Let $\varepsilon>0$. It follows from $10.92 b$ that there is some $n_{0} \in \mathbb{N}$ such that $\left|a_{m}-0\right|=\left|\sum_{k=m}^{m} a_{k}\right|<\varepsilon$ for all $m \geqq n_{0}$. But this means that the sequence $a_{n}$ converges to zero.

Here is a second corollary.
Corollary 10.2 (Dominance criterion ${ }^{113}$ ). Let $N \in \mathbb{N}$ and let $\sum a_{j}$ and $\sum b_{j}$ be two series such that $\left|b_{k}\right| \leqq a_{k}$ for all $k \geqq N$. It follows that if $\sum a_{k}$ converges then $\sum b_{k}$ converges.

In particular, if $\left|b_{k}\right| \leqq a_{k}$ for all $k \in \mathbb{N}$ then $\left|\sum_{k=1}^{\infty} b_{j}\right| \leqq \sum_{k=1}^{\infty} a_{j}$
PROOF: Let $\varepsilon>0$. It follows from $10.92 b$ that there is some $n_{0} \in \mathbb{N}$ such that $\left|\sum_{k=m}^{n} a_{k}\right|<\varepsilon$ for all $m, n \geqq n_{0}$. Let $M:=\max \left(n_{0}, N\right)$. We obtain

$$
\left|\sum_{k=i+1}^{j} b_{j}\right| \leqq \sum_{k=i+1}^{j}\left|b_{j}\right| \leqq \sum_{k=i+1}^{j} a_{j}<\varepsilon \quad \text { for all } j>i \geqq M .
$$

We conclude from (10.92b) that $\sum b_{k}$ converges.

[^83]Now assume that $\left|b_{k}\right| \leqq a_{k}$ for all $k \in \mathbb{N}$. Let

$$
s_{n}:=\sum_{k=1}^{n}\left|a_{k}\right|, \quad s:=\lim _{n \rightarrow \infty} s_{n}, \quad t_{n}:=\sum_{k=1}^{n} b_{k}, \quad t:=\lim _{n} t_{n} .
$$

It follows from the triangle inequality for real numbers that $\left|t_{n}\right| \leqq s_{n}$ for all $n \in \mathbb{N}$. We apply cor.8.4 on $p .148$ to deduce that $|t|=\lim _{n}\left|t_{n}\right| \leqq \lim _{n} s_{n}=s$. This completes the proof.

Remark 10.24. It is very important to remember that a series either converges to a finite number or it diverges. If it diverges it may be the case that $\sum_{k=1}^{\infty} a_{k}=\infty$ or $\sum_{k=1}^{\infty} a_{k}=-\infty$ or there is no limit at all. As an example for a series which has no limit, look at the oscillating sequence

$$
\begin{equation*}
a_{0}=1 ; \quad a_{1}=-1 ; \quad a_{2}=1 ; \quad a_{3}=-1 ; \cdots \quad s_{n}=\sum_{k=0}^{n}(-1)^{n} \tag{10.94}
\end{equation*}
$$

The above is an example of a series that starts with an index other than 1 (zero). $s_{n}$ obviously does not have limit $+\infty$ or $-\infty$ because $s_{n}$ is 1 for all even $n$ and 0 for all odd $n$. Do not make the mistake of thinking that the limit of the series is zero because you fail to notice the odd indices and only see that $s_{0}=s_{2}=s_{4}=\cdots=s_{2 j}=0$.

Note that for any $j \in \mathbb{N}$ we have $\left|s_{j}-s_{j-1}\right|=1$ because at each step we either add or subtract 1. This means that no matter what real number $a$ and how big a number $n_{0} \in \mathbb{N}$ we choose, it will never be true that $\left|a-s_{j}\right|<1$ for all $j \in \mathbb{N}$ and $a$ cannot be a limit of the series. ${ }^{114}$

Just so you understand the difference between limits and contact points (see (def.10.25) on p.201): Even though neither $\left(a_{j}\right)_{j}$ nor $\left(s_{j}\right)_{j}$ has a limit, the tail sets for both have two contact points each. The ones for $\left(a_{j}\right)_{j}$ have the contact points $\{1,-1\}$ and the ones for $\left(s_{j}\right)_{j}$ have the contact points $\{0,1\}$.

We now turn our attention to convergence properties of series.
Definition 10.40 (Finite permutations). Let $N \in \mathbb{N}$ and let $[N]:=\{1,2,3, \ldots, N\}$ denote the set of the first $N$ integers. ${ }^{115}$ A permutation of $[N]$ is a mapping

$$
\pi(\cdot):[N] \rightarrow[N] ; \quad j \mapsto \pi(j)
$$

which is both surjective: each element $k$ of $[N]$ is the image $\pi(j)$ for a suitable $j \in[N]$ and injective: different arguments $i \neq j \in[N]$ will always map to different images $\pi(i) \neq \pi(j) \in[N]$ (see (4.9) on p.79). Remember that

$$
\text { surjective }+ \text { injective }=\text { bijective }
$$

and that under our assumptions the inverse mapping

$$
\pi^{-1}(\cdot):[N] \rightarrow[N] ; \quad \pi(j) \mapsto \pi^{-1} \pi(j)=j,
$$

[^84]which associates with each image $\pi(j)$ the unique argument $j$ which maps into $\pi(j)$, exists (see def. 4.9 on p. 79 for properties of the inverse mapping).

It is customary to write
$i_{1}$ instead of $\pi(1), \quad i_{2}$ instead of $\pi(2), \quad \ldots, \quad i_{j}$ instead of $\pi(j), \quad \ldots$
Definition 10.41 (Permutations of $\mathbb{N}$ ). A permutation of $\mathbb{N}$ is a bijective function

$$
\pi(\cdot): \mathbb{N} \rightarrow \mathbb{N} ; \quad j \mapsto \pi(j) .
$$

Permutations are the means of describing a A rearrangement or reordering of the members of a finite or infinite sequence or series. Look at any sequence $\left(a_{j}\right)$. Given a permutation $\pi(\cdot)$ of the natural numbers, we can form the sequence $\left(b_{k}\right):=\left(a_{\pi(k)}\right)$, i.e.,

$$
b_{1}=a_{\pi(1)}, \quad b_{2}=a_{\pi(2)}, \quad \ldots, \quad b_{k}=a_{\pi(k)}, \quad \ldots
$$

We can use the inverse permutation, $\pi^{-1}(\cdot)$, to regain the $a_{j}$ from the $b_{j}$ because

$$
b_{\pi^{-1}(k)}=a_{\pi^{-1}(\pi(k))}=a_{k} .
$$

Proposition 10.29. Let $\left(a_{n}\right)$ be a sequence of non-negative members: $a_{n} \geqq 0$ for all $n \in \mathbb{N}$. Then exactly one of the following is true:
a. the series $\sum a_{n}$ converges (to a finite number). In that case

$$
\sum_{n=1}^{\infty} a_{n}=\sum_{n=1}^{\infty} a_{\pi(n)} \quad \text { for any permutation } \pi(\cdot) \text { of } \mathbb{N} .
$$

b. the series $\sum_{n=1}^{\infty} a_{n}$ has limit $\infty$. In that case it is true for any permutation $\pi(\cdot)$ of $\mathbb{N}$ that the reordered series $\sum_{n=1}^{\infty} a_{\pi(n)}$ also has limit $\infty$.

PROOF of $\boldsymbol{a}$ : Let $b_{j}:=a_{\pi(j)}$ and, hence, $a_{k}=b_{\pi^{-1}(j)}$. Let $N \in \mathbb{N}$. Let

$$
\begin{equation*}
\alpha:=\max \{\pi(j): j \leqq N\} \quad \text { and } \quad \beta:=\max \left\{\pi^{-1}(k): k \leqq N\right\} . \tag{10.95}
\end{equation*}
$$

Note that $\alpha \geqq N$ and $\beta \geqq N$. Because all terms $a_{j}, b_{k}$ are non-negative it follows that

$$
\begin{aligned}
& \sum_{j=1}^{N} b_{j}=\sum_{j=1}^{N} a_{\pi(j)} \leqq \sum_{k=1}^{\alpha} a_{k} \leqq \sum_{k=1}^{\alpha} a_{k}+\sum_{k=\alpha+1}^{\infty} a_{k}=\sum_{k=1}^{\infty} a_{k}, \\
& \sum_{k=1}^{N} a_{k}=\sum_{k=1}^{N} b_{\pi^{-1}(k)} \leqq \sum_{j=1}^{\beta} b_{j} \leqq \sum_{j=1}^{\beta} b_{j}+\sum_{j=\beta+1}^{\infty} b_{j}=\sum_{j=1}^{\infty} b_{j} .
\end{aligned}
$$

We take limits as $N \rightarrow \infty$ and it follows from cor.8.4 on $p .148$ that

$$
\sum_{j=1}^{\infty} b_{j} \leq \sum_{k=1}^{\infty} a_{k} \text { and } \sum_{k=1}^{\infty} a_{k} \leqq \sum_{j=1}^{\infty} b_{j}, \quad \text { hence } \sum_{k=1}^{\infty} a_{k}=\sum_{j=1}^{\infty} b_{j} .
$$

This proves part a of the proposition.
PROOF of $\boldsymbol{b}$ : Assume that $\sum a_{j}$ diverges. Because all terms $a_{j}$ are non-negative, the sequence $s_{n}$ of the partial sums is non-decreasing and hence has a limit s. $s \notin \mathbb{R}$ because we assumed that $\sum a_{j}$ is not convergent and we can rule out $s=-\infty$ because $s \geqq a_{0} \geqq 0$. It follows that $s=\infty$.

Assume to the contrary that there is a rearrangement $\sum b_{j}:=\sum a_{\pi(j)}$ of $\sum a_{j}$ which converges to a limit $t \in \mathbb{R}$. According to the already proved part a the rearrangement $\sum a_{j}=\sum b_{\pi^{-1}(j)}$ converges to the same (finite) limit $t$. We have reached a contradiction.

Definition 10.42 (absolutely convergent series). A series $\sum a_{j}$ is absolutely convergent if the corresponding series $\sum\left|a_{j}\right|$ of its absolute values converges.

Proposition 10.30. Let $\sum a_{k}$ be an absolutely convergent series. Then $\sum a_{k}$ converges and

$$
\begin{equation*}
\left|\sum_{k=1}^{\infty} a_{k}\right| \leqq \sum_{k=1}^{\infty}\left|a_{k}\right| \tag{10.96}
\end{equation*}
$$

PROOF: This follows from the dominance criterion (cor.10.2)
It follows from prop. 10.29 on p. 231 that if a series of non-negative terms converges then its value is invariant under rearrangements of that series. The next theorem states that any absolutely convergent series also has that property and we will see later ${ }^{116}$ that the reverse is also true: Any series whose value is invariant under rearrangements is absolutely convergent.

Theorem 10.17. ${ }^{117}$ Let $\sum a_{k}$ be an absolutely convergent series. Let $\pi: \mathbb{N} \rightarrow \mathbb{N}$ be a permutation of $\mathbb{N}$, i.e., the series $\sum b_{k}$ with $b_{k}:=a_{\pi(k)}$ is a rearrangement of the series $\sum a_{k}$. Then $\sum b_{k}$ converges and has the same limit as $\sum a_{k}$. ${ }^{118}$

PROOF: Let $\varepsilon>0$. Since $\sum\left|a_{k}\right|$ converges, there exists $n_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\sum_{k=n_{0}+1}^{n_{0}+m}\left|a_{k}\right| \leqq \sum_{k=n_{0}+1}^{\infty}\left|a_{k}\right|<\varepsilon \quad \text { for all } m \in \mathbb{N} . \tag{10.97}
\end{equation*}
$$

For $n \in \mathbb{N}$ let $s_{n}:=\sum_{k=1}^{n} a_{k}$ and $t_{n}:=\sum_{k=1}^{n} b_{k}$.
Let $A:=\left\{\pi(j): 1 \leqq j \leqq n_{0}\right\}$ and $p_{0}:=\max (A)$. This maximum exists because the set $A$ is finite.
Then $p_{0} \geqq n_{0}$. Each of $a_{1}, a_{2}, \ldots, a_{n_{0}}$ is a term of $s_{n_{0}}$, hence of $s_{p_{0}}$.
Moreover each of $b_{1}=a_{\pi(1)}, b_{2}=a_{\pi(2)}, \ldots, b_{p_{0}}=a_{\pi\left(p_{0}\right)}$ is a term of $t_{p_{0}}$.
Let $n, p \geqq p_{0}$. Then each of $a_{1}, a_{2}, \ldots, a_{n_{0}}$ is a term of $s_{n}$
and each of $b_{1}=a_{\pi(1)}, b_{2}=a_{\pi(2)}, \ldots, b_{p_{0}}=a_{\pi\left(p_{0}\right)}$ is a term of $t_{p}$.
We recall that $p_{0}$ was chosen so big that each of $a_{1}, \ldots, a_{n_{0}}$ is one of $b_{1}, \ldots, b_{p_{0}}$.
It follows from all this that each of $a_{1}, \ldots, a_{n_{0}}$ is a term both of $s_{n}$ and $t_{p}$, hence none of those terms appears in the difference $s_{n}-t_{p}$. We obtain for big enough $m \in \mathbb{N}$ (the bigger of $\max (\{\pi(j): 1 \leqq j \leqq n\})$

[^85]and $p$ )
$$
\left|s_{n}-t_{p}\right| \leqq \sum_{k=n_{0}+1}^{n_{0}+m}\left|a_{k}\right|<\varepsilon
$$

This implies

$$
\left|s-t_{p}\right| \leqq\left|s-s_{n}\right|+\left|s_{n}-t_{p}\right| \leqq\left|s-s_{n}\right|+\sum_{k=n_{0}+1}^{n_{0}+m}\left|a_{k}\right|<\left|s-s_{n}\right|+\varepsilon .
$$

We had chosen $n \geqq n_{0}$ and it follows from (10.97) that $\left|s-s_{n}\right|<\varepsilon$, hence $\left|s-t_{p}\right|<2 \varepsilon$.
We remember how $p$ was defined and see that it could be any integer $\geqq p_{0}$, a number which depends (via $n_{0}$ ) only on $\varepsilon$.

To summarize: for all $\varepsilon>0$ there exists $p_{0}$ such that $p \geqq p_{0}$ implies $\left|s-t_{p}\right|<2 \varepsilon$. But then $\lim _{p \rightarrow \infty} t_{p}=s$. On the other hand, $\lim _{p \rightarrow \infty} t_{p}=t=\sum_{p \rightarrow \infty} b_{k}$.

This concludes the proof that $\sum_{p \rightarrow \infty} a_{k}=\sum_{p \rightarrow \infty} b_{k}$.
There are series which are convergent but not absolutely convergent. Such series are given a special name:
Definition 10.43 (conditionally convergent series). A series $\sum a_{j}$ is called conditionally convergent if it is convergent but not absolutely convergent.

We introduce alternating series to give a simple example of a conditionally convergent series.
Definition 10.44 (alternating series). A series $\sum a_{j}$ is called an alternating series if it is of the form $\sum(-1)^{j} a_{j}$ with either all terms $a_{j}$ being strictly positive or all of them being strictly negative.

Proposition 10.31 (Leibniz test for alternating series). Let $\sum a_{k}$ be an alternating series such that the sequence $\left|a_{k}\right|$ is non-increasing: $\left|a_{1}\right| \geqq\left|a_{2}\right| \geqq\left|a_{3}\right| \geqq \ldots$. Then $\sum a_{k}$ converges.

PROOF: Left as an exercise.
Example 10.14 (Alternating series). The series $\sum(-1)^{n}$ and the alternating harmonic series $\sum(-1)^{n} / n$ are examples of alternating series.

It is known from calculus that the harmonic series $\sum 1 / n$ is divergent: $\sum_{j=1}^{\infty} \frac{1}{n}=\infty$. On the other hand, according to the Leibniz test, $\sum(-1)^{n} / n$ converges. It follows that the alternating harmonic series is convergent but not absolutely convergent, i.e., it is conditionally convergent.

We are going to prove Riemann's Reordering Theorem, from which it can be easily deduced that if $\sum a_{j}$ is conditionally convergent and $x \in \mathbb{R}$, a rearrangement $\sum a_{\pi_{j}}$ can be found which converges to $x$. In preparation we must prove a lemma.
Lemma 10.2. Let $\sum a_{k}$ be a series. We split it into two series $\sum p_{k}$ and $\sum q_{k}$ as follows.
$p_{j}$ is the $j$ th strictly positive member of the sequence $\left(a_{k}\right)_{k}$ and $q_{j}$ is the $j$ th strictly negative member of that sequence.

The following is true:
a. If $\sum a_{k}$ is absolutely convergent then both $\sum p_{k}$ and $\sum q_{k}$ are (absolutely) convergent.
b. If $\sum a_{k}$ is convergent but not absolutely convergent then $\sum p_{k}$ has limit $\infty$ and $\sum q_{k}$ has limit $-\infty$.

PROOF of $\boldsymbol{a}$ : Let $\alpha:=\sum_{i=1}^{\infty}\left|a_{i}\right|$ and let $j \in \mathbb{N}$.
Let $m$ be the index such that $a_{m}$ is the $j$ th (not mth!) strictly positive member of the sequence $\left(a_{k}\right)_{k}$. Then each $p_{i}$ for $i \leqq j$ is some $\left|a_{k}\right|$ for a suitable $k \leqq m$. It follows from $m \geqq j$ that

$$
\sum_{i=1}^{j} p_{i} \leqq \sum_{i=1}^{m}\left|a_{i}\right| \leqq \sum_{i=1}^{\infty}\left|a_{i}\right|<\infty
$$

The above is true for all $j \in \mathbb{N}$ and it follows that $\sum_{i=1}^{\infty} p_{i}<\infty$. The proof that $\sum q_{k}$ has a finite limit is similar.
PROOF of $\boldsymbol{b}$ : The proof will be done in three parts. In part 1 we will show that not both $\sum p_{k}$ and $\sum q_{k}$ can converge. In part 2 we will show that $\sum p_{k}=\infty$ and $\sum q_{k} \in \mathbb{R}$ leads to a contradiction. In part 3 we will show that $\sum q_{k}=-\infty$ and $\sum p_{k} \in \mathbb{R}$ leads to a contradiction.

Part 1: Let us assume that $\sum p_{k}<\infty$ and $\sum q_{k}>-\infty$.
For any $n \in \mathbb{N}$ we have

$$
\sum_{k=1}^{n}\left|a_{k}\right| \leqq \sum_{k=1}^{n} p_{k}+\sum_{k=1}^{n}\left(-q_{k}\right) .
$$

This is true because each one of $a_{1}, \ldots, a_{n}$ is one of the first $n$ strictly positive numbers $p_{1}, \ldots, p_{n}$ or one of the strictly positive numbers $-q_{1}, \ldots,-q_{n}$ or it is zero, in which case it contributes nothing to the series. Both series $\sum a_{k}$ and $\sum\left(-q_{k}\right)$ are non-decreasing, hence for each fixed $n$,

$$
\sum_{k=1}^{n}\left|a_{k}\right| \leqq \sum_{k=1}^{\infty} p_{k}-\sum_{k=1}^{\infty} q_{k} .
$$

It follows that if both $\sum p_{k}$ and $\sum q_{k}$ are convergent then so is $\sum\left|a_{k}\right|$, i.e., this series is absolutely convergent. We have a contradiction.

Part 2: Let us assume that $\sum p_{k}=\infty$ and $\sum q_{k} \in \mathbb{R}$.
We fix $n \in \mathbb{N}$. Let $M_{n}$ be the index of $p_{n}$, i.e., $M_{n}$ is the smallest index $j$ such that $a_{j}=p_{n}$. Note that

$$
a_{M_{n}}=p_{n} \quad(\star) \quad \text { and } \quad M_{n} \geqq n . \quad(\star \star)
$$

Let

$$
I_{n}:=\left\{i \leqq M_{n}: a_{i}>0\right\}, \quad J_{n}:=\left\{j \leqq M_{n}: a_{i}<0\right\} .
$$

Then

$$
\begin{equation*}
\sum_{k=1}^{M_{n}} a_{k}=\sum_{i \in I_{n}} a_{k}+\sum_{j \in J_{n}} a_{k} \stackrel{(\star)}{=} \sum_{i=1}^{n} p_{i}+\sum_{j \in J_{n}} a_{k} \geqq \sum_{i=1}^{n} p_{i}+\sum q_{k} . \tag{10.98}
\end{equation*}
$$

It follows from ( $* \star$ ) that if $n \rightarrow \infty$ then $M_{n} \rightarrow \infty$. Hence, from (10.98),

$$
\sum a_{k}=\lim _{n \rightarrow \infty} \sum_{k=1}^{n} a_{k}=\lim _{n \rightarrow \infty} \sum_{k=1}^{M_{n}} a_{k} \geqq \lim _{n \rightarrow \infty}\left(\sum_{i=1}^{n} p_{i}+\sum q_{k}\right)=\infty
$$

contrary to the assumption that $\sum a_{k}$ converges. We have reached a contradiction.
Part 3: Let us assume that $\sum q_{k}=\infty$ and $\sum p_{k} \in \mathbb{R}$.
We obtain a contradiction by applying part 2 to the series $\sum\left(-a_{k}\right)$.
Theorem 10.18 (Riemann's Rearrangement Theorem). ${ }^{119}$
Let $\alpha, \beta \in \mathbb{R}$ such that $\alpha \leqq \beta$. and let the series $\sum a_{k}$ be conditionally convergent. Then a rearrangement $\sum b_{k}$ of $\sum a_{k}$ exists such that

$$
\liminf _{n \rightarrow \infty} \sum_{k=1}^{n} b_{k}=\alpha \quad \text { and } \quad \limsup _{n \rightarrow \infty} \sum_{k=1}^{n} b_{k}=\beta
$$

PROOF: We may assume that $a_{j} \neq 0$ for all $j \in \mathbb{N}$ because those terms do not contribute anything to the partial sums, hence leave the limit of the series and any rearrangement unchanged.

We split $\sum a_{j}$ into the series $\sum p_{j}$ of its positive members and $\sum q_{j}$ of its negative members in the same way as was done in lemma 10.2:
$p_{j}$ is the $j$ th strictly positive member of the sequence $\left(a_{k}\right)_{k}$;
$q_{j}$ is the $j$ th strictly negative member of $\left(a_{k}\right)_{k}$.
It was proved in lemma 10.2 that $\sum_{k=1}^{\infty} p_{k}=\infty$ and $\sum_{k=1}^{\infty} q_{k}=-\infty$.
case 1: $\beta \geqq 0$.
Let $U_{1}:=\left\{k \in \mathbb{N}: p_{1}+p_{2}+\cdots+p_{k}>\beta\right\}$. $U_{1}$ is not empty because $\sum p_{j}$ has limit $\infty$, hence $u_{1}:=\min \left(U_{1}\right)$ exists. We call the list $p_{1}, p_{2}, \ldots, p_{u_{1}}$ the first upcrossing of the (unfinished) series $\sum b_{k}$.

We now construct the first piece of the desired rearrangement $\sum b_{k}$. Let

$$
n_{1}:=u_{1} ; \quad b_{1}:=p_{1}, b_{2}:=p_{2}, \ldots, b_{n_{1}}:=p_{u_{1}} ; \quad \sigma_{1}:=\sum_{j=1}^{n_{1}} b_{j} .
$$

Note that $n_{1}$ is the first (and so far, only) index $n$ of the series $\sum b_{k}$ for which $\sum_{k=1}^{n} b_{k}$ exceeds $\beta$.
Let $L_{1}:=\left\{k \in \mathbb{N}: \sigma_{1}+\sum_{j=1}^{k} q_{j}<\alpha\right\}$. $L_{1}$ is not empty because $\sum q_{j}$ has limit $-\infty$, hence $l_{1}:=\min \left(L_{1}\right)$ exists. We call the list $q_{1}, q_{2}, \ldots, q_{l_{1}}$ the first downcrossing of $\sum b_{k}$.

We add more terms to $b_{1}, b_{2}, \ldots, b_{n_{1}}$.

$$
n_{2}:=n_{1}+l_{1} ; \quad b_{n_{1}+1}:=q_{1}, b_{n_{1}+2}:=q_{2}, \ldots, b_{n_{2}}:=q_{l_{1}} ; \quad \sigma_{2}:=\sum_{j=1}^{n_{2}} b_{j} .
$$

Note that $n_{2}$ is the first index $n$ of $\sum b_{k}$ for which $\sum_{k=1}^{n} b_{k}$ drops below $\alpha$.

[^86]Let $U_{2}:=\left\{k \in \mathbb{N}: k>u_{1}\right.$ and $\left.\sigma_{2}+\sum_{j=u_{1}+1}^{u_{1}+k} p_{j}>\beta\right\} . U_{2}$ is not empty because $\sum_{j=u_{1}+1}^{\infty} p_{j}$ has limit $\infty$, hence $u_{2}:=\min \left(U_{2}\right)$ exists. We call $p_{u_{1}+1}, p_{u_{1}+2}, \cdots, p_{u_{2}}$ the second upcrossing of $\sum b_{k}$.

We add more terms to $b_{1}, b_{2}, \ldots, b_{n_{2}}$.

$$
n_{3}:=n_{2}+u_{2} ; \quad b_{n_{2}+1}:=p_{u_{1}+1}, b_{n_{2}+2}:=p_{u_{1}+2}, \ldots, b_{n_{3}}:=p_{u_{1}+u_{2}} ; \quad \sigma_{3}:=\sum_{j=1}^{n_{3}} b_{j} .
$$

Note that $n_{3}$ is the second index $n$ of the series $\sum b_{k}$ for which $\sum_{k=1}^{n} b_{k}$ exceeds $\beta$.
Let $L_{2}:=\left\{k \in \mathbb{N}: k>l_{1}\right.$ and $\left.\sigma_{3}+\sum_{j=l_{1}+1}^{l_{1}+k} q_{j}<\alpha\right\} . L_{2}$ is not empty because $\sum_{j=l_{1}+1}^{\infty} q_{j}$ has limit $-\infty$, hence $l_{2}:=\min \left(L_{2}\right)$ exists. We call $q_{l_{1}+1}, q_{l_{1}+2}, \cdots, q_{l_{2}}$ the second downcrossing of $\sum b_{k}$.

We add more terms to $b_{1}, b_{2}, \ldots, b_{n_{3}}$.

$$
n_{4}:=n_{3}+l_{2} ; \quad b_{n_{3}+1}:=q_{l_{1}+1}, b_{n_{3}+2}:=q_{l_{1}+2}, \ldots, b_{n_{4}}:=q_{l_{1}+l_{2}} ; \quad \sigma_{4}:=\sum_{j=1}^{n_{4}} b_{j} .
$$

Note that $n_{4}$ is the second index $n$ of the series $\sum b_{k}$ for which $\sum_{k=1}^{n} b_{k}$ drops below $\alpha$.
It should be clear how we proceed. Let us assume that we have constructed the Nth upcrossing $p_{u_{N-1}+1}, p_{u_{N-1}+2}, \cdots, p_{u_{N}}$ and from it

$$
\begin{aligned}
n_{(2 N-1)} & :=n_{(2 N-2)}+u_{N} ; \\
b_{\left(n_{(2 N-2)}+1\right)} & :=p_{\left(u_{(N-1)}+1\right)}, b_{\left(n_{(2 N-2)}+2\right)}:=p_{\left(u_{(N-1)}+2\right)}, \ldots, b_{\left(n_{(2 N-1)}\right)}:=p_{u_{N}}, \\
\sigma_{(2 N-1)} & :=\sum_{j=1}^{n_{(2 N-1)}} b_{j} .
\end{aligned}
$$

Let us further assume that we have constructed the Nth downcrossing $q_{l_{N-1}+1}, q_{l_{N-1}+2}, \cdots, q_{l_{N}}$ and from it

$$
\begin{aligned}
n_{(2 N)} & :=n_{(2 N-1)}+l_{N} ; \\
b_{\left(n_{(2 N-1)}+1\right)} & :=q_{\left(l_{(N-1)}+1\right)}, b_{\left(n_{(2 N-1)}+2\right)}:=q_{\left(l_{(N-1)}+2\right)}, \ldots, b_{n_{(2 N)}}:=q_{l_{N}}, \\
\sigma_{2 N} & :=\sum_{j=1}^{n_{(2 N)}} b_{j} .
\end{aligned}
$$

We proceed to construct the $(N+1)$ th upcrossing and the $(N+1)$ th downcrossing as follows.

$$
\text { Let } U_{N+1}:=\left\{k \in \mathbb{N}: k>u_{N} \text { and } \sigma_{2 N}+\sum_{j=u_{N}+1}^{u_{N}+k} p_{j}>\beta\right\} . U_{N+1} \text { is not empty because } \sum_{j=u_{N}+1}^{\infty} p_{j} \text { has }
$$

limit $\infty$, hence $u_{N+1}:=\min \left(U_{N+1}\right)$ exists. We call $p_{\left(u_{N}+1\right)}, p_{\left(u_{N}+2\right)}, \cdots, p_{u_{(N+1)}}$ the $(N+1)$ th upcrossing of $\sum b_{k}$.

We add more terms to $b_{1}, b_{2}, \ldots, b_{n_{2 N}}$.

$$
\begin{aligned}
n_{(2 N+1)} & :=n_{(2 N)}+u_{(N+1)} ; \\
b_{\left(n_{(2 N)}+1\right)} & :=p_{\left(u_{N}+1\right)}, b_{\left(n_{(2 N)}+2\right)}:=p_{\left(u_{N}+2\right)}, \ldots, b_{\left(n_{(2 N+1)}\right)}:=p_{u_{(N+1)}}, \\
\sigma_{(2 N+1)} & :=\sum_{j=1}^{n_{(2 N+1)}} b_{j} .
\end{aligned}
$$

Let $L_{N+1}:=\left\{k \in \mathbb{N}: k>l_{N}\right.$ and $\left.\sigma_{2 N+1}+\sum_{j=l_{N}+1}^{l_{N}+k} q_{j}<\alpha\right\} . L_{N+1}$ is not empty because $\sum_{j=l_{N}+1}^{\infty} q_{j}$ has limit $\infty$, hence $l_{N+1}:=\min \left(L_{N+1}\right)$ exists. We call $q_{\left(l_{N}+1\right)}, q_{\left(l_{N}+2\right)}, \cdots, q_{l_{(N+1)}}$ the $(N+1)$ th downcrossing of $\sum b_{k}$.

We add more terms to $b_{1}, b_{2}, \ldots, b_{n_{(2 N+1)}}$.

$$
\begin{aligned}
n_{(2(N+1))} & :=n_{(2 N+1)}+l_{(N+1)} \\
b_{\left(n_{(2 N+1)}+1\right)} & :=q_{\left(l_{N}+1\right)}, b_{\left(n_{(2 N+1)}+2\right)}:=q_{\left(l_{N}+2\right)}, \ldots, b_{\left(n_{2(N+1)}\right)}:=q_{l_{(N+1)}} \\
\sigma_{2(N+1)} & :=\sum_{j=1}^{n_{2(N+1)}} b_{j}
\end{aligned}
$$

We have defined by recursion $\sum_{k=1}^{n_{N}} b_{k}$ for all $N \in \mathbb{N}$
We now show that the increasing sequence $\left(n_{N}\right)_{N \in \mathbb{N}}$ is not bounded above. We observe that $n_{(2 N)}$ is the number of terms that belong to the first $N$ upcrossings plus the first $N$ downcrossings. Each upcrossing and each downcrossing must have at least one term because at least one term $p_{j}$ is needed to move a partial sum from below $\alpha$ to above $\beta$ and at least one term $q_{j}$ is needed to move a partial sum from above $\beta$ to below $\alpha$. Hence $n_{2 N} \geqq 2 N$ and this proves that the sequence $\left(n_{N}\right)_{N \in \mathbb{N}}$ is indeed not bounded above.

It follows that $\sum b_{k}$ indeed has infinitely many terms.
We note that all positive terms $p_{j}$ and all negative terms $q_{j}$ are being used in sequence, starting with the first one. This shows that each one of the terms of $\sum a_{k}$ has become part of $\sum b_{k}$ and it follows that $\sum b_{k}$ is indeed a rearrangement of $\sum a_{k}$.

Let $s_{n}:=\sum_{j=1}^{n} b_{j} . n_{1}, n_{3}, n_{5}, \ldots$ are (precisely the) integers $n$ for which $s_{n}>\beta$ and $n_{2}, n_{4}, n_{6}, \ldots$ are (precisely the) integers $n$ for which $s_{n}<\alpha$. There are infinitely many of each and it follows from thm.8.2 (Characterization of limsup and liminf) on p. 145 that

$$
\begin{equation*}
\liminf _{n \rightarrow \inf } s_{n} \leqq \alpha \quad \text { and } \quad \limsup _{n \rightarrow \inf } s_{n} \geqq \beta \tag{10.99}
\end{equation*}
$$

We now prove that for any $\varepsilon>0$

$$
\begin{equation*}
\liminf _{n \rightarrow \inf } s_{n} \geqq \alpha-\varepsilon \quad \text { and } \quad \limsup _{n \rightarrow \inf } s_{n} \leqq \beta+\varepsilon . \tag{10.100}
\end{equation*}
$$

Let $\varepsilon>0$. The terms $\left(a_{n}\right)_{n}$ of the original series $\sum a_{k}$ converge to zero because $\sum a_{k}$ converges (see cor.10.1 on $p$.229). It follows that there exists $n_{0} \in \mathbb{N}$ such that $\left|a_{j}\right|<\varepsilon$ for all $j \geqq n_{0}$. We show next that

$$
\begin{equation*}
\left|p_{j}\right|=p_{j}<\varepsilon \text { and }\left|q_{j}\right|=-q_{j}<\varepsilon \text { for all } j \geqq n_{0} . \tag{10.101}
\end{equation*}
$$

$\left|p_{j}\right|=p_{j}<\varepsilon$ is true whenever $j \geqq n_{0}$ because $p_{j}$ is the $j$ th positive member of $\left(a_{n}\right)_{n}$, hence $p_{j}=a_{i}$ for some $i \geqq j \geqq n_{0}$. Likewise, $\left|q_{j}\right|=-q_{j}<\varepsilon$ whenever $j \geqq n_{0}$ because $q_{j}$ is the $j$ th negative member of $\left(a_{n}\right)_{n}$, hence $q_{j}=a_{i}$ for some $i \geqq j \geqq n_{0}$. We have proved (10.101).

We recall that $n_{1}, n_{3}, n_{5}, \ldots$ are precisely the integers $n$ for which $s_{n}>\beta$, so

$$
s_{\left(n_{1}-1\right)} \leqq \beta, s_{\left(n_{3}-1\right)} \leqq \beta, \ldots, s_{\left(n_{(2 j-1)}-1\right)} \leqq \beta, \ldots
$$

But then $s_{\left(n_{(2, j-1)}\right)} \leqq \beta+\varepsilon$ because less than $\varepsilon$ was added to the previous term (which is no bigger than $\beta$ ) for any $j$ so big that the last item in the $j$ th upcrossing is less than $\varepsilon$

It follows from (10.101) that $j$ is certainly big enough if $j \geqq n_{0}$ because each upcrossing has size of at least 1 . This shows that there are at most finitely many indices $n$ such that $s_{n}>\beta+\varepsilon$ and we conclude that $\lim \sup _{n} s_{n} \leqq \beta+\varepsilon$. A similar reasoning allows us to conclude that $\lim _{\inf }^{n}{ }_{n} \geqq \alpha-\varepsilon$.

We have proved (10.100) and this impliess, together with (10.99), that

$$
\begin{equation*}
\liminf _{n \rightarrow \inf } s_{n}=\alpha \quad \text { and } \quad \limsup _{n \rightarrow \text { inf }} s_{n}=\beta \tag{10.102}
\end{equation*}
$$

We have proved the theorem for case 1: $\beta \geqq 0$
case 2: $\beta<0$. We proceed exactly as in case 1. The only difference is that we start with a downcrossing that gets us below $\alpha$ rather than an upcrossing to obtain a rearrangement $\sum c_{k}$ for which a partial sum $\sum_{j=1}^{n} a_{j}$ exceeds $\alpha$ when $n$ is the last term of an upcrossing and it drops below $\beta$ when $n$ is the last term of a downcrossing.

Because $a_{j}$ converges to zero there will again only be finitely many upcrossings and downcrossings with terms that exceed $\varepsilon$. For all others the partial sums cannot exceed $\beta$ or drop below $\alpha$ by more than $\varepsilon$ and we conclude as before that

$$
\begin{equation*}
\liminf _{n \rightarrow \inf } \sum_{k=1}^{n} c_{k}=\alpha \quad \text { and } \limsup _{n \rightarrow \inf } \sum_{k=1}^{n} c_{k}=\beta \tag{10.103}
\end{equation*}
$$

Corollary 10.3. Let the series $\sum a_{k}$ be conditionally convergent and let $\alpha \in \mathbb{R}$. Then a rearrangement $\sum b_{k}$ of $\sum a_{k}$ exists such that

$$
\lim _{n \rightarrow \infty} \sum_{k=1}^{n} b_{k}=\alpha
$$

PROOF: We apply Riemann's Reordering Theorem to the special case $\beta=\alpha$ : There is a rearrangement $\sum b_{j}$ of $\sum a_{j}$ such that

$$
\liminf _{n \rightarrow \infty} \sum_{k=1}^{n} b_{k}=\alpha \quad \text { and } \quad \limsup _{n \rightarrow \infty} \sum_{k=1}^{n} b_{k}=\alpha .
$$

It follows now from thm.8.3 on p. 147 that $\sum b_{j}$ converges to $\alpha$.
We have seen that if a series is absolutely convergent then it is convergent and each rearrangement converges to the same limit. Here is the reverse.
Corollary 10.4. Let $\sum a_{k}$ be a convergent series with limit $\alpha \in \mathbb{R}$ such that each rearrangement $\sum b_{k}$ also converges to $\alpha$.

Then $\sum a_{k}$ is absolutely convergent.

PROOF: We assume to the contrary that $\sum a_{k}$ is not absolutely convergent. This series is convergent but not absolutely convergent, hence conditionally convergent. We apply Riemann's Reordering Theorem and find that there is a rearrangement of $\sum a_{j}$ which converges to a different real number, contrary to our assumption.

Corollary 10.5 (Dichotomy for convergent series). Let series $\sum a_{k}$ be a convergent series. Then either $\boldsymbol{a}$ or $\boldsymbol{b}$ is true:
a. All rearrangements of $\sum a_{k}$ converge to the same limit.
b. For any $\alpha \in \mathbb{R}$ there is a rearrangement of $\sum a_{k}$ which converges to $\alpha$.

PROOF: Either $\sum a_{k}$ is absolutely convergent and $\boldsymbol{a}$ is true according to Riemann's Reordering Theorem or the series it is conditionally convergent and $\boldsymbol{b}$ is true according to cor.10.3.

### 10.4 Exercises for Ch. 10

### 10.4.1 Exercises for Ch.10.1

### 10.4.1.1 Exercises for Ch.10.1.1 (Definition and Examples of Metric Spaces)

Exercise 10.1. Prove prop. 10.1 on p.184: Let $(X, d)$ be a metric space. Let $n \in \mathbb{N}$ and $x_{1}, x_{2}, \ldots, x_{n} \in$ $X$. Then

$$
\begin{equation*}
d\left(x_{1}, x_{n}\right) \leqq \sum_{j=1}^{n-1} d\left(x_{j}, x_{j+1}\right)=d\left(x_{1}, x_{2}\right)+d\left(x_{2}, x_{3}\right)+d\left(x_{n-1}, x_{n}\right) . \tag{10.104}
\end{equation*}
$$

Exercise 10.2. Prove thm. 10.1 (Norms define metric spaces) on p.185: Let $(V,\|\cdot\|)$ be a normed vectors space. Then the function

$$
\begin{equation*}
d_{\|\cdot\|}(\cdot, \cdot): V \times V \rightarrow \mathbb{R}_{\geqq 0} ; \quad(x, y) \mapsto d_{\|\cdot\|}(x, y):=\|y-x\| \tag{10.105}
\end{equation*}
$$

defines a metric space $\left(V, d_{\|\cdot\|}\right)$.
Hint: This proof is very easy. Even the triangle inequality for the metric $d(x, y)=\|x-y\|$ follows easily from the triangle inequality for the norm.

### 10.4.1.2 Exercises for Ch.10.1.2 (Measuring the Distance of Real Valued Functions)

### 10.4.1.3 Exercises for Ch.10.1.3 (Neighborhoods and Open Sets)

Exercise 10.3. Prove prop. 10.32 on p.242: Let $a, b \in \mathbb{R}$ such that $a<b$. Then the open interval $] a, b[$ is an open set in $\left(\mathbb{R}, d_{|\cdot|}\right)$.

Exercise 10.4. Let $A:=\left\{\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}: x_{1}>0, x_{2}>0\right\}$ be the first quadrant in the plane (the points on the coordinate axes are excluded). Prove that each element of $A$ is an inner point, i.e., $A$ is open in $\mathbb{R}^{2}$. See the picture for a hint.

Hint: Find for $\vec{a}=\left(a_{1}, a_{2}\right)$ small enough $\varepsilon$ such that $N_{\varepsilon}(\vec{a}) \subseteq A$


Exercise 10.5. Let $a, b \in \mathbb{R}$ such that $a<b$.
a. The closed interval $[a, b]$ is not open in $\left(\mathbb{R}, d_{|\cdot|}\right)$.
b. The complement of the closed interval $[a, b]$ is open in $\left(\mathbb{R}, d_{|\cdot|}\right)$.

Exercise 10.6. a. Let $m \in \mathbb{Z}$, viewed as a subset of the metric space $\left(\mathbb{R}, d_{|\cdot|}\right)$. Prove that $m$ is a boundary point of $\mathbb{Z}$.
b. Prove that the above also holds both for the set of rational numbers and the set of all irrational numbers.

### 10.4.1.4 Exercises for Ch.10.1.4 (Convergence)

Exercise 10.7. Given is a metric space $(X, d)$.
Prove the following: A sequence $\left(x_{n}\right)$ of elements of $X$ converges to $a \in X$ as $n \rightarrow \infty$ iff for any neighborhood $U$ of $a$ there exists some $n_{0} \in \mathbb{N}$ such that the $n_{0}$-tail set $T_{n_{0}}=\left\{x_{j}: j \geqq n_{0}\right\}$ is contained in $U$ (see def.8.13 (Tail sets of a sequence) on p.140.)
Exercise 10.8. Prove prop. 10.6 on p.193:
Let $x_{n}$, $y_{n}$ be two sequences in a metric space $(X, d)$. Assume there is $K \in \mathbb{N}$ such that $x_{n}=y_{n}$ for all $n \geqq K$. Let $L \in X$ Then

$$
\lim _{n \rightarrow \infty} x_{n}=L \Leftrightarrow \lim _{n \rightarrow \infty} y_{n}=L .
$$

Exercise 10.9. Prove prop. 10.5 on p.192:
Let $x_{n}$ be a convergent sequence in a metric space $(X, d)$ with limit $L \in E$. Let $K \in \mathbb{N}$. For $n \in \mathbb{N}$ let $y_{n}:=x_{n+K}$. Then $\lim _{n \rightarrow \infty}\left(y_{n}\right)_{n}=L$.
Exercise 10.10. Let $f_{n}, f \in \mathscr{B}([0,1], \mathbb{R}) \quad n \in \mathbb{N}$ be continuous such that $f=\lim _{n \rightarrow \infty} f_{n}$ in $\left(\mathscr{B}\left([0,1], d_{\|\cdot\|_{\infty}}\right)(!)\right.$ Prove $\lim _{n \rightarrow \infty} \int_{0}^{1} f_{n}(x) d x=\int_{0}^{1} f(x) d x$. You must use the $\varepsilon, N$ definition of convergence.
Hints: a. No need to mention that continuous functions are both bounded and integrable and that they attain both max and min on closed and bounded intervals. b. Use the mean value theorem: For cont. $h(\cdot)$ on $[0,1]$ let $\alpha:=\min _{x \in[0,1]} h(x), \beta:=\max _{x \in[0,1]} h(x)$. Then $\exists \lambda \in[\alpha, \beta]$ such that $\lim _{n \rightarrow \infty} \int_{0}^{1} h(x) d x=\lambda\left(=\lambda(1-0)\right.$. c. Use without proof that $\left|\int_{0}^{1} h(x) d x\right| \leqq \int_{0}^{1}|h(x)| d x$ for any integrable $h(\cdot)$ on $[0,1] \mathbf{d}$. Apply $\mathbf{b}$ and $\mathbf{c}$ to $h_{n}(x)=\left|f_{n}(x)-f(x)\right|$. (So you deal with $\alpha_{n}, \lambda_{n}, \beta_{n}$ ).

### 10.4.1.5 Exercises for Ch.10.1.5 (Abstract Topological spaces)

Exercise 10.11. It was stated in example. 10.5 on p. 194 that the discrete topology which is induced by the discrete metric $d(x, y)=1$ if $x \neq y$ and 0 if $x=y$ is the entire powerset $2^{X}$ of $X$. Prove it.

### 10.4.1.6 Exercises for Ch.10.1.6 (Neighborhood Bases)

Exercise 10.12. Let $(X, d)$ be a metric space and let $\mathfrak{B}:=\left\{N_{1 / k}(x): x \in X, k \in \mathbb{N}\right\}$. Then $\mathfrak{B}$ is a base of the topology for the associated topological space ( $X, \mathfrak{U}_{d}$ ).

### 10.4.1.7 Exercises for Ch.10.1.7 (Metric and Topological Subspaces)

Exercise 10.13. Prove prop. 10.11 on p.200: Let ( $X, \mathfrak{U}$ ) be a topological space, $A \subseteq X$, and let $\mathfrak{U}_{A}$ be the collection of all open sets in $A$. Then $\left(A, \mathfrak{U}_{A}\right)$ is a topological space, i.e., it satisfies the definition def. 10.12 on p. 193 of an abstract topological space.

### 10.4.1.8 Exercises for Ch.10.1.8 (Bounded Sets and Bounded Functions)

### 10.4.1.9 Exercises for Ch.101.9 (Contact Points and Closed Sets)

Exercise 10.14. Let $A \subseteq \mathbb{R}$ be a closed, nonempty set which is bounded above. Prove that the maximum of $A$ exists and that $\sup (A)=\max (A)$.

Exercise 10.15. Prove prop. 10.17 on p.204: Let $(X, \mathfrak{U})$ be a topological space and $A \subseteq X$. Then $\partial A=\bar{A} \cap \overline{A^{\complement}}$.

Exercise 10.16. Prove prop. 10.15 on p.204: Let ( $X, \mathfrak{U}$ ) be a topological space and $A \subseteq B \subseteq X$. Then $\bar{A} \subseteq \bar{B}$.

Exercise 10.17. Prove rem. 10.10 on p.195: Let ( $X, \mathfrak{U}$ ) be a topological space and $A \subseteq B \subseteq X$. Then $A^{o} \subseteq B^{o}$.

Exercise 10.18. Prove parts $\mathbf{c}$ and $\mathbf{d}$ of prop. 10.19 (Closure of a set as a hull operator) on p.205: Let $A$ be a subset of a topological space $(X, \mathfrak{U})$. Then $\quad$ c. $\bar{A}=\bar{A}, \quad$ d. $\overline{A \cup B}=\bar{A} \cup \bar{B}$.

### 10.4.1.10 Exercises for Ch.10.1.10 (Completeness in Metric Spaces)

Exercise 10.19. Let $(X, d)$ be a metric space and $A \subseteq X, A \neq \emptyset$. Let

$$
\begin{equation*}
\gamma:=\gamma(A):=\inf \{d(x, y): x, y \in A \text { and } x \neq y\} . \tag{10.106}
\end{equation*}
$$

a. Prove that if $\gamma>0$ then $A$ is complete.
b. The reverse is not true. Find a counterexample.

Exercise 10.20. Let $(X, d)$ be a metric space and let $A \subseteq X$ be a finite subset. Prove that $A$ is complete.

Exercise 10.21. Given is $\mathbb{R}$ with the Euclidean metric $d(x, y)=|x-y|$. We look at $\mathbb{N}$ and $\mathbb{Q}$ as metric subspaces of $\mathbb{R}$. We know that $\mathbb{Q}$ is not complete.
a. Is $\mathbb{N}$ complete as a subspace of $\mathbb{Q}$ ?
b. Is $\mathbb{N}$ complete as a subspace of $\mathbb{R}$ ?

Prove your answer.
Exercise 10.22. Let $X$ be a nonempty set with the discrete metric $d(x, y)=1-1_{\{x\}}(y)$, i.e., $d(x, y)=0$ if $x=y$ and 1 else. Prove that $(X, d)$ is complete.

### 10.4.2 Exercises for Ch.10.2

Exercise 10.23. Prove prop. 10.7 (Opposite of continuity) on p.193:
A sequence $\left(x_{k}\right)_{k}$ with values in $(X, d)$ does not have $L \in X$ as its limit if and only if there exists some $\varepsilon>0$ and $n_{1}<n_{2}<n_{3}<\cdots \in \mathbb{N}$ such that $d\left(x_{n_{j}}, L\right) \geq \varepsilon$ for all $j$.

Exercise 10.24. Prove prop.10.25 (All polynomials are continuous) on p. 218.
Exercise 10.25. Let $n \in \mathbb{N}$ and let $p(x):=\sum_{j=0}^{n} a_{j} x^{j}\left(a_{j} \in \mathbb{R}, a_{n} \neq 0\right)$ be a polynomial. Let $A \subseteq \mathbb{R}$ be unbounded. Prove that the direct image $p(A)$ is unbounded. Hint: Show that there is a sequence $x_{n} \in A$ such that $x_{n} \rightarrow \infty$ or $x_{n} \rightarrow-\infty$. Examine $y_{k}:=\frac{p\left(x_{k}\right)}{a_{n} x_{k}^{n}}$ and use prop.8.10 (rules of arithmetic for limits).

Exercise 10.26. In example 10.13 on p. 225 the functions $f_{n}(\cdot)$ were defined as follows on the closed unit interval [ 0,1 ]:

$$
f_{n}(x):= \begin{cases}n^{2} x & \text { for } 0 \leqq x \leqq \frac{1}{n} \\ \frac{1}{x} & \text { for } \frac{1}{n} \leqq x \leqq 1\end{cases}
$$

Prove that $f_{n}$ is continuous for all $n \in \mathbb{N}$.

### 10.5 Addenda to Ch. 10 (Metric spaces)

The following theorem belongs before def. 10.9 (Neighborhoods in Metric Spaces) on p. 190 .

Proposition 10.32 (Open intervals are open in $\left(\mathbb{R}, d_{|\cdot|}\right)$ ). Let $a, b \in \mathbb{R}$ such that $a<b$. Then the open interval $] a, b\left[\right.$ is an open set in $\left(\mathbb{R}, d_{|\cdot|}\right)$.

PROOF:
left as exercise 10.3 on p.239.
Example 10.15 (Cover page of the $B / G$ book). Geometric meaning of

$$
\sum_{j=0}^{n}(1 / 4)^{j}=\frac{4}{3}=1+\sum_{j=1}^{n}(1 / 4)^{j}, \quad \text { i.e., } \sum_{j=1}^{n}(1 / 4)^{j} ;=\frac{1}{3}
$$

a. each subsequent iteration has half the height (similar triangles), hence the triangles of each iteration have $1 / 4$ the area of the previous one. That means that if $(1 / 4)^{j}$ is the area of each of the triangles in iteration $j$ then $(1 / 4) \cdot(1 / 4)^{j}$ is the area of each of the triangles in iteration $j+1$.
b. in each horizontal slice the shaded triangle has $1 / 3$ the area of the entire slice. That means that the limit of the sum of the shaded triangles is $1 / 3$ of the total area.

Proposition 10.33. Let the metric space $\left(\mathbb{Q}, d_{|\cdot|}\right)$ (Euclidean metric) is not complete.
PROOF: Let us work for the time being in the metric space $\left(\mathbb{R}, d_{|\cdot|}\right)$ of all real numbers, not in the subspace $\left(\mathbb{Q}, d_{|\cdot|}\right)$ which we are interested in.

Let $x \in \mathbb{R}$ be any positive irrational number, e.g., $x=\pi=3.1415 \ldots$ or $x=\sqrt{2}=1.414 \ldots$. $x$ has a decimal representation $x=m+\sum_{j=1}^{\infty} d_{j} 10^{-j}$ where each $d_{j} \in \mathbb{Z}_{\geq 0}$ is a digit, i.e., $0 \leqq d_{j} \leqq 9$. Let $s_{n}=m+\sum_{j=1}^{n} d_{j} 10^{-j}$ Then

$$
\begin{equation*}
\left|x-s_{n}\right|=x-s_{n}=\sum_{j=n+1}^{\infty} d_{j} 10^{-j} \leqq 9 \cdot \sum_{j=n+1}^{\infty} 10^{-j}=10^{-n} \tag{10.107}
\end{equation*}
$$

and it follows, not surprisingly, that $x=\lim _{n \rightarrow \infty} x_{n}$.
We know from thm. 10.12 (Convergent sequences are Cauchy) on $p .210$ that the sequence $s_{n}$ is Cauchy in $\left(\mathbb{R}, d_{|\cdot|}\right)$. But $s_{n} \in \mathbb{Q}$ for all $n$ and the distance $d_{|\cdot|}\left(s_{n}, s_{m}\right)=\left|s_{n}-s_{m}\right|$ is the same in $\left(\mathbb{R}, d_{|\cdot|}\right)$ and $\left(\mathbb{Q}, d_{|\cdot|}\right)$.

It follows that $s_{n}$ is Cauchy in $\left(\mathbb{Q}, d_{|\cdot|}\right)$. We had constructed this sequence in such a way that it does not have a limit in $\mathbb{Q}$, and it follows that $\left(\mathbb{Q}, d_{|\cdot|}\right)$ is not complete.

## 11 Compactness

### 11.1 Introduction: Closed and Bounded sets in Euclidean space

One of the results that are true for $N$-dimensional space is the "sequence compactness" of closed and bounded subsets: Any sequence that lives in such a set has a convergent subsequence. We will discuss that next.
Theorem 11.1 (Convergent subsequences in closed and bounded sets of $\mathbb{R}$ ). Let $A$ be a bounded and closed set of real numbers and let $\left(z_{n}\right)$ be an arbitrary sequence in $A$. Then there exists $z \in A$ and a subset

$$
n_{1}<n_{2}<\ldots<n_{j}<\ldots \text { of indices such that } z=\lim _{j \rightarrow \infty} z_{n_{j}}
$$

i.e., the subsequence $\left(z_{n_{j}}\right)$ converges to $z$.

PROOF: Let $m$ be the midpoint between $a:=\inf (A)$ and $b:=\sup (A)$. Because $A$ is bounded, $a$ and $b$ must exist as finite numbers. Let

$$
\begin{equation*}
A_{\star 1}:=A \cap[a, m] ; \quad A^{\star}{ }_{1}:=A \cap[m, b] . \tag{11.1}
\end{equation*}
$$

Then at least one of $A_{\star 1}, A^{\star}{ }_{1}$ must contain infinitely many of the $z_{n}$ because $A_{\star 1}$ and $A^{\star}{ }_{1}$ form a "covering" of $A$ (the formal definition will be given later in def.11.4 on p.252), i.e., $A_{\star 1} \cup A^{\star}{ }_{1} \supseteq A$. We pick one with infinitely many elements and call it $A_{1}$. In case both sets contain infinitely many of the $z_{n}$, it does not matter which one we pick. We observe that $\operatorname{diam}\left(A_{1}\right) \leqq \operatorname{diam}(A) / 2$

Let $m_{1}$ be the midpoint between $a_{1}:=\inf \left(A_{1}\right)$ and $b_{1}:=\sup \left(A_{1}\right)$. Let

$$
\begin{equation*}
A_{\star 2}:=A_{1} \cap\left[a_{1}, m_{1}\right] ; \quad A^{\star} 2:=A \cap\left[m_{1}, b_{1}\right] . \tag{11.2}
\end{equation*}
$$

Then at least one of $A_{\star 2}, A^{\star}{ }_{2}$ must contain infinitely many of the $z_{n}$. We pick one with infinitely many elements and call it $A_{2}$. In case both sets contain infinitely many of the $z_{n}$, it does not matter which one we pick. Note that

$$
\operatorname{diam}\left(A_{2}\right) \leqq \operatorname{diam}\left(A_{1}\right) / 2 \leqq \operatorname{diam}(A) / 2^{2}
$$

We keep picking the midpoints $m_{j}$ of the sets $A_{j}$ each of which has at most half the diameter of the previous one. (Why?) In other words, we have constructed a sequence

$$
\begin{align*}
& A \supseteq A_{1} \supseteq A_{2} \supseteq \ldots \supseteq A_{n} \supseteq \ldots \quad \text { such that } \\
& \operatorname{diam}(A) \geqq 2 \operatorname{diam}\left(A_{1}\right) \geqq 2^{2} \operatorname{diam}\left(A_{2}\right) \ldots \geqq 2^{n} \operatorname{diam}\left(A_{n}\right) \ldots \tag{11.3}
\end{align*}
$$

which means that $\operatorname{diam}\left(A_{n}\right) \leqq \operatorname{diam}(A) / 2^{n} \rightarrow 0$ as $n \rightarrow \infty$.
We pick a subsequence $\left(x_{j}\right)_{j}=\left(z_{n_{j}}\right)_{j}$ of the original sequence $\left(z_{n}\right)_{n}$ such that $z_{n_{j}} \in A_{j}$ for all $j \in \mathbb{N}$. This is possible because the sets $A_{j}$ were picked in such a way that each one of them contains infinitely many of the $z_{k}$.

The sequence of sets $\left(A_{j}\right)$ is "nested": $A_{j} \subseteq A_{j-1}$ for each index $j$. It follows that $A_{m}$ contains all $A_{k}$ for any $k>m$ and this implies that $A_{m}$ contains all members $x_{k}=z_{n_{k}}$, for all $k>m$. Thus

$$
\left|x_{m}-x_{k}\right| \leqq \operatorname{diam}\left(A_{m}\right) \leqq \frac{\operatorname{diam}(A)}{2^{m}} \quad \text { for all } m \text { and } k \text { such that } k>m
$$

It follows that $\left(x_{n}\right)$ is a Cauchy sequence ( $p .205$ ). According to theorem 10.8 about the completeness of $\mathbb{R}$ (p.206) there is a contact point $x$ such that $x_{n} \rightarrow x$ for $n \rightarrow \infty$.

Because $A$ is a closed set it contains all its contact points. It follows that $x \in A$ and we have found $a$ subsequence of the original sequence $\left(z_{n}\right)$ which converges to an element of $A$.

Theorem 11.2 (Convergent subsequences in closed and bounded sets of $\mathbb{R}^{N}$ ). Let $A$ be a bounded and closed set of $\mathbb{R}^{N}$ and let $\left(\vec{z}_{n}\right)$ be an arbitrary sequence of $N$-dimensional vectors in $A$. Then there exists $\vec{z} \in A$ and a subset

$$
n_{1}<n_{2}<\ldots<n_{j}<\ldots \text { of indices such that } \vec{z}=\lim _{j \rightarrow \infty} \vec{z}_{n_{j}}
$$

i.e., the subsequence $\left(\vec{z}_{n_{j}}\right)$ converges to $\vec{z}$.

Proof (outline): We review the above proof for $\mathbb{R}$ :
The base idea was to chop $A$ in half during each step to obtain a sequence of sets $A_{n}$ which become smaller and smaller in diameter but yet contain infinitely many points. of the original sequence $z_{n}$.

In higher dimensions we would still find the center point $\vec{m}_{n}$ which is determined by the fact that it is the center of a $\gamma$-neighborhood ( $N$-dimensional ball) that contains $A_{n}$ and does so with the smallest radius possible. We then take the minimal square (in $\mathbb{R}^{2}$ ) or the minimal $N$-dimensional cube (in $\mathbb{R}^{N}$ ) that is parallel to the coordinate axes and still contains that sphere or ball.

We then divide that $N$-dimensional cube (a square in 2 dimensions, a cube in 3 dimensions) into $2^{N}$ sectors ( 4 quadrants in $\mathbb{R}^{2}, 8$ sectors in $\mathbb{R}^{3}$ ) and partition $A_{n}$ into at most $2^{N}$ pieces by intersecting it with those $2^{N}$ sectors).

The set $A_{n+1}$ is then chosen from one of those pieces which contain infinitely many of the $z_{n}$. Again, we get a nested sequence $A_{n}$ whose diameters contract towards 0 . You'll find more detail about the messy calculations required in the proof of prop.11.2 on p.247.

Each $A_{n}$ contains infinitely many of the $\vec{z}_{k}$. Now pick $\vec{x}_{j}:=\vec{z}_{n_{j}}$ where $\vec{z}_{n_{j}}$ is one of the infinitely many members of the original sequence $\left(\vec{z}_{n}\right)$ which are contained in $A_{j}$.

Because $A_{j} \subseteq A_{K}$ for $j \geqq K$ and $\lim _{K \rightarrow \infty} \operatorname{diam}\left(A_{K}\right)=0$, we do the following for a given $\varepsilon>0$ : choose $K$ so big that $\operatorname{diam}\left(A_{K}\right) \leqq \varepsilon$. Note that

$$
\text { if } i, j \geqq K \text { then } d\left(\vec{x}_{i}, \vec{x}_{j}\right)=d\left(\vec{z}_{n_{i}}, \vec{z}_{n_{j}}\right) \leqq \operatorname{diam}\left(A_{K}\right) \leqq .
$$

This is true because $n_{i} \geqq i$ (and $n_{j} \geqq j$ ), hence $\vec{x}_{i}, \vec{x}_{j} \in A_{K}$. It follows that the sequence ( $\vec{x}_{j}$ ) is Cauchy. We have seen in thm. 10.9 on $p .207$ that $\mathbb{R}^{N}$ is complete. It follows that $\vec{L}:=\lim _{j \rightarrow \infty} \vec{x}_{j}$ exists in $\mathbb{R}^{N}$.

The proof is complete if it can be shown that $\vec{L} \in A$.
But we know that all $\vec{x}_{i}=\vec{z}_{n_{i}}$ belong to $A$. $\vec{L}$ is a contact point of $A$ according to the sequence criterion for contact points in metric spaces (thm.10.6 on p.201).
any neighborhood $N_{\varepsilon}(\vec{L})$ contains an entire tail set of the sequence $\left(\vec{x}_{i}\right)_{i}$. As the closed set $A$ owns all its contact points, it follows that $L \in A$ and the theorem is proved.

Theorem 11.3. Let $A$ be a bounded and closed set of real numbers and let $f(\cdot): A \rightarrow \mathbb{R}$ be a continuous function on $A$. Then $f(\cdot)$ is a bounded function.

PROOF: Let us assume that $f(\cdot)$ is not bounded and show that this leads to a contradiction.
An unbounded function is not bounded from above, from below, or both. We might as well assume that $f(\cdot)$ is not bounded from above because otherwise it is not bounded from below and we can work with $-f(\cdot)$ which then is not bounded from above. We conclude that there is a sequence $\left(z_{n}\right) \in A$ such that

$$
\begin{equation*}
f\left(z_{n}\right)>n \quad \text { for all } n \in \mathbb{N} \tag{11.4}
\end{equation*}
$$

According to thm.11.1 (Convergent subsequences in closed and bounded sets) there exists a subsequence $\left(x_{j}\right)=\left(z_{n_{j}}\right)$ and $x_{0} \in A$ such that $x_{0}=\lim _{n \rightarrow \infty} x_{n}$.

In particular, $f\left(x_{0}\right)$ exists as a finite value and $f\left(x_{n}\right) \rightarrow f\left(x_{0}\right)$ because $f(\cdot)$ is continuous in $x_{0}$. But the $x_{n}$ were constructed as a subsequence of the $z_{j}$ which have the property that $f\left(z_{j}\right)>j$ for all $j$.

The subsequence $\left(f\left(x_{n}\right)\right)$ cannot converge to $f\left(x_{0}\right)$ because $f\left(x_{j}\right)=f\left(z_{n_{j}}\right)>n_{j}$, i.e., $\lim _{j \rightarrow \infty} f\left(z_{n_{j}}\right)=\infty$. We have reached a contradiction and it follows that $f(\cdot)$ is bounded.

Corollary 11.1. Let $a<b$ be two real numbers and let $f(\cdot):[a, b] \rightarrow \mathbb{R}$ be a continuous function on $[a, b]$. Then $f(\cdot)$ is a bounded function.

PROOF: The interval interval $[a, b]$ is closed and bounded (diam $([a, b])=b-a$ ). The proof follows from thm.11.3.

### 11.2 Four Definitions of Compactness

We now look at ways to extend those results to general metric spaces by looking at the concept of compactness.
Compact sets are a wonderful thing to deal with because they allow you in some sense to go from dealing with "arbitrarily many" to dealing with "countably many" and even "finitely many". This chapter will show that $A, B$ and $C$ below are equivalent statements for any subspace $\left(K,\left.d\right|_{K \times K}\right)$ of a metric space $(X, d)$ :
A. any sequence in $K$ has a convergent subsequence
B. $K$ is complete and contains only finitely many point of a grid of length $\varepsilon$
C. any open covering of $K$ has a finite subcovering
D. $K$ is bounded and closed - ONLY works in $\mathbb{R}^{N}$ !

Such metric spaces $K$ will be called "compact" (see def. 11.3 (Sequence compactness) on $p .249$ and def. 11.5 (Compact sets) on p.253).

We can now state theorems 11.1 and 11.2 of the introduction as follows: In $\mathbb{R}^{N}$ statement D implies $A$.
When you take a course on real analysis you will probably be given the definition of compactness as that in C: any open covering of $K$ has a finite subcovering. In this document this definition is pushed into the background as it is the most difficult to understand. But full proofs will be given of the equivalence of sequence compactness (def.A) on the one hand and completeness plus "total boundedness" (def.B) on the other hand.

One of the important results of this chapter on compactness is that, if you look at $\mathbb{R}^{N}$ with the Euclidean norm and its associated metric

$$
d(\vec{x}, \vec{y})=\sqrt{\sum_{j=1}^{n}\left(x_{j}-y_{j}\right)^{2}} \quad\left(\vec{x}=\left(x_{1}, x_{2}, \ldots\right), \vec{y}=\left(y_{1}, y_{2}, \ldots\right) \in \mathbb{R}^{N}\right)
$$

(see (9.3) on $p$.158) then all four statements $A, B, C, D$ coincide.

## $11.3 \varepsilon$-Nets and Total Boundedness

We now briefly discuss $\varepsilon$-nets and decreasing sequences of closed sets which contract to a single point.

Definition 11.1 ( $\varepsilon$-nets). Let $\varepsilon>0$. Let $(X, d)$ be a metric space and $A \subseteq X$. let $G \subseteq A$ be a subset of $A$ with the following property:

$$
\begin{equation*}
\text { For each } x \in A \text { there exists } g \in G \text { such that } x \in N_{\varepsilon}(g) \text {. } \tag{11.5}
\end{equation*}
$$

In other words, the points of $G$ form a "grid" or "net" fine enough so that no matter what point $x$ of $A$ you choose, you can always find a "grid point" $g$ with distance less than $\varepsilon$ to $x$, because that is precisely the meaning of $x \in N_{\varepsilon}(g)$.

We call $G$ an $\varepsilon$-net or $\varepsilon$-grid for $A$ and we call $g \in G$ a grid point of the net.
Proposition 11.1 ( $\varepsilon$-nets and coverings). Let $\varepsilon>0$. Let $(X, d)$ be a metric space and $A \subseteq X$. Let $G \subseteq A$ be an $\varepsilon$-grid for $A$. Then $\left\{N_{\varepsilon}(g)\right\}_{g \in G}$ is an open covering of $A$ in the sense of def.11.4 on p.252: It is a collection of open sets whose union "covers", i.e., contains, $A$.

PROOF: Let $x \in A$. We can choose a point $g=g(x) \in G$ such that $x \in N_{\varepsilon}(g(x))$. It follows from $\{x\} \subseteq N_{\varepsilon}(g(x))$ and $g(x) \in G$ for all $x \in A$ that

$$
A=\bigcup_{x \in A}\{x\} \subseteq \bigcup_{x \in A} N_{\varepsilon}(g(x)) \subseteq \bigcup_{g \in G} N_{\varepsilon}(g) .
$$

Proposition $11.2\left(\varepsilon\right.$-nets in $\left.\mathbb{R}^{N}\right)$. Let $(X, d)$ be $\mathbb{R}^{N}$ with the Euclidean metric.
A. Let

$$
\mathbb{Z}^{N}=\left\{\vec{z}=\left(z_{1}, z_{2}, \ldots z_{N}\right): z \in \mathbb{Z}\right\}
$$

In other words, $\mathbb{Z}^{N}$ is the set of all points in $\mathbb{R}^{N}$ with integer coordinates. ${ }^{120}$
Then $\mathbb{Z}^{N}$ is a $\sqrt{N}$-net of $\mathbb{R}^{N}$.
B. Let $\varepsilon>0$ and $G_{\varepsilon}^{\mathbb{R}^{N}}:=\left\{\varepsilon \vec{z}: \vec{z} \in \mathbb{Z}^{N}\right\}$. Then $G_{\varepsilon}^{\mathbb{R}^{N}}$ is an $\varepsilon \sqrt{N}$-net of $\mathbb{R}^{N}$.
C. Let $A$ be a bounded set in $\mathbb{R}^{N}$ and $\varepsilon>0$. Then there is $n \in \mathbb{N}$ and $g_{1}, \ldots g_{n} \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}$ such that

$$
A \subseteq N_{\varepsilon}\left(g_{1}\right) \cup N_{\varepsilon}\left(g_{2}\right) \cup \ldots, \cup N_{\varepsilon}\left(g_{n}\right)
$$

i.e., $A$ is covered by finitely many $\varepsilon$-neighborhoods of points in the $(\varepsilon / \sqrt{N})$-grid $G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}$.
(Skip this proof!) (all three parts A, B, C)

## PROOF of $A$.

Let $\vec{x}=\left(x_{1}, x_{2}, \ldots x_{N}\right) \in \mathbb{R}^{N}$. For each $x_{j}$ let $x_{j}^{\star}$ be the integer closest to $x_{j}$.
Before we continue, let's have an example, if $N=5$ and $\vec{x}=(12.85,-12.35,2 / 3,9,-\pi)$ then its associated grid point is $\overrightarrow{x^{\star}}=(13,-12,1,9,-3)$. Let's compute the distance:

$$
d\left(\vec{x}, \overrightarrow{x^{\star}}\right)=\sqrt{.15^{2}+.35^{2}+1 / 3^{2}+0+(\pi-3)^{2}} \leqq \sqrt{(1 / 2+1 / 2+1 / 2+0+1 / 2)} \leqq \sqrt{N}
$$

and we see that part $\boldsymbol{A}$ of the lemma is true for this specific example.
Now to the real proof. It is not really more complicated if you notice that $\left|x_{j}-x_{j}^{\star}\right|<1$ for all $1 \leqq j \leqq N$. We get

$$
d\left(\vec{x}, \overrightarrow{x^{\star}}\right)=\sqrt{\sum_{j=1}^{N}\left(x_{j}-x_{j}^{\star}\right)^{2}}<\sqrt{N \cdot 1}=\sqrt{N}
$$

[^87]So, for each point you can find a grid point with integer coordinates at a distance of less than $\sqrt{N}$. That proves that $\mathbb{Z}^{N}$ is a $\sqrt{N}$-net of $\mathbb{R}^{N}$.

## PROOF of $\boldsymbol{B}$.

Let $\vec{y} \in \mathbb{R}^{N}$. Let $\vec{x}:=(\sqrt{N} / \varepsilon) \vec{y}$ and let $\overrightarrow{x^{\star}}$ be the vector where we discard the decimal parts of $\vec{x}$. According to part $\boldsymbol{A}$. we know that $\left\|\vec{x}-\vec{x}^{\star}\right\|_{2}=d\left(\overrightarrow{x^{\star}}, \vec{x}\right)<\sqrt{N}$. Thus

$$
d\left(\vec{y}, \frac{\varepsilon}{\sqrt{N}} \vec{x}^{\star}\right)=d\left(\frac{\varepsilon}{\sqrt{N}} \vec{x}, \frac{\varepsilon}{\sqrt{N}} \vec{x}^{\star}\right)=\left\|\frac{\varepsilon}{\sqrt{N}} \vec{x}-\frac{\varepsilon}{\sqrt{N}} \vec{x}^{\star}\right\|_{2}=\frac{\varepsilon}{\sqrt{N}}\left\|\vec{x}-\vec{x}^{\star}\right\|_{2}<\frac{\varepsilon}{\sqrt{N}} \sqrt{N}=\varepsilon
$$

In other words, for any $\vec{y} \in \mathbb{R}^{N}$ there is a vector $\vec{x}^{\star} \in \mathbb{Z}^{N}$ such that $d\left(\vec{y},(\varepsilon \sqrt{N}) \vec{x}^{\star}\right)<\varepsilon$.
Rephrase that: For any $\vec{y} \in \mathbb{R}^{N}$ there is a vector $\vec{g} \in G_{(\varepsilon / \sqrt{N})}^{\mathbb{R}^{N}}=\left\{(\varepsilon / \sqrt{N}) \vec{z}: \vec{z} \in \mathbb{Z}^{N}\right\}$
such that $d(\vec{y}, \vec{g})<\varepsilon$ (choose $\vec{g}=\vec{x}^{\star}$ ).
So, for each point you can find a grid point in $G_{\varepsilon}^{\mathbb{R}^{N}}$ at a distance of less than $\varepsilon \sqrt{N}$. It follows that $G_{\varepsilon}^{\mathbb{R}^{N}}$ is an $\varepsilon \sqrt{N}$-net of $\mathbb{R}^{N}$.

## PROOF of $C$.

Intuitively clear but very messy. Here is an outline.
a. You can choose a radius $R_{1}$ so big that $A \subseteq N_{R_{1}}(\overrightarrow{0})$ (see prop. 10.12 on p.200).
b. We enlarge the radius by $\varepsilon$ : Let $R:=R_{1}+\varepsilon$. The enlarged " $N$-dimensional ball" of radius $R N_{R}(\overrightarrow{0})$ is contained in the " $N$-dimensional cube"

$$
Q_{R}:=\left\{\vec{x}=\left(x_{1}, x_{2}, \ldots x_{N}\right):-R \leqq x_{j} \leqq R \text { for all } 1 \leqq j \leqq N\right\}
$$

This is true because if $\vec{x}=\left(x_{1}, x_{2}, \ldots x_{N}\right) \in N_{R}(\overrightarrow{0})$ then $\left|x_{j}\right|=\sqrt{x_{j}^{2}} \leqq \sqrt{\sum_{i} x_{i}^{2}}=\|\vec{x}\|_{2}=R$.
c. Let $\vec{z}=\left(z_{1}, z_{2}, \ldots z_{N}\right)$ be a grid point of $G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}$, i.e., for the $j$-th coordinate $(1 \leqq j \leqq N)$ there exists $m_{j} \in \mathbb{Z}$ such that $z_{j}=m_{j} \frac{\varepsilon}{\sqrt{N}}$. We observe that there are only finitely many integers $m$, say $K$, for which $-R \leqq m \frac{\varepsilon}{\sqrt{N}} \leqq R$. We further observe that $K$ does not depend on $j$.
$\boldsymbol{d}$. Hence there are only $K$ possible values for the first coordinate $z_{1}=m_{1} \frac{\varepsilon}{\sqrt{N}}$. For each one of those there are only $K$ possible values for $z_{2}$, so there are at most $K^{2}$ possible combinations $\left(z_{1}, z_{2}\right)$ for which $\vec{z} \in A$. We keep going and find that there are at most $K^{N}$ possible grid points $\vec{z} \in Q_{R}$.
e. Any point in $\mathbb{R}^{N}$ with distance less than $\varepsilon$ from some point in A must belong to $B_{R}(\overrightarrow{0})$ (now you know why we chose to augment $R_{1}$ by $\varepsilon$ ). In particular, all grid points $g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}$ whose neighborhoods $N_{\varepsilon}(g)$ intersect $A$ belong to $B_{R}(\overrightarrow{0})$ and hence to $Q_{R}$. We conclude that $A \cap N_{\varepsilon}(g)=\emptyset$ for all grid points outside $Q_{R}$, hence

$$
\begin{equation*}
\bigcup\left[A \cap N_{\varepsilon}(g): g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}\right]=\bigcup\left[\left(A \cap N_{\varepsilon}(g)\right) \cap Q_{R}: g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}\right] . \tag{11.6}
\end{equation*}
$$

We know from part B which was already proved that $\mathbb{R}^{N}=\bigcup\left[N_{\varepsilon}(g): g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}\right]$. Hence,

$$
\begin{aligned}
A=A \cap \mathbb{R}^{N} & =A \cap \bigcup\left[N_{\varepsilon}(g): g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}\right] \\
& =\bigcup\left[A \cap N_{\varepsilon}(g): g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}\right]=\bigcup\left[\left(A \cap N_{\varepsilon}(g)\right) \cap Q_{R}: g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}}\right] .
\end{aligned}
$$

We obtain the third equality from prop. 5.2 (Distributivity of unions and intersections) on $p .98$ and the last one follows from (11.6) above.

It follows that $A \subseteq \bigcup\left[N_{\varepsilon}(g): g \in G_{\varepsilon / \sqrt{N}}^{\mathbb{R}^{N}} \cap Q_{R}\right]$. We have proved $C$ as there are only finitely many grid points in $Q_{R}$.

Remark 11.1. The observant reader will have noted that, in part C. of the previous proposition, it was not stated that the gridpoints belong to the subset $A$ of $\mathbb{R}^{N}$. Here is a trivial counterexample. Look at the "standard" $\varepsilon$-grid $G_{\varepsilon}^{\mathbb{R}^{N}}=\left\{\varepsilon \vec{z}: \vec{z} \in \mathbb{Z}^{N}\right\}$ defined in prop.11.2, part B. Take any $A \subseteq \mathbb{R}^{N}$ you like and look at $B:=A \backslash G_{\varepsilon}^{\mathbb{R}^{N}}$, i.e., we have removed all grid points. It is clear that $B$ cannot be covered by $\sqrt{N} \varepsilon$ balls belonging to grid points in $B$.
Definition 11.2 (Total boundedness). Let $(X, d)$ be a metric space and let $A$ be a subset of $X$. We say that $A$ is totally bounded if for each $\varepsilon>0$ there is a finite collection $\mathscr{G}_{\varepsilon}:=\left\{g_{1}, \ldots g_{n}\right\}$ of points in $A$ whose $\varepsilon$-balls $N_{\varepsilon}\left(g_{j}\right)$ cover $A$ : For any $a \in A$ there is $j=j(a)$ such that $d\left(a, g_{j}\right)<\varepsilon$.

We will use this definition in connection with sequence compactness which is defined in the next section.

### 11.4 Sequence Compactness

We saw in the introductory section that, for the space $\mathbb{R}^{N}$ with the Euclidean metric, closed and bounded sets have the property that any sequence contains a convergent subsequence. We named this property in section 11.2, p. 246 on Four Definitions of Compactness "sequence compactness" and we will examine that property in this chapter.
Definition 11.3 (Sequence compactness). Let ( $X, d$ ) be a metric space and let $A$ be a subset of $X$. We say that $A$ is sequence compact or sequentially compact if it has the following property: Given any sequence $\left(x_{n}\right)$ of elements of $A$, there exists $L \in A$ and a subset

$$
n_{1}<n_{2}<\ldots<n_{j}<\ldots \text { of indices such that } L=\lim _{n \rightarrow \infty} x_{n_{j}}
$$

i.e., there exists a subsequence ${ }^{121}\left(x_{n_{j}}\right)$ which converges to $L$.

Proposition 11.3 (Sequence compactness implies total boundedness). Let ( $X, d$ ) be a metric space and let $A$ be a sequentially compact subset of $X$. Then $A$ is totally bounded.

PROOF: This is proved by contradiction. We may assume that $A \neq \emptyset$ as nothing needs to be shown otherwise. a. Assume that $A$ is not totally bounded. Then there is $\varepsilon>0$ such that the following holds for any $n \in \mathbb{N}$ : If $z_{1}, z_{2}, \ldots z_{n} \in A$ then the union $\bigcup_{1 \leq j \leq n} N_{\varepsilon}\left(z_{j}\right)$ does not cover $A$ : There exists $z \in A$ outside any one of those $\varepsilon$-neighborhoods, i.e., $z \in A \backslash \bigcup\left[\bar{N}_{\varepsilon}\left(z_{j}\right): 1 \leqq j \leqq n\right]$.

This allows us to create an infinite sequence $\left(x_{j}\right)_{j \in \mathbb{N}}$ such that $d\left(x_{j}, x_{n}\right) \geqq \varepsilon$ for all $j, n \in \mathbb{N}$ such that $j \neq n$, say, $j<n$, as follows: We pick

$$
x_{1} \in A ; \quad x_{2} \in A \backslash N_{\varepsilon}\left(x_{1}\right) ; \quad x_{3} \in A \backslash\left(N_{\varepsilon}\left(x_{1}\right) \cup N_{\varepsilon}\left(x_{2}\right)\right) ; \ldots x_{n} \in A \backslash \bigcup_{j<n} N_{\varepsilon}\left(x_{j}\right) ; \ldots
$$

b. The proof is done if we can show that $\left(x_{j}\right)_{j \in \mathbb{N}}$ does not possess a convergent subsequence. Assume to the contrary that there is $L \in A$ and $n_{1}<n_{2}<\ldots$ such that $\lim _{j \rightarrow \infty} x_{n_{j}}=L$.

[^88]We pick the number $\varepsilon>0$ that was used in part a of the proof. There exists $N=N(\varepsilon)$ such that $d\left(x_{n_{m}}, L\right)<\varepsilon / 2$ for all $m \geqq N$. Because $i \leqq n_{i}$ and $j \leqq n_{j}$, it follows for all $i, j \geqq N$ that

$$
d\left(x_{n_{i}}, x_{n_{j}}\right) \leqq d\left(x_{n_{i}}, L\right)+d\left(L, x_{n_{j}}\right)<\varepsilon / 2+\varepsilon / 2=\varepsilon .
$$

But the $x_{n}$ were constructed in such a fashion that $d\left(x_{m}, x_{k}\right) \geqq \varepsilon$ for all $m \neq k$, in particular for $m:=n_{i}$ and $k:=n_{j}$ if $i \neq j$. We have arrived at a contradiction because $n_{i} \neq n_{j}$ whenever $i \neq j$.

Proposition 11.4 (Sequence compact implies completeness). Let ( $X, d$ ) be a metric space and let $A$ be a sequence compact subset of $X$. Then $A$ is complete, i.e., any Cauchy sequence $\left(x_{n_{j}}\right)$ in $A$ converges to a limit $L \in A$.

PROOF: Let $\left(x_{n}\right)$ be a Cauchy sequence in $A$. Because $A$ is sequence compact, we can extract a subsequence $z_{j}:=x_{n_{j}}$ and find $L \in A$ such that $z_{j} \rightarrow L$ as $j \rightarrow \infty$.

## EMPTY LABEL! MUST BE REMOVED! EMPTY LABEL! MUST BE REMOVED!

It follows from prop. 10.21 on $p .210$ that the entire Cauchy sequence $\left(x_{n}\right)$ converges to $L$.
The last two propositions have proved that any sequence compact set in a metric space is both totally bounded and complete. The reverse is also true:

Theorem 11.4 (Sequence compact iff totally bounded and complete). Let $A$ be a subset of a metric space $(X, d)$. Then $A$ is sequence compact if and only if $A$ is totally bounded and complete.

PROOF: We have already seen in prop. 11.3 on $p .249$ and prop. 11.4 on $p .250$ that if $A$ is sequentially compact then $A$ is totally bounded and complete. We now show the other direction.

Let $A$ be totally bounded and complete and let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence in $A$. All we need to show is the existence of a subsequence $z_{j}=x_{n_{j}}$ which is Cauchy: As $A$ is complete, such a Cauchy sequence must converge to a limit $L \in A$, i.e., $x_{n_{j}} \rightarrow L$ as $n \rightarrow \infty$. We now are going to extract a convergent subsequence $\left(x_{n_{j}}\right)_{j}$ from $\left(x_{n}\right)_{n}$.
$\boldsymbol{a}$. Because $A$ is totally bounded, there will be a finite net for $\varepsilon=1 / 2$ : there exists

$$
\mathscr{G}_{1}=\left\{g_{1,1}, g_{1,2}, \ldots, g_{1, k_{1}}\right\} \subseteq A \text { such that } A \subseteq U_{1}:=\bigcup\left[N_{1 / 2}\left(g_{1, j}\right): j \leqq k_{1}\right] .
$$

It follows that $x_{k} \in U_{1}$ for each $k$. There are infinitely many indices $k$ for our sequence but only finitely many points in $\mathscr{G}_{1}$. Hence there exists some $g_{1} \in \mathscr{G}_{1}$ such that $B_{1}:=N_{1 / 2}\left(g_{1}\right)$ contains $x_{1, j}:=x_{n_{j}}$ for an entire (infinite) subsequence $n_{j}$. ${ }^{122}$
b. Because $A$ is totally bounded, there will be a finite net for $\varepsilon=1 / 3$ : there exists

$$
\mathscr{G}_{2}=\left\{g_{2,1}, g_{2,2}, \ldots, g_{2, k_{2}}\right\} \subseteq A \text { such that } A \subseteq U_{2}:=\bigcup\left[N_{1 / 3}\left(g_{2, j}\right): j \leqq k_{2}\right] .
$$

It follows that $x_{1, k} \in U_{2}$ for each $k$. There are infinitely many indices $k$ for our sequence but only finitely many points in $\mathscr{G}_{2}$. Hence there exists some $g_{2} \in \mathscr{G}_{2}$, such that $N_{1 / 3}\left(g_{2}\right)$ contains $x_{2, j}:=x_{1, n_{j}}$ for an entire subsequence $n_{j}$. As the entire sequence ( $x_{1, k}$ ) belongs to $B_{1}$, it follows that our new subsequence $\left(x_{2, j}\right)$ of $\left(x_{1, k}\right)$ belongs to $B_{2}:=B_{1} \cap N_{1 / 3}\left(g_{2}\right)$.

[^89]c. Having constructed a subsequence $\left(x_{n-1, j}\right)$ of the original sequence $\left(x_{k}\right)$ which lives in a set $N_{n-1}$ contained in $N_{1 / n}\left(g_{n-1}\right)$ for a suitable $g_{n-1} \in A$, total boundedness of $A$, guarantees the existence of a finite net for $\varepsilon=1 /(n+1)$ : there exists
$$
\mathscr{G}_{n}=\left\{g_{n, 1}, g_{n, 2}, \ldots, g_{n, k_{n}}\right\} \subseteq A \text { such that } A \subseteq U_{n}:=\bigcup\left[N_{1 /(n+1)}\left(g_{n, j}\right): j \leqq k_{n}\right]
$$

It follows that $x_{n, k} \in U_{n}$ for each $k$. There are infinitely many indices $k$ for our sequence but only finitely many points in $\mathscr{G}_{n}$. Hence there must be at least one of those which we name $g_{n}$, such that $N_{1 /(n+1)}\left(g_{n}\right)$ contains $x_{n, j}:=x_{n-1, n_{j}}$ for an entire subsequence $n_{j}$. As the entire sequence ( $x_{n-1, k}$ ) belongs to $N_{n-1}$, it follows that our new subsequence $\left(x_{n, j}\right)$ of $\left(x_{n-1, k}\right)$ belongs to $B_{n}:=N_{n-1} \cap N_{1 /(n+1)}\left(g_{n}\right)$.

We note that the maximal distance $d\left(x_{n, i}, x_{n, j}\right)$ between any two members of that new subsequence is bounded by $2 /(n+1)$, the diameter of $N_{1 /(n+1)}$. It follows that

$$
\operatorname{diam}\left(B_{n}\right) \leqq 2 /(n+1)
$$

d. Diagonalization procedure: The following trick is employed quite frequently in real analysis. We create the "diagonal sequence" $z_{1}:=x_{1,1}, z_{2}:=x_{2,2}, \ldots$ which is a subsequence of the original sequence $\left(x_{n}\right)$. If we can show that it is Cauchy then the proof is complete.

By construction, if $j \geqq n$ then

$$
z_{j} \in B_{j} \subseteq B_{n} \subseteq N_{n-1} \subseteq \cdots \subseteq B_{2} \subseteq B_{1} \quad \text { and } \operatorname{diam}\left(B_{j}\right) \leq \frac{2}{j+1} \leq \frac{2}{n+1} \quad \text { (see (??) }
$$

Let $\varepsilon>0$. We can find $N \in \mathbb{N}$ such that $\frac{1}{N+1}<\frac{\varepsilon}{2}$. We remember from part $c$ of this proof that

$$
B_{j}=N_{j-1} \cap N_{1 /(j+1)}\left(g_{j}\right) \subseteq N_{1 /(j+1)}\left(g_{j}\right) \text { for some suitable } g_{j} \in A \text {, }
$$

hence all its points have distance from $g_{j}$ bounded by $(j+1)^{-1}$. We obtain for any $i, j \geqq N$ that

$$
d\left(z_{i}, z_{j}\right) \leq d\left(z_{i}, g_{N}\right)+d\left(g_{N}, z_{j}\right) \leqq \frac{1}{N+1}+\frac{1}{N+1}<\varepsilon
$$

It follows that $\left(z_{n}\right)_{n}$ is indeed Cauchy and the proof is completed.
Corollary 11.2 (Sequence compact sets are complete). Let ( $X, d$ ) be a metric space and let $K$ be a sequence compact subset of $X$. Then $K$ is complete.

PROOF: Immediate from the last theorem.
Theorem 11.5 (Sequence compact sets are closed and bounded). Let $A$ be sequence compact subset of a metric space $(X, d)$. Then $A$ is a bounded and closed set.

## a. Proof of boundedness:

We may assume that $A$ is not empty because otherwise there is nothing to prove. We assume that $A$ is not bounded, i.e., $\operatorname{diam}(A)=\infty$. It will be proved by induction that there exists a sequence $x_{n} \in A$ such that $d\left(x_{i}, x_{j}\right) \geqq 1$ for any $i \neq j$.

Base case: Let $x_{0} \in A$. There exists $x_{1} \in A$ such that $r_{1}:=d\left(x_{0}, x_{1}\right) \geqq 1$.

Induction step: We assume that $n$ elements $x_{1}, \ldots x_{n}$ such that $d\left(x_{i}, x_{j}\right) \geqq 1$ for any $1 \leqq i<j \leqq n$ have aready been chosen. Let $\beta:=\max \left\{d\left(x_{0}, x_{j}\right): j \leqq n\right\}$ and $r:=\beta+1$. As $A$ is not bounded, we can pick $x_{n+1} \in A \backslash N_{r}\left(x_{0}\right)$. We obtain

$$
\beta+1 \leq d\left(x_{n+1}, x_{0}\right) \leq d\left(x_{n+1}, x_{j}\right)+d\left(x_{j}, x_{0}\right) \leq d\left(x_{n+1}, x_{j}\right)+\beta \text {, i.e., } 1 \leq d\left(x_{n+1}, x_{j}\right) .
$$

We have constructed a sequence $\left(x_{n}\right)$ for which any two items have distance no less than 1 . It follows that there is no Cauchy subsequence, hence no convergent subsequence and we have a contradiction to the sequence compactness of $A$.
b. Proof of closedness:

If $A$ is not closed then $A$ has a contact point $x \in A^{\complement}$.
As $N_{1 / m}(x) \cap A \neq \emptyset$ we can pick a sequence $x_{m} \in A$ such that $d\left(x_{m}, x\right)<1 / m$ for all $m \in \mathbb{N}$.
Clearly $x_{m}$ converges to $x \in A^{\complement}$. Sequence compactness of $A$ allows us to extract a subsequence $z_{j}=x_{n_{j}}$ which converges to $z \in A$.

Both $z$ and $x$ are limit of $z_{j}$. According to thm.10.3 on $p .192, x=z$. It follows that both $x \in A^{\complement}$ and $x \in A$, a contradiction. This proves that sequence compact sets are closed.

Corollary 11.3 (Sequence compact sets are bounded). Let $(X, d)$ be a metric space and let $K$ be a sequence compact subset of $X$. Then $K$ is bounded, i.e., $\operatorname{diam}(K)=\sup \{d(x, y): x, y \in K\}<\infty$.

PROOF: Obvious from thm.11.5.
Remark 11.2. It follows from the results of this chapter and the introductory chapter on Closed and bounded sets in Euclidean space (11.1 on p.244) that the following three of the four definitions of compactness given in section 11.2 on (p. 246 are equivalent in $\mathbb{R}^{N}$ :

A subset of $\mathbb{R}^{N}$ is sequentially compact
$\Leftrightarrow$ it is totally bounded and complete
$\Leftrightarrow$ it is bounded and closed.
We will see later that any metric space is sequentially compact if and only if it is compact, i.e., covering compact (thm. 11.8 on p.255).

In other words, in $\mathbb{R}^{\mathbb{N}}$ all four of the definition given in section 11.2 on p .246 coincide.

### 11.5 Open Coverings and the Heine-Borel Theorem

We now discuss families of open sets called "open coverings". You should review the concept of an indexed family and how it differs from that of a set (see (4.17) on p.86).
Definition 11.4 (Coverings and open coverings). Let $X$ be an arbitrary nonempty set and $A \subseteq$ $X$. Let $\left(U_{i}\right)_{i \in I}$ be an indexed family of subsets of $X$ such that $A \subseteq \bigcup_{i \in I} U_{i}$. Then we call $\left(U_{i}\right)_{i \in I}$ a covering of $A$.

A finite subcovering of a covering $\left(U_{i}\right)_{i \in I}$ of the set $A$ is a finite collection

$$
\begin{equation*}
U_{i_{1}}, U_{i_{2}}, U_{i_{3}}, \ldots, U_{i_{n}} \quad\left(i_{j} \in I \quad \text { for } 1 \leqq j \leqq n\right) \quad \text { such that } \quad A \subseteq U_{i_{1}} \cup U_{i_{2}} \cup \ldots \cup U_{i_{n}} . \tag{11.9}
\end{equation*}
$$

Assume in addition that $X$ is a topological space - this includes any normed vector space or metric space. If all members $U_{i}$ of the family $\left(U_{i}\right)_{i \in I}$ are open then we call this family an open covering of $A$.

## Remark 11.3. Remark:

a. Partitions ${ }^{123}$ are coverings.
b. If $\left(U_{i}\right)_{i \in I}$ is a covering of $A$ then $\left(U_{i} \cap A\right)_{i \in I}$ is a covering of $A$ which satisfies

$$
\begin{equation*}
\bigcup_{j \in I}\left(U_{j} \cap A\right)=A . \tag{11.10}
\end{equation*}
$$

Definition 11.5 (Compact sets). Let $(X, \mathfrak{U})$ be a topological space and $K \subseteq X$. We say that $K$ is compact if it possesses the "extract finite open subcovering" property Given any open covering $\left(U_{i}\right)_{i \in I}$ of $K$, one can extract a finite subcovering. In other words, there is $n \in \mathbb{N}$ and indices

$$
i_{1}, i_{2}, \ldots, i_{n} \in I \quad \text { such that } \quad A \subseteq \bigcup_{j=1}^{n} U_{i_{j}}
$$

## Remark 11.4.

a. An open covering for the entire space $X$ is a collection of open sets $\left(U_{i}\right)_{i \in I}$ such that $X=\bigcup\left[U_{i}: i \in I\right]$.
b. Any subcovering of an open covering necessarily consists exclusively of open sets, i.e., it is again an open covering of $A$.
c. Let $(X, d)$ be a metric space. Then

$$
K \subseteq(X, d) \text { is compact if and only if the metric subspace }\left(K,\left.d\right|_{K \times K}\right) \text { is compact, }
$$

i.e., for any collection of subsets $\left(U_{i}\right)_{i \in I}$ of $K$ which are open in $K$ there exist finitely many indices $i_{1}, \ldots, i_{n} \in I$ such that $K=U_{i_{1}} \cup \cdots \cup U_{i_{n}}$. This is true because the open subsets of ( $K, d$ ) are the traces in $K$ of sets which are open in ( $X, d$ ) (see def.10.22 on p.198).

Example 11.1. Here are some simple examples.
a. Any finite topological space is compact.
b. Any topological space that only contains finitely many open sets is compact. In particular a set with the indiscrete topology (example 10.6 on p.194) is compact
c. A space with the discrete metric (def.10.3 on p.186) is compact if and only if it is finite.

And here is a counterexample.
The open interval $] 0,1[$ with the Euclidean metric is not compact because it is not possible to extract an finite covering from the open covering (]$\frac{1}{n}, 1[)_{n \in \mathbb{N}}$.

Example 11.2. Here are the corresponding results for sequence compact metric spaces.
a. Any finite metric space is sequence compact.
b. Any metric space that only contains finitely many open sets is compact. ${ }^{124}$
c. A space with the discrete metric is compact if and only if it is finite.

The counterexample also fits in:

[^90]The open interval $] 0,1[$ with the Euclidean metric is not sequence compact because it is not possible to extract a convergent subsequence from the sequence $x_{n}:=1 / n$ (the limit zero does not belong to $] 0,1[$ ).

We will now see that the correspondence in the above two examples is not a coincidence. We will see that (subspaces of) metric spaces are compact if and only if they are sequentially compact. We will prove each direction separately.

Theorem 11.6 (Compact metric spaces are sequence compact). Let $(X, d)$ be a compact metric space. Then $X$ is sequence compact.

PROOF: We assume to the contrary that $X$ is compact and that there is a sequence $\left(x_{n}\right)_{n}$ in $X$ from which one cannot extract a convergent subsequence.

Let $F:=\left\{x \in X: x=x_{j} \text { for some } j \in \mathbb{N}\right\}^{125}$ be the set of distinct(!) members of $\left(x_{n}\right)_{n}$. Let $z \in X$. We first prove that there exists an open neighborhood $U_{z}$ of $z$ such that $U_{z} \cap F$ is finite.

This is true because otherwise for each for each $m \in \mathbb{N}$ there exists some index $j_{m}$ such that $x_{j_{m}} \neq z$ and $x_{j_{m}} \in N_{1 / m}(z)$. We have constructed a sequence $\left(x_{j_{m}}\right)_{m}$ which converges to $z$, i.e., $\left(x_{n}\right)_{n}$ possesses a convergent subsequence, contrary to our assumption.

It follows from $\{z\} \subseteq U_{z}$ that $\left(U_{z}\right)_{z \in X}$ is an open covering of $X . X$ is compact and we can extract a finite subcovering $U_{1}, U_{2}, \ldots, U_{k}$. We proved that each $U_{j}$ contains at most finitely many distinct members of the sequence $\left(x_{n}\right)_{n}$. We conclude that this entire sequence consists of only finitely many distinct members.

It follows that at least one of those members, say $x_{k^{*}}$, will appear infinitely often in that sequence: there is $k_{1}<k_{2}<\ldots$ such that $x_{k_{1}}=x_{k_{2}}=\cdots=x_{k^{\star}}$. This constant subsequence converges (to $x_{k^{\star}}$ ). We have reached a contradiction.

Proposition 11.5. Let $(X, d)$ be a sequence compact metric space. Let $\left(U_{i}\right)_{i \in I}$ be an open covering of $X$. Then there exists $\rho>0$ as follows: For each $x \in X$ there exists $i \in I$ such that $N_{\rho}(x) \subseteq U_{i}$. ${ }^{126}$

PROOF: Assume to the contrary that no such $\rho>0$ exists. We then can find for any $n \in \mathbb{N}$ some $x_{n} \in X$ such that $N_{1 / n}\left(x_{n}\right)$ is not contained in any of the $U_{i} . X$ is sequence compact, so there exists $x \in X$ and $a$ subsequence $\left(x_{n_{j}}\right)_{j}$ which converges to $x$. $\left(U_{i}\right)_{i \in I}$ covers $X$, so there exists $i_{0} \in I$ such that $x \in U_{i_{0}}$.
$(\star)$ Because $U_{i_{0}}$ is an open neighborhood of $x$ there exists $\varepsilon>0$ such that $N_{\varepsilon}(x) \subseteq U_{i_{0}}$.
( $\star \star$ ) Because $\left(x_{n_{j}}\right)_{j}$ converges to $x$ there are infinitely many $j \in \mathbb{N}$ such that $d\left(x_{n_{j}}, x\right)<\varepsilon / 2$, hence there is at least one $j$ such that $j>2 / \varepsilon$. It follows from $n_{j} \geqq j$ that $n_{j}>2 / \varepsilon$, i.e. $1 / n_{j}<\varepsilon / 2$.
$(\star \star \star)$ It follows from $d\left(x_{n_{j}}, x\right)<\varepsilon / 2$ and lemma 10.3 on $p .190$ that $N_{\varepsilon / 2}\left(x_{n_{j}}\right) \subseteq N_{\varepsilon / 2+\varepsilon / 2}(x)=$ $N_{\varepsilon}(x)$.

We apply first $\star \star$, then $\star \star \star$, then $\star$ and obtain

$$
N_{1 / n_{j}}\left(x_{n_{j}}\right) \subseteq N_{\varepsilon / 2}\left(x_{n_{j}}\right) \subseteq N_{\varepsilon}(x) \subseteq U_{i_{0}} .
$$

But this contradicts our assumption that each $x_{n_{j}}$ was chosen in such a fashion that $N_{1 / n}\left(x_{n}\right)$ is not contained in any of the $U_{i}$.

[^91]We now can prove the converse of thm.11.6.
Theorem 11.7. Sequence compact metric spaces are compact.
PROOF: Let $(X, d)$ be a sequence compact metric space and let $\left(U_{i}\right)_{i \in I}$ be an open covering of $X$. According to prop.11.5 there exists $\rho>0$ as follows: For each $x \in X$ there exists $i(x) \in I$ such that $N_{\rho}(x) \subseteq U_{i(x)}$.

Since $X$ is totally bounded (see thm.11.4 on $p$.250) there exist finitely many $x_{1}, \ldots, x_{k} \in X$ such that $\left\{N_{\rho}\left(x_{j}\right): j=1, . . k\right\}$ forms an open covering of $X$. It then follows from $N_{\rho}\left(x_{j}\right) \subseteq U_{i\left(x_{j}\right)}$ that $U_{i\left(x_{1}\right)}, U_{i\left(x_{2}\right)}, \ldots, U_{i\left(x_{k}\right)}$ also forms an open covering of $X$. We have extracted a finite subcovering from $\left(U_{i}\right)_{i \in I}$.

Theorem 11.8 (Sequence compact is same as compact in metric spaces). Let ( $X, d$ ) be a metric space and let $A$ be a subset of $X$. Then $A$ is sequence compact if and only if $A$ is compact, i.e., every open covering of $A$ has a finite subcovering.

## PROOF: Theorems 11.6 and 11.7.

An easy consequence is the Heine-Borel theorem.
Theorem 11.9 (Heine-Borel). Let $(X, d)$ be $\mathbb{R}^{N}$ with the Euclidean norm and its associated metric. A subset $K \subseteq \mathbb{R}^{N}$ is compact if and only if it is closed and bounded.

For a general metric space it is still true that any compact subset is closed and bounded.
PROOF: We have seen in thm. 11.2 on $p .245$ that closed and bounded subsets of $\mathbb{R}^{N}$ are sequence compact. We have further seen that sequence compact sets are closed and bounded (thm.11.5 on p.251). Our assertion now follows from the equivalence of compactness and sequence compactness in metric spaces (thm. 11.8 on p.255).

### 11.6 Continuous Functions and Compact Spaces

Theorem 11.10 (Closed subsets of compact topological spaces are compact). Let $A$ be a closed subset of a compact topological space $(X, \mathfrak{U})$. Then $A$ is a compact subspace, i.e., the open sets

$$
\mathfrak{U}_{A}=\{V \cap A: V \in \mathfrak{U}\}
$$

possess the "extract finite open subcovering" property of def. 11.5 on p. 253.
PROOF: Let $\left(U_{j}\right)_{j \in J}$ be a family of sets open in A whose union is A. According to def. 10.23 on $p .199$ there are open sets $V_{j}$ in $X$ such that $U_{j}=V_{j} \cap A$. It follows that $\bigcup_{j \in J} V_{j} \supseteq A$, hence the family $\left(V_{j}\right)_{j \in J}$, augmented by the (open!) set $X \backslash A$ is an open cover of $(X, d)$.

As $X$ is compact, we can extract finitely many members from that extended family such that they still cover $X$. If one of them happens to be $X \backslash A$ then we remove it and we still obtain that the remaining ones, say, $V_{i_{1}}, V_{i_{2}}, \ldots, V_{i_{n}}$, cover $A$. But then the traces in $A$ (def.10.22 on $p .198$ )

$$
U_{i_{1}}=V_{i_{1}} \cap A, U_{i_{2}}=V_{i_{2}} \cap A, \ldots, U_{i_{n}}=V_{i_{n}} \cap A
$$

of those open sets in $X$ are open in $A$ and hence form an open covering of the subspace $\left(A, \mathfrak{U}_{A}\right)$. We have proved that the given open covering of $A$ has contains a finite subcover of $A$.

Corollary 11.4 (Closed subsets of compact metric spaces are compact). Let $A$ be a closed subset of a compact metric space $(X, d)$. Then $\left(A,\left.d\right|_{A \times A}\right)$ is a compact subspace.

PROOF: Immediate from thm.11.10.
Let $\left(X, \mathfrak{U}_{1}\right)$ and $\left(Y, \mathfrak{U}_{2}\right)$ be topological spaces and $A \subseteq X$. We recall that continuity for functions $f: A \rightarrow\left(Y, \mathfrak{U}_{2}\right)$ was defined in def. 10.32 on $p .215$.

Theorem 11.11 (Continuous images of compact topological spaces are compact). Let ( $X, \mathfrak{U}_{1}$ ) and $\left(Y, \mathfrak{U}_{2}\right)$ be two topological spaces. and let $f: X \rightarrow Y$ be continuous on $X$. If $X$ is compact then the direct image $f(X)$ is compact, i.e., the topological subspace $\left(f(X), \mathfrak{U}_{2 f(X)}\right)$ of $Y$ is compact.

PROOF: Let $\left(V_{j}\right)_{j \in J}$ be a family of sets open in $Y$ whose union contains $B:=f(X)$. Let the sets $W_{j}:=$ $V_{j} \cap f(X)$ be the traces of $V_{j}$ in $f(X)$. Then the $W_{j}$ are open in the subspace $\left(f(X), \mathfrak{U}_{2 f(X)}\right)$ of $Y$ and they form an open cover of $f(X)$. We note that any open cover of $f(X)$ is obtained in this manner from open sets in $Y$.

Let $U_{j}:=f^{-1}\left(V_{j}\right)$. Then

$$
\begin{equation*}
\bigcup_{j \in J} U_{j}=\bigcup_{j \in J} f^{-1}\left(V_{j}\right)=f^{-1}\left(\bigcup_{j \in J} V_{j}\right) \supseteq f^{-1}(B)=f^{-1}(f((X))=X \tag{11.11}
\end{equation*}
$$

In the above the second equation follows from prop. $6.2\left(f^{-1}\right.$ is compatible with all basic set ops) on $p .104$ and the last one follows from the fact that $f^{-1}(f((\Gamma)) \supseteq \Gamma$ for any subset $\Gamma$ of the domain of $f$ (see cor. 6.1 on $p$. 107).

According to prop.10.22 (" $f^{-1}$ (open) $=$ open" continuity) on $p .215$, each $U_{j}$ is open as the inverse image of the open set $V_{j}$ under the continuous function $f$.

It follows from (11.11) that $\left(U_{j}\right)_{j \in J}$ is an open covering of the compact space $X$. We can extract a finite subcover $U_{i_{1}}, U_{i_{2}}, \ldots, U_{i_{n}}$.

It follows from the interchangeability of unions with direct images (see (6.19) on p.105) that

$$
\begin{aligned}
f(X) & =f\left(U_{j_{1}} \cup \cdots \cup U_{j_{n}}\right)=f\left(U_{j_{1}}\right) \cup \cdots \cup f\left(U_{j_{n}}\right) \\
& =f\left(f^{-1}\left(V_{j_{1}}\right)\right) \cup \cdots \cup f\left(f^{-1}\left(V_{j_{n}}\right)\right) \subseteq V_{j_{1}} \cup \cdots \cup V_{j_{n}} .
\end{aligned}
$$

The inclusion relation above follows from the fact that $f\left(f^{-1}(B)\right)=B \cap f(X)$ for any subset $B$ of the codomain of $f$ (see prop.6.6 on p. 107).

We have proved that the arbitrary open cover $\left(V_{j}\right)_{j \in J}$ of $f(X)$ contains a finite subcover $V_{j_{1}}, \ldots, V_{j_{n}}$ and it follows that $f(X)$ is indeed a compact metric subspace of $Y$.

Corollary 11.5 (Continuous images of compact metric spaces are compact). Let ( $X, d_{1}$ ) and ( $Y, d_{2}$ ) be two metric spaces. and let $f: X \rightarrow Y$ be continuous on $X$. If $X$ is compact then the direct image $f(X)$ is compact, i.e., the metric subspace $\left(f(X), d_{2}\right)$ of $Y$ is compact.

PROOF: Let $\left(V_{j}\right)_{j \in J}$ be a family of sets open in $Y$ whose union contains $B:=f(X)$. Let the sets $W_{j}:=$ $V_{j} \cap f(X)$ be the traces of $V_{j}$ in $f(X)$. Then the $W_{j}$ are open in the subspace $\left(f(X), d_{2}\right)$ of $Y$ and they form an open cover of $f(X)$. We note that any open cover of $f(X)$ is obtained in this manner from open sets in $Y$. Let $U_{j}:=f^{-1}\left(V_{j}\right)$. Then

$$
\begin{equation*}
\bigcup_{j \in J} U_{j}=\bigcup_{j \in J} f^{-1}\left(V_{j}\right)=f^{-1}\left(\bigcup_{j \in J} V_{j}\right) \supseteq f^{-1}(B)=f^{-1}(f((X))=X \tag{11.12}
\end{equation*}
$$

In the above the second equation follows from prop. 6.2 ( $f^{-1}$ is compatible with all basic set ops) on $p .104$ and the last one follows from the fact that $f^{-1}(f((\Gamma)) \supseteq \Gamma$ for any subset $\Gamma$ of the domain of $f$ (see cor. 6.1 on $p$. 107).

According to prop.10.22 (" $f^{-1}$ (open) $=$ open" continuity) on p.215, each $U_{j}$ is open as the inverse image of the open set $V_{j}$ under the continuous function $f$. It follows from (11.12) that $\left(U_{j}\right)_{j \in J}$ is an open covering of the compact space $X$. We can extract a finite subcover $U_{i_{1}}, U_{i_{2}}, \ldots, U_{i_{n}}$.

It follows from the interchangeability of unions with direct images (see (6.19) on p.105) that

$$
\begin{aligned}
f(X) & =f\left(U_{j_{1}} \cup \cdots \cup U_{j_{n}}\right)=f\left(U_{j_{1}}\right) \cup \cdots \cup f\left(U_{j_{n}}\right) \\
& =f\left(f^{-1}\left(V_{j_{1}}\right)\right) \cup \cdots \cup f\left(f^{-1}\left(V_{j_{n}}\right)\right) \subseteq V_{j_{1}} \cup \cdots \cup V_{j_{n}} .
\end{aligned}
$$

The inclusion relation above follows from the fact that $f\left(f^{-1}(B)\right)=B \cap f(X)$ for any subset $B$ of the codomain of $f$ (see prop.6.6 on $p$. 107).

We have proved that the arbitrary open cover $\left(V_{j}\right)_{j \in J}$ of $f(X)$ contains a finite subcover $V_{j_{1}}, \ldots, V_{j_{n}}$ and it follows that $f(X)$ is indeed a compact metric subspace of $Y$.

Read the following remark for an easier way to prove the above theorem.
Remark 11.5. We could have proved the last two theorems more easily using sequence compactness instead of covering compactness but the proofs that were given generalize to abstract topological spaces.

Alternate proof of cor. 11.4 which uses sequence compactness.
Given is a sequence $x_{n} \in A . X$ is compact, hence sequence compact and it follows that there is $x \in X$ and a subsequence $x_{n_{j}} \in A$ such that $x_{n_{j}}$ converges to $x$. It follows from theorem 10.6 (Sequence criterion for contact points in metric spaces) on p .201 that $x$ is a contact point of $A$ and hence $x \in \bar{A}=A$. This proves that $A$ is (sequence) compact.

Alternate proof of theorem 11.5 which uses sequence compactness (outline).
Given a sequence $y_{n} \in f(X)$ we construct a convergent subsequence $y_{n_{j}}$ as follows: For each $n$ there is some $x_{n} \in X$ such that $y_{n}=f\left(x_{n}\right) X$ is compact, hence sequence compact and it follows that there is $x \in X$ and a subsequence $x_{n_{j}}$ such that $x_{n_{j}}$ converges to $x$. We now use (sequence) continuity of $f$ at $x$ to conclude that $y_{n_{j}}=f\left(x_{n_{j}}\right)$ converges to $f(x) \in f(X)$.
Corollary 11.6 (Continuous real valued functions attain their max and min on a compact domain). Let $(X, \mathfrak{U})$ be a topological space and let $A \subseteq X$ be a compact subspace. Let $f: A \rightarrow \mathbb{R}$ be continuous on $A$. Then there exist $x_{\star}, x^{\star} \in A$ such that

$$
f\left(x_{\star}\right)=\min _{x \in A} f(x) \text { and } f\left(x^{\star}\right)=\max _{x \in A} f(x) .
$$

PROOF: It follows from thm. 11.11 on p. 256 and thm.thm-x:seq-conv-is-closed-bound (Sequence compact sets are closed and bounded) on $p .251$ that $f(X)$ is closed and bounded in $\mathbb{R}$. It follows from exercise 10.14 on p. 241 that $\min ((f(X))$ and $\max (f(X))$ exist, i.e., according to the definition of preimages, there exist elements in the domain $A$ of $f$ which are mapped to those two values.

The following theorem relates compactness and uniform continuity. ${ }^{127}$
Theorem 11.12 (Uniform continuity on sequence compact spaces). Let $\left(X, d_{1}\right),\left(Y, d_{2}\right)$ be metric spaces and let $A$ be a compact subset of $X$. Then any continuous real valued function on $A$ is uniformly continuous on $A$.

[^92]PROOF: Let us assume that $f(\cdot)$ is continuous but not uniformly continuous and find a contradiction. Because $f(\cdot)$ is not uniformly continuous, there exists $\varepsilon>0$ such that no $\delta>0$, however small, will satisfy (10.66) on $p .220$ for all pairs $x$, $y$ such that $d_{1}(x, y)<\delta$. Looking specifically at $\delta:=1 / j$ for all $j \in \mathbb{N}$, we can find $x_{j}, x_{j}^{\prime} \in A$ such that

$$
\begin{equation*}
d_{1}\left(x_{j}, x_{j}^{\prime}\right)<\frac{1}{j} \quad \text { but } \quad d_{2}\left(f(x), f\left(x^{\prime}\right)\right) \geqq \varepsilon . \tag{11.13}
\end{equation*}
$$

Because $A$ is compact, it is sequence compact. There is a subsequence $\left(x_{j_{k}}\right)$ of the $x_{j}$ which converges to an element $x \in A$ We have

$$
\begin{equation*}
d_{1}\left(x_{j_{k}}^{\prime}, x\right) \leqq d_{1}\left(x_{j_{k}}^{\prime}, x_{j_{k}}\right)+d_{1}\left(x_{j_{k}}, x\right) \leqq \frac{1}{j_{k}}+d_{1}\left(x_{j_{k}}, x\right) \tag{11.14}
\end{equation*}
$$

Both right hand terms converge to zero as $k \rightarrow \infty$. This is obvious for $1 / j_{k}$ because $j_{k} \geqq k$ for all $k$ and it is true for $d_{1}\left(x_{j_{k}}, x\right)$ because $x_{j_{k}}$ converges to $x$.

It follows from (11.14) that $\left(x_{j_{k}}^{\prime}\right)$ also converges to $x$. It follows from the ordinary continuity of $f(\cdot)$ that

$$
f(x)=\lim _{k \rightarrow \infty} f\left(x_{j_{k}}^{\prime}\right)=\lim _{k \rightarrow \infty} f\left(x_{j_{k}}\right) .
$$

It follows from the "ordinary" (non-uniform) convergence of sequences that there exist $N, N^{\prime} \in \mathbb{N}$ such that

$$
d_{2}\left(f(x), f\left(x_{j_{k}}\right)\right)<\frac{\varepsilon}{2} \text { for } k \geqq N ; \quad d_{2}\left(f(x), f\left(x_{j_{k}}^{\prime}\right)\right)<\frac{\varepsilon}{2} \text { for } k \geqq N^{\prime}
$$

Both inequalities are true whenever $k \geqq \max \left(N, N^{\prime}\right)$. It follows for all such $k$ that

$$
\begin{equation*}
d_{2}\left(f\left(x_{j_{k}}\right), f\left(x_{j_{k}}^{\prime}\right)\right)<d_{2}\left(f\left(x_{j_{k}}\right), f(x)\right)+d_{2}\left(f(x), f\left(x_{j_{k}}^{\prime}\right)\right)<\frac{\varepsilon}{2}+\frac{\varepsilon}{2}=\varepsilon \tag{11.15}
\end{equation*}
$$

and we have a contradiction to (11.13).
Corollary 11.7 (Uniform continuity on closed intervals). Let $a, b$ be two real numbers such that $a \leqq b$. Any continuous real valued function on the closed interval $[a, b]$ is uniformly continuous on $[a, b]$ :

For any $\varepsilon>0$, there exists $\delta>0$ such that

$$
\begin{equation*}
|f(x)-f(y)|<\varepsilon \text { for all } x, y \in[a, b] \text { such that } \mid x-y) \mid<\delta \tag{11.16}
\end{equation*}
$$

PROOF: This follows from the previous theorem (11.12) because closed intervals $[a, b]$ are closed and bounded sets and, in $\mathbb{R}$, any closed and bounded set is sequence compact.

### 11.7 Exercises for Ch. 11

Exercise 11.1. Let $N \in \mathbb{N}$. Let $X:=\left\{x_{1}, x_{2}, \ldots, x_{N}\right\}$ be a finite set with a metric $d(\cdot, \cdot)$ (so $(X, d)$ is a metric space). Prove that $X$ is compact three different ways:
a. Show sequence compactness to prove that $X$ is compact.
b. Show that $X$ has the "extract finite open subcovering" property to prove that it is compact.
c. Show that $X$ is complete and totally bounded to prove that it is compact.

## Hints:

a. ANY sequence in $X$ possesses a constant subsequence (WHY?)
b. If $\left(U_{i}\right)_{i}$ covers $X$ then for each $x$ there exists (at least one) $i$ such that $x \in U_{i}$ (WHY?) How many of those $U_{i}$ do you need to cover $X$ if $X$ has only $N$ elements?
c. Prop. 10.21 on $p .210$ should prove useful.

## 12 Applications of Zorn's Lemma

### 12.1 More on Partially Ordered Sets (Study this!)

## Some of this was copied almost literally from [6] Dudley.

Definition 12.1 (Linear orderings). Given is a nonempty set $X$ and a partial ordering $\preceq$ on $X$ (see def.4.4 on p.69). $\preceq$ is a linear ordering, also called a total ordering of $X$ if and only if, for all $x$ and $y \in X$ such that $x \neq y$, either $x \preceq y$ or $y \preceq x$. We call $(X, \preceq)$ is called a linearly ordered set or a totally ordered set set.
Example 12.1. The real numbers line $(\mathbb{R}, \leq)$ with its usual " $\leqq$ " ordering is a linearly ordered set. So is $(\mathbb{R}, \geq)(!)$

Definition 12.2. A chain is a linearly ordered subset of a partially ordered set (with the same ordering). In a partially ordered set ( $X, \preceq$ ), an element $m$ of $X$ is called maximal iff there is no $x \neq m$ with $m \preceq x$. A maximum of $X$ is an $m \in X$ such that $x \preceq m$ for all $x \in X$.

Note 12.1 (Notes on maximal elements and maxima).
a. If $X$ is not linearly ordered, it may have many maximal elements. For example, for the trivial partial ordering $x \preceq y$ if and only if $x=y$, every element is maximal. A maximum is a maximal element, but the converse is often not true.
b. If an ordering is not specified, then we always mean set inclusion.
c. If $m \in X$ is a maximum of $X$ then this implies that $m$ must be related to all other elements of $X$.

Example 12.2 (Maximal elements and maxima). Let $X$ be the collection of all intervals $[a, b] \in \mathbb{R}$ of length $b-a \leq 2$ such that $a \leq b$. These intervals are partially ordered by inclusion. Any interval of length equal to 2 is a maximal element. There is no maximum.

Axiom 12.1 (Zorn's Lemma). A hundred years ago the following was seen as extremely controversial by mathematicians who specialize in the foundations of mathematics.

Zorn's Lemma: Let ( $X, \preceq$ ) be a partially ordered set with the ZL property:
Every chain $C \subseteq X$, possesses an upper bound $u \in X$, i.e., $x \preceq u$ for all $x \in C$.
Then $X$ has a maximal element.
Remark 12.1. Zorn's lemma is an axiom rather than a theorem or a proposition in the following sense: It is impossible to verify its truth or falsehood from the axioms of "a" (meaning there are more than one) "reasonable" axiomatic set theory. In that sense mathematicians are free to accept or reject this statement when building their mathematical theories. Two notes on that remark:
a. Today the mathematicians who refuse to accept proofs which make use, directly or indirectly, of Zorn's lemma, are a very small minority.
b. It can be proved that if one accepts (rejects) Zorn's lemma as a mathematical tool then this is equivalent to accepting (rejecting) the Axiom of Choice ${ }^{128}$ which is so abstract in its precise formulation that it was relegated to a footnote.

[^93]c. Moreover the Axiom of Choice, hence Zorn's lemma, is equivalent to prop.4.3b on p.82: If $A, B$ are not empty and $\varphi: A \rightarrow B$ is surjective then $\varphi$ has a right inverse, i.e., a function $\psi: B \rightarrow A$ such that $\varphi \circ \psi=i d_{B}$ For a proof see http:/ / planetmath.org / surjectionandaxiomofchoice.

We will see now how Zorn's lemma allows a surprisingly simple proof to the effect that any vector space has a basis.

### 12.2 Existence of Bases in Vector Spaces (Study this!)

The following is thematically a continuation of the material in chapter 9 (Vectors and vector spaces).
We now focus on proving that every vector space, even if it does not possess a finite subset which spans the entire space, possesses a basis.

For the remainder of this chapter we assume that $V$ is a vector space and that $\mathfrak{B}$ denotes the set

$$
\begin{equation*}
\mathfrak{B}:=\{A \subseteq V: A \text { is linearly independent }\} . \tag{12.1}
\end{equation*}
$$

$\mathfrak{B}$ is a partially ordered set with respect to set inclusion. The next lemma allows us to apply Zorn's lemma.
Lemma 12.1. Every chain ${ }^{129} \mathfrak{C}$ in $(\mathfrak{B}, \subseteq)$ possesses an upper bound.
Proof:
Let $U:=\bigcup[C: C \in \mathfrak{C}]$. We will show that $U \in \mathfrak{B}$. As $U \supseteq C$ for all $C \in \mathfrak{C}$ it then follows that $U$ is an upper bound of $\mathfrak{C}$ and the proof is finished.

Let $x_{1}, x_{2}, \ldots x_{k} \in U$ and $\alpha_{1}, \alpha_{2}, \ldots \alpha_{k} \in \mathbb{R}(k \in \mathbb{N})$ such that

$$
\begin{equation*}
\sum_{j=1}^{k} \alpha_{j} x_{j}=0 \tag{12.2}
\end{equation*}
$$

We must show that each $\alpha_{j}$ is zero. For each $0 \leqq j \leqq k$ there is some $C_{j} \in \mathfrak{C}$ such that $x_{j} \in C_{j}$. $\mathfrak{C}$ is totally ordered, hence $C_{i} \subseteq C_{j}$ or $C_{j} \subseteq C_{i}$ for any two indices $0 \leqq i, j \leqq k$. But then there exists an index $j_{0}$ such that $C_{j_{0}} \supseteq C_{j}$ for all $j$, hence $x_{1}, x_{2}, \ldots x_{k} \in C_{j_{0}}$. But $C_{j_{0}}$ is linearly independent because $C_{j_{0}} \in \mathfrak{C} \subseteq \mathfrak{A}$. It follows that $\alpha_{1}=\cdots=\alpha_{k}=0$.

Theorem 12.1. Every vector space $V$ has a basis.
PROOF: It follows from lemma 12.1 and Zorns Lemma (axiom 12.1 on $p$. 259) that the set $\mathfrak{B}$ of all independent subsets of the vector space $V$ contains a maximal element (subset of $V$ ) which we denote by $B$. As membership in $\mathfrak{B}$ guarantees its linear independence we only need to prove that $\operatorname{span}(B)=V$.

Let us assume to the contrary that there exists $y \in \operatorname{span}(B)^{\complement}$. It follows from lemma 9.2 on $p .167$ that the set $B^{\prime}:=B \cup\{y\}$ is linearly independent, hence $B^{\prime} \in \mathfrak{B}$. Clearly, $B \subsetneq B^{\prime}$. This contradicts the maximality of $B$ in the partially ordered set $(\mathfrak{B}, \subseteq)$.

### 12.3 The Cardinal Numbers are a totally ordered set

As another application of Zorn's Lemma we now prove thm.7.6, p.125, of ch.7.1.1 (Cardinality as a Partial Ordering).

[^94]Theorem 12.2. Let $X, Y \subseteq \Omega$. Then $\operatorname{card}(X) \leqq \operatorname{card}(Y) \operatorname{or} \operatorname{card}(Y) \leqq \operatorname{card}(X)$
PROOF: ${ }^{130}$ The result is immediate if $X=\emptyset$ or $y=\emptyset$. Assume $X$ and $Y$ are not empty. Let

$$
\mathscr{F}:=\left\{D_{f} \xrightarrow{f} C_{f}: D_{f} \subseteq X, C_{f} \subseteq Y, \text { and } f \text { is injective }\right\}
$$

be the set of all injective functions with domain contained in $X$ and codomain contained in $Y$. We define a partial order $\preceq$ on $\mathscr{F}$ as follows: Let $D_{f} \xrightarrow{f} C_{f}$ and $D_{g} \xrightarrow{g} C_{g}$. Then

$$
f \preceq g \text { if and only if } D_{f} \subseteq D_{g}, C_{f} \subseteq C_{g} \text {, and }\left.g\right|_{D_{f}}=f \text {. }
$$

Let $\mathfrak{C}$ be a chain in $\mathscr{F}$. Let $D_{u} \xrightarrow{u} C_{u}$ be defined as follows:

$$
D_{u}:=\bigcup\left[D_{f}: f \in \mathfrak{C}\right], C_{u}:=\bigcup\left[C_{f}: f \in \mathfrak{C}\right], u(x):=g(x) \text { for } x \in D_{g} .
$$

Note that the assignment $u(x)=g(x)$ for $x \in D_{g}$ is unambiguous because, if there also is $g^{\prime} \in \mathfrak{C}$ such that $x \in D_{g^{\prime}}$ then $g^{\prime}(x)=g(x)$ because $g^{\prime} \preceq g$, i.e., $g$ is an extension of $g^{\prime}$, or $g \preceq g^{\prime}$,i.e., $g^{\prime}$ is an extension of $g$.
$u$ has been constructed in such a fashion that it extends any $g \in \mathfrak{C}$ and hence is an upper bound of $\mathfrak{C}$. It follows that the partially ordered set ( $\mathscr{F}, \preceq$ ) possesses the ZL property (see axiom 12.1 (Zorn's Lemma) on p.259), hence there exists a maximal element $D_{h} \xrightarrow{h} C_{h}$ in $\mathscr{F}$.

We claim that $D_{h}=X$ or $C_{h}=Y$ (or both). Otherwise there would be $x_{0} \in D_{h}^{\complement}$ and $y_{0} \in C_{h}^{\complement}$ and the function

$$
\psi: D_{h} \cup\left\{x_{0}\right\} \longrightarrow C_{h} \cup\left\{y_{0}\right\}, \quad \begin{cases}x \mapsto h(x) & \text { if } x \in D_{h} \\ x \mapsto y_{0} & \text { if } x=x_{0}\end{cases}
$$

is an injective extension of $h$, i.e., $\psi \in \mathscr{F}$. This contradicts the maximality of $h$.
Case 1: $\quad D_{h}=X$. Then $h$ is an injective function which maps $X$ into $Y$, i.e., $\operatorname{card}(X) \leqq \operatorname{card}(Y)$.
Case 2: $\quad C_{h}=Y$. Then $h: D_{h} \rightarrow Y$ is a bijective function whose inverse $f^{-1}: Y \rightarrow D_{h}$ is an injection from $Y$ into the subset $D_{h}$ of $X$, i.e., $\operatorname{card}(Y) \leqq \operatorname{card}(X)$.

### 12.4 Extensions of Linear Functions in Arbitrary Vector Spaces

We now turn our attention to extending a linear real valued function $f$ from a subspace $F \subseteq V$ to the entire vector space $V$.

Lemma 12.2. Let $V$ be a vector space and let $F$ be a (linear) subspace of $V$. Let $f: F \rightarrow \mathbb{R}$ be linear.
Let $\mathscr{G}:=\left\{\left(W, f_{W}\right): W\right.$ is a subspace of $V, W \supseteq F, f_{W}: W \rightarrow \mathbb{R}$ is a linear extension of $f$ to $\left.W\right\}$.
Then the following defines a partial ordering on $\mathscr{G}:\left(U, f_{U}\right) \preceq\left(W, f_{W}\right) \Leftrightarrow V \subseteq W$ and $\left.f_{W}\right|_{U}=f_{U}$.
Moreover this ordering satisfies the requirements of Zorn's Lemma: Every chain in $(\mathscr{G}, \preceq)$ possesses an upper bound (in $\mathscr{G}$ ).

PROOF: Reflexivity and transitivity of " $\preceq$ " are trivial. The latter is true because the extension of an extension is again an extension.

[^95]Antisymmetry: If both $\left(U, f_{U}\right) \preceq\left(W, f_{W}\right)$ and $\left(W, f_{W}\right) \preceq\left(U, f_{U}\right)$ then both $U \subseteq W$ and $W \subseteq U$, hence $U=W$. But then $f_{W}$ is an extension of $f_{U}$ to itself, i.e., $f_{U}=f_{W}$. It follows that $\preceq$ is indeed a partial order on $\mathscr{G}$.

Now let $\mathscr{C}$ be a chain in $\mathscr{G}$. We must find an upper bound for $\mathscr{C}$. Let $W:=\bigcup\left[U:\left(U, f_{U}\right) \in \mathscr{C}\right]$.
We show that $W$ is a subspace of $E$ : If $x, y \in W$ and $\lambda \in \mathbb{R}$ then there are $\left(C_{1}, f_{1}\right),\left(C_{2}, f_{2}\right) \in \mathscr{C}$ such that $x \in C_{1}$ and $y \in C_{2}$. Because $\mathscr{C}$ is a chain we have $C_{1} \subseteq C_{2}$ or $C_{2} \subseteq C_{1}$, say, $C_{1} \subseteq C_{2}$. It follows that $x, y \in C_{2}$. But $C_{2}$ is a subspace of $V$ and we conclude that $x+\lambda y \in C_{2}$, hence $x+\lambda y \in W$. It follows that $W$ is a subspace of $V$.

Let $f_{W}: W \rightarrow \mathbb{R}$ be defined as follows: If $x \in W$ then there is some $\left(C, f_{C}\right) \in \mathscr{C}$ such that $x \in C$. We define $f_{W}(x):=f_{C}(x)$. This definition is unambiguous even if $x$ belongs to (possibly infinitely) many elements of $\mathscr{C}$. To see this let $\left(C, f_{C}\right),\left(D, f_{D}\right) \in \mathscr{C}$ such that $x \in C$ and $x \in D$. Then $C \subseteq D$ or $D \subseteq C$. We may assume that $C \subseteq D$. But as $\left.f_{D}\right|_{C}=f_{C}$ we conclude that $f_{C}(x)=f_{D}(x)$, i.e., the definition of $f_{W}(\cdot)$ is unambiguous. The above specifically holds for $x \in F$ and we note that $f_{W}$ is an extension of $f$.

Next we show linearity of $f_{W}$. Let $x, y \in W$ and $\alpha \in \mathbb{R}$. Then there are $\left(C, f_{C}\right),\left(D, f_{D}\right) \in \mathscr{C}$ such that $x \in C$ and $y \in D$. Again we may assume that $C \subseteq D$. It now follows from the linearity of $f_{D}$ that

$$
f_{W}\left((x+\alpha y)=f_{D}\left((x+\alpha y)=f_{D}\left((x)+\alpha f_{D}(y)=f_{W}\left((x)+\alpha f_{W}(y) .\right.\right.\right.\right.
$$

and we have proved that $f_{W}$ is linear (on all of $W$ ).
To summarize, $W$ is a subspace of $V$ and $f_{W}$ is a linear extension of $f$ to $W$. But then $\left(W, f_{W}\right) \in \mathscr{G}$ and $\left(W, f_{W}\right) \succeq\left(C, f_{C}\right)$ for all $\left(C, f_{C}\right) \in \mathscr{C}$. It follows that $\left(W, f_{W}\right)$ is an upper bound of $\mathscr{C}$.

Lemma 12.3. Let $V$ be a vector space and let $F$ be a (linear) subspace of $V$ such that $F \varsubsetneqq V$. Let $f: F \rightarrow \mathbb{R}$ be linear. Let $a \in F^{\complement}$ and $\gamma \in \mathbb{R}$.

Let $\tilde{F}:=\operatorname{span}(F \uplus\{a\})$. Then $a$ and $\gamma$ define a linear extension $\tilde{f}$ of $f$ to $\tilde{F}$ as follows:

$$
\tilde{f}(x+\alpha a):=f(x)+\alpha \gamma
$$

The proof is left as exercise 12.3.
Theorem 12.3 (Extension theorem for linear real valued functions). Let $V$ be a vector space and let $F$ be a (linear) subspace of $V$. Let $f: F \rightarrow \mathbb{R}$ be a linear mapping.

Then there is an extension of $f$ to a linear mapping $\tilde{f}: V \rightarrow \mathbb{R}$.

## PROOF:

Let $\mathscr{G}:=\left\{\left(W, f_{W}\right): W\right.$ is a subspace of $V, W \supseteq F, f_{W}: W \rightarrow \mathbb{R}$ is a linear extension of $f$ to $\left.W\right\}$
and let $\left(U, f_{U}\right) \preceq\left(W, f_{W}\right) \Leftrightarrow U \subseteq W$ and $\left.f_{W}\right|_{U}=f_{U}$.
We have seen in lemma 12.2 that $\preceq$ is a partial ordering on $\mathscr{G}$ such that any chain in $(\mathscr{G}, \preceq)$ possesses an upper bound. We apply Zorn's Lemma (axiom 12.1 on $p .259$ ) to conclude that $\mathscr{G}$ possesses a maximal element ( $F^{\prime}, f^{\prime}$ ).

We now show that $F^{\prime}=V$.
If this was not true then we could find $a \in V \backslash F^{\prime}$ and, according to lemma 12.3 extend $f^{\prime}$ to a linear function $\tilde{f}$ on $F^{\prime} \uplus\{a\}$. It follows that $\left(F^{\prime} \uplus\{a\}, \tilde{f}\right) \in \mathscr{G}$ and $\left(F^{\prime}, f^{\prime}\right) \supsetneqq\left(F^{\prime} \uplus\{a\}, \tilde{f}\right)$.

This contradicts the maximality of $\left(F^{\prime}, f^{\prime}\right)$. and we have reached a contradiction.

### 12.5 The Hahn-Banach Extension Theorem ( $\star$ )

Note that this chapter is starred, hence optional.
Let $V$ be a vector space and let $F$ be a (linear) subspace of $V$. Let $f: F \rightarrow \mathbb{R}$ be linear mapping. This chapter shows how to extend $f$ from its domain $F$ to the entire space $V$ subject to some condition which guarantees that if $V$ is a normed space (hence a metric space) and if $f$ is continuous on $F$ then this linear extension will be continuous on all of $V$. The Hahn-Banach Extension Theorem provides such conditions.

In preparation for this we must study sublinearity, a generalization of linearity and norms.

### 12.5.1 Sublinear Functionals

Definition 12.3 (Sublinear functionals). Let $V$ be a vector space and $p: V \rightarrow \mathbb{R}$ such that

$$
\begin{array}{lll}
\text { a. if } \lambda \in \mathbb{R}_{\geq 0} \text { and } x \in V \text { then } p(\lambda x)=\lambda p(x) & \text { (positive homogeneity); } \\
\text { b. } & \text { if } x, y \in V \text { then } p(x+y) \leqq p(x)+p(y) & \text { (subadditivity). }
\end{array}
$$

Then we call $p$ a sublinear functional on $V$.
Proposition 12.1. Let $V$ be a vector space and $p: V \rightarrow \mathbb{R}$ sublinear. Let $x \in V$. Then
a. $\quad p(0)=0$,
b. $-p(x) \leqq p(-x)$,

PROOF of $\boldsymbol{a}$ : $p(0)=p(0 \cdot 0)=0 \cdot p(0)=0$.
PROOF of $\boldsymbol{b}$ : This follows from $0=p(0)=p(x+(-x)) \leqq p(x)+p(-x)$.
Example 12.3 (Norms are sublinear). Let $(V,\|, \cdot\|)$ be a normed vector space. Then the function $p(x):=\|x\|$ is sublinear.

Indeed, norms are homogenous: We have $\|\lambda x\|=|\lambda| \cdot\|x\|$ not only for $\lambda \geqq 0$ but for all $\lambda \in \mathbb{R}$.
Further subadditivity is just the validity of the triangle inequality.
Example 12.4 (Linear functions are sublinear). Let $V$ be a vector space and let $f:=V \rightarrow \mathbb{R}$ be a linear function. Then $f$ is sublinear.

Indeed, linear functions $f$ satisfy $f(\lambda x)=\lambda \cdot f(x)$ not only for $\lambda \geqq 0$ but for all $\lambda \in \mathbb{R}$.
Further linear functions satisfy additivity: $f(x+y)=f(x)+f(y)$, hence also subadditivity $f(x+y) \leqq f(x)+f(y)$.

More about sublinearity can be found in chapter 12.6 on p. 266

### 12.5.2 The Hahn-Banach extension theorem and its Proof

This chapter follows closely [5] Choquet.
As mentioned previously, the subject of the Hahn-Banach extension theorem is extension of a continuous, linear, real-valued function from a subspace to the entire vector space. The following remark is about first extending it to "one more dimension".

Remark 12.2. Let $V$ be a vector space, let $F$ be a linear subspace of $V$ and let $f:=F \rightarrow \mathbb{R}$ be a linear function. Let $a \in \underset{\tilde{f}}{V} \backslash F$. Then any linear extension $\tilde{f}$ of $f$ to $\operatorname{span}(F \uplus\{a\})$ is uniquely determined by its value $k:=\tilde{f}(a)$.

Indeed, any $x \in \operatorname{span}(F \uplus\{a\})$ can be written as $u+\lambda a$ for some $u \in F$ and $\lambda \in \mathbb{R}$. It follows from the linearity of $\tilde{f}$ that

$$
\begin{equation*}
\tilde{f}(x+\lambda a)=\tilde{f}(x)+\lambda \tilde{f}(a)=f(x)+\lambda k \tag{12.3}
\end{equation*}
$$

Theorem 12.4 (Hahn-Banach extension theorem). Let $V$ be a vector space and $p: V \rightarrow \mathbb{R}$ a sublinear function. Suppose $F$ is a (linear) subspace of $V$ and $f: F \rightarrow \mathbb{R}$ is a linear mapping with $f \leqq p$ on $F$. Then there is an extension of $f$ to a linear map $\tilde{f}: V \rightarrow \mathbb{R}$ such that $\tilde{f} \leqq p$ on $V$.

Before proving this theorem, first we prove two lemmata.
Lemma 12.4. Suppose $F$ is a subspace of $V, f: F \rightarrow \mathbb{R}$ is a linear mapping and $a \in V \backslash F$. Let $k \in \mathbb{R}$ and

$$
\begin{equation*}
\tilde{f}(x+\lambda a):=f(x)+\lambda k, \quad \text { i.e. }, \quad k=\tilde{f}(a) \tag{12.4}
\end{equation*}
$$

(see (12.3)). Then

$$
\begin{align*}
k \leqq \inf _{u \in F}\{p(u+a)-f(u)\} & \Leftrightarrow \tilde{f}(x+\lambda a) \leqq p(x+\lambda a) \text { for all } \lambda>0 \text { and } x \in F,  \tag{12.5}\\
k \leqq \sup _{v \in F}\{f(v)-p(v-a)\} & \Leftrightarrow \tilde{f}(x+\lambda a) \leqq p(x+\lambda a) \text { for all } \lambda<0 \text { and } x \in F \tag{12.6}
\end{align*}
$$

Proof of $(12.5), \Rightarrow)$ : Let us assume that $\lambda>0$. Then, on account of the left side of (12.5),
$\tilde{f}(x+\lambda a)=f(x)+\lambda k=\lambda(f(x / \lambda)+k) \leqq \lambda(f(x / \lambda)+(p(x / \lambda+a)-f(x / \lambda)))=\lambda p(x / \lambda+a)$
We use the positive homogeneity of $p: \lambda p(x / \lambda+a)=p(x+\lambda a)$ to obtain $\tilde{f}(x+\lambda a) \leqq p(x+\lambda a)$.
Proof of $(12.6), \Rightarrow)$ : Let us assume that $\lambda<0$. Because of the left side of (12.6) and $\lambda<0$ and positive homogeneity of $p$,

$$
\begin{aligned}
k \geqq f(v)-p(v-a) & \Rightarrow \lambda k \leqq f(\lambda v)-\lambda p(v-a) \\
& \Rightarrow-f(\lambda v)+\lambda k \leqq(-\lambda) p(v-a)=p((-\lambda)(v-a))=p((-\lambda) v+\lambda a) .
\end{aligned}
$$

We substitute $v:=x / \lambda \in F$ :

$$
-f(x)+\lambda k \leqq p(-x+\lambda a), \text { hence } \tilde{f}(-x+\lambda a)=f(-x)+\lambda k \leqq p(-x+\lambda a)
$$

We can switch from $-x$ to $x$ as the above holds for all $x$ in the subspace $F$ and because $-x \in F$ if and only if $x \in F$. It follows that $p$ indeed dominates $\tilde{f}$ for all $x \in F$ and $\lambda<0$.

Proof of $(12.5), \Leftarrow)$ : we assume $\tilde{f}(x+\lambda a) \leqq p(x+\lambda a)$ for all $\lambda>0$ and $x \in F$. We now show that $k=\tilde{f}(a) \leqq p(u+a)-f(u)$ for all $u \in F$.

$$
p(u+a)-f(u) \geqq \tilde{f}(u+a)-f(u)=\tilde{f}(u)+\tilde{f}(a)-f(u)=f(u)+\tilde{f}(a)-f(u)=\tilde{f}(a)=k
$$

Proof of $(12.6), \Leftarrow)$ : we assume $\tilde{f}(x+\lambda a) \leqq p(x+\lambda a)$ for all $\lambda<0$ and $x \in F$. We now show that $k=\tilde{f}(a) \geqq f(v)-p(v-a)$ for all $v \in F$.

$$
-p(v-a)+f(v) \leqq-\tilde{f}(v-a)+f(v)=\tilde{f}(a-v)+f(v)=\tilde{f}(a)-\tilde{f}(v)+f(v)=\tilde{f}(a)=k
$$

Lemma 12.5. Let $F \subset V$ be a genuine subspace of $V$ and $a \in V \backslash F$. Let $G:=\operatorname{span}(F \uplus\{a\})$ be the subspace of all linear combinations that can be created by a and or vectors in $F$. Then
a. there exists a linear extension $\tilde{f}$ of $f$ to $G$.
b. This extension is unique if and only if $\sup _{v \in E}\{f(v)-p(v-a)\}=\inf _{u \in E}\{p(u+a)-f(u)\}$.

Proof of $a$. For $u, v \in F$ we have

$$
f(u)+f(v)=f(u+v) \leqq p(u+v)=p((u+a)+(v-a)) \leqq p(u+a)+p(v-a)
$$

and hence $f(v)-p(v-a) \leqq p(u+a)-f(u)$. Therefore

$$
\sup _{v \in F}\{f(v)-p(v-a)\} \leqq \inf _{u \in F}\{p(u+a)-f(u)\} .
$$

Now for a fixed $k \in \mathbb{R}$, we define $\tilde{f}(x+\lambda a)=f(x)+\lambda k$. We claim that $\tilde{f} \leqq p$ if and only if we have

$$
\begin{equation*}
\sup _{v \in F}\{f(v)-p(v-a)\} \leqq k \leqq \inf _{u \in F}\{p(u+a)-f(u)\} \tag{12.7}
\end{equation*}
$$

which will conclude the proof of a since such a $k$ exists by the above work. Our claim holds because $f(x)+$ $\lambda k=\tilde{f}(x+\lambda a) \leqq p(x+\lambda a)$ for all $\lambda$ if and only if

$$
\begin{aligned}
k & \leqq p(u+a)-f(u) \\
\text { and } & \text { for all } u \in F \\
\text { and } & \geqq f(v)-p(v-a)
\end{aligned} \text { for all } v \in F
$$

(the cases $\lambda>0$ and $\lambda<0$ respectively). This is proved above in lemma 12.4.
Proof of b. From (12.7) we deduce that $k$ is unique if and only if $\sup _{v \in E}\{f(v)-p(v-a)\}=$ $\inf _{u \in E}\{p(u+a)-f(u)\}$. This exactly the case where $\tilde{f}$ which was defined by $\tilde{f}(x+\lambda a)=f(x)+\lambda k$ is unique.
PROOF of thm. 12.4 (Hahn-Banach):
Let $\mathscr{G}=\{(V, g): V$ is a subspace of $E, g: V \rightarrow \mathbb{R}$ is linear and $g \leqq p$ on $V\}$.
We define a partial order " $\preceq$ " on $\mathscr{G}$ as follows:

$$
\begin{equation*}
\left(V_{1}, g_{1}\right) \preceq\left(V_{2}, g_{2}\right) \Leftrightarrow V_{1} \subseteq V_{2} \text { and } g_{2} \text { is an extension of } g_{1} \text {. } \tag{12.8}
\end{equation*}
$$

a. We first prove that any chain $\mathscr{C} \subseteq(\mathscr{G}, \preceq)$ has an upper bound: Let $W:=\bigcup[V: V \in \mathscr{C}]$. Then $W$ is a subspace of $E$ because if $x, y \in W$ and $\lambda \in \mathbb{R}$ then there are $\left(V_{1}, g_{1}\right),\left(V_{2}, g_{2}\right) \in \mathscr{C}$ such that $x \in V_{1}$ and $y \in V_{2}$. Because $\mathscr{C}$ is a chain we have $V_{1} \subseteq V_{2}$ or $V_{2} \subseteq V_{1}$, say, $V_{1} \subseteq V_{2}$. It follows that $x, y \in V_{2}$. But $V_{2}$ is a subspace of $E$ and we conclude that $x+\lambda y \in V_{2}$, hence $x+\lambda y \in W$. It follows that $W$ is a subspace of $E$.

We now find a linear $h: W \rightarrow \mathbb{R}$ such that $h \leqq p$ on $W$ and $\left.h\right|_{V}=g$ for all $(V, g) \in \mathscr{C}$, i.e., $h$ is a linear extension of $g$ for all $(V, g) \in \mathscr{C}$. If we find such $h$ then it follows that $(W, h) \in \mathscr{G}$ and $(W, h)$ is an upper bound of $\mathscr{C}$. Let $x \in W$. Then $x \in V_{1}$ for some $\left(V_{1}, g_{1}\right) \in \mathscr{C}$. We define $h(x):=g_{1}(x)$. This assignment is unambiguous because if $x \in V_{2}$ for some other $\left(V_{2}, g_{2}\right) \in \mathscr{C}$ then one of them, say $V_{1}$, is contained in the other and $g_{2}$ is an extension of $g_{1}$, i.e., $h(x)=g_{1}(x)=g_{2}(x)$. As $\left(V_{1}, g_{1}\right) \in \mathscr{G}$ we conclude that $h(x)=g_{1}(x) \leqq p(x)$, i.e., $h \leqq p$ on $W$.

Next we show that $h$ is linear. Let $x, y \in W$ and $\lambda \in \mathbb{R}$. We repeat the argument from the proof that $W$ is a subspace of $V$ to conclude that both $x, y$ belong to some $(V, g) \in \mathscr{C}$. We obtain

$$
h(x+\lambda y)=g(x+\lambda y)=g(x)+\lambda g(y)=h(x)+\lambda h(y)
$$

This completes the proof that $(W, h) \in \mathscr{G}$. Let $(V, g) \in \mathscr{C}$. Clearly, $V \subseteq W=\bigcup[U: U \in \mathscr{C}]$. Further $h$ is linear, dominated by $p$ and is constructed in such a manner that $h(x)=g(x)$ for all $x \in V$. It follows that $(W, h) \succeq(V, g)$ for all $(V, g) \in \mathscr{C}$ and we have proved that $\mathscr{C}$ has an upper bound in $(\mathscr{G}, \preceq)$.
b. We are now in a position to apply Zorn's Lemma (axiom 12.1 on $p .259$ ) to conclude that $\mathscr{G}$ possesses a maximal element $\left(F^{\prime}, f^{\prime}\right)$.

We now show that $F^{\prime}=E$. If this was not true then we could find $a \in E \backslash F^{\prime}$ and, according to lemma 12.5 on $p .265$, extend $f^{\prime}$ to a linear function $\tilde{f}$ on $F^{\prime} \uplus\{a\}$ in such a fashion that $\tilde{f} \leqq p$. It follows that $\left(F^{\prime} \uplus\{a\}, \tilde{f}\right) \in \mathscr{G}$ and $\left(F^{\prime}, f^{\prime}\right) \supsetneqq\left(F^{\prime} \uplus\{a\}, \tilde{f}\right)$. This contradicts the maximality of $\left(F^{\prime}, f^{\prime}\right)$ and we have reached a contradiction.

Corollary 12.1 (Continuous extensions of continuous linear functions). Let $(V,\|\cdot\|)$ be a normed vector space. Let $F$ be a (linear) subspace of $V$ and let $f: F \rightarrow \mathbb{R}$ be a continuous linear mapping on $F$. Then there is an extension of $f$ to a continuous linear map $\tilde{f}: V \rightarrow \mathbb{R}$.

PROOF: Let $p(x):=\|f\| \cdot\|x\|$ (see def.10.35 on $p .221$ ). Because $p$ is a positive multiple of a norm it also is a norm on $V$ (see prop.9.11) on $p .171$ ), hence sublinear by example 12.3 on $p .263$. According to the Hahn-Banach extension theorem there exists a linear extension $\tilde{f}$ of $f$ to all of $V$ such that

$$
\begin{equation*}
\tilde{f}(x) \leqq p(x) \text { for all } x \in V . \tag{12.9}
\end{equation*}
$$

We replace $x$ with $-x$ and obtain from the linearity of $\tilde{f}$ that $-\tilde{f}(x)=\tilde{f}(-x) \leqq p(-x)$.
We note that $p(x)=p(-x)$ because $p$ is a norm. Hence

$$
-\tilde{f}(x) \leqq p(-x)=p(x) \text {, i.e., } \tilde{f}(x) \geqq-p(x) \text {. }
$$

Together with (12.9) this shows that

$$
-p(x) \leqq \tilde{f}(x) \leqq p(x) \text { for all } x \in V
$$

and thus

$$
\begin{equation*}
|\tilde{f}(x)| \leqq p(x)=\|f\| \cdot\|x\| \text { for all } x \in V \tag{12.10}
\end{equation*}
$$

It follows from (12.10) that $f$ has been extended in such a way that $\|\tilde{f}\| \leqq\|f\|$. We aply the continuity criterion for linear functions (thm. 10.15 on $p .222$ ) twice in a row to finish the proof as follows: It follows from the continuity of $f$ that $\|f\|<\infty$. But then $\|\tilde{f}\|<\infty$ and this proves the continuity of $\tilde{f}$.

### 12.6 Convexity ( $\star$ )

Note that this chapter is starred, hence optional.
Definition 12.4 (Concave-up and convex functions). Let $f: \mathbb{R} \rightarrow \mathbb{R}$.
a. The epigraph of $f$ is the set epi $(f):=\left\{\left(x_{1}, x_{2}\right) \in \mathbb{R}^{2}: x_{2} \geqq f\left(x_{1}\right)\right\}$ of all points in the plane that lie above the graph of $f$.
b. $f$ is convex if for any two vectors $\vec{a}, \vec{b} \in e p i(f)$ the entire line segment $S:=\{\lambda \vec{a}+(1-\lambda) \vec{b}\}$ : $0 \leqq \lambda \leqq 1$ is contained in epi $(f)$.
c. Let $f$ be differentiable at all points $x \in \mathbb{R}$. Then $f$ is concave-up, if the function $f^{\prime}: x \mapsto f^{\prime}(x)$ is non-decreasing.

Proposition 12.2 (Convexity criterion). $f$ is convex if and only if the following is true: For any $a \leqq x_{0} \leqq b$ let $S\left(x_{0}\right)$ be the unique number such that the point $\left(x_{0}, S\left(x_{0}\right)\right)$ is on the line segment that connects the points $(a, f(a))$ and $(b, f(b))$. Then

$$
\begin{equation*}
f\left(x_{0}\right) \leqq S\left(x_{0}\right) \tag{12.11}
\end{equation*}
$$

Note that any $x_{0}$ between $a$ and $b$ can be written as $x_{0}=\lambda a+(1-\lambda)$ for some $0 \leqq \lambda \leqq 1$ and that the corresponding $y$-coordinate $S\left(x_{0}\right)=S(\lambda a+(1-\lambda) b$ ) on the line segment that connects $(a, f(a))$ and $(b, f(b))$ then is $S(\lambda a+(1-\lambda) b)=\lambda f(a)+(1-\lambda) f(b)$. Hence we can rephrase the above as follows:
$f$ is convex if and only if for any $a<b$ and $0 \leqq \lambda \leqq 1$ it is true that

$$
\begin{equation*}
f(\lambda a+(1-\lambda) b) \leqq \lambda f(a)+(1-\lambda) f(b) \tag{12.12}
\end{equation*}
$$

PROOF of " $\Rightarrow$ ": Any line segment $S$ that connects the points $(a, f(a))$ and $(b, f(b))$ in such a way that is entirely contained in the epigraph of $f$ will satisfy $f\left(x_{0}\right) \leqq S\left(x_{0}\right)$ for all $a \leqq x_{0} \leqq b$ and it follows that convexity of $f$ implies (12.11).
PROOF of " $\Leftarrow "$ : Now assume that (12.11) is valid. Let $\vec{a}=\left(a_{1}, a_{2}\right), \vec{b}=\left(b_{1}, b_{2}\right) \in$ epi $(f)$, i.e.,

$$
\begin{equation*}
a_{2} \geqq f\left(a_{1}\right) \text { and } b_{2} \geqq f\left(b_{1}\right) . \tag{12.13}
\end{equation*}
$$

We must show that the entire line segment $S:=\{\lambda \vec{a}+(1-\lambda) \vec{b}\}: 0 \leqq \lambda \leqq 1$ is contained in epi $(f)$.
Let $\vec{a}^{\prime}:=\left(a_{1}, f\left(a_{1}\right)\right)$. Let $S^{\prime}:=\left\{\lambda \vec{a}^{\prime}+(1-\lambda) \vec{b}\right\}: 0 \leqq \lambda \leqq 1$ be the line segment obtained by leaving the right endpoint $\vec{b}$ unchanged and pushing the left one downward until $a_{2}$ matches $f\left(a_{1}\right)$. Clearly, $S^{\prime}$ nowhere exceeds $S$.

Let $\vec{b}^{\prime \prime}:=\left(b_{1}, f\left(b_{1}\right)\right)$. Let $S^{\prime \prime}:=\left\{\lambda \vec{a}^{\prime}+(1-\lambda) \overrightarrow{b^{\prime \prime}}\right\}: 0 \leqq \lambda \leqq 1$ be the line segment obtained by leaving the left endpoint $\vec{a}^{\prime}$ unchanged and pushing the right one downward until the $b_{2}$ matches $f\left(b_{1}\right)$. Clearly, $S^{\prime \prime}$ nowhere exceeds $S^{\prime \prime}$.

We view any line segment $T$ between two points with abscissas $a_{1}$ and $b_{1}$ as a function $T(\cdot):\left[a_{1}, b_{1}\right] \rightarrow \mathbb{R}$ which assigns to $x \in\left[a_{1}, b_{1}\right]$ that unique value $T(x)$ for which the point $(x, T(x)$ lies on $T$.

The segment $S^{\prime \prime}$ connects the points $(a, f(a))$ and $(b, f(b))$ and it follows from assumption $\boldsymbol{b}^{\prime}$ that for any $a \leqq x_{0} \leqq b$ we have $f\left(x_{0}\right) \leqq S^{\prime \prime}\left(x_{0}\right)$. We conclude from $S(\cdot) \geqq S^{\prime}(\cdot) \geqq S^{\prime \prime}(\cdot)$ that $S\left(x_{0}\right) \geqq f\left(x_{0}\right)$, i.e. $\left(x_{0}, S\left(x_{0}\right)\right) \in$ epi $(f)$. As this is true for any $a \leqq x_{0} \leqq b$ it follows that the line segment $S$ is entirely contained in the epigraph of $f$.

Example 12.5 (Sublinear functions are convex). Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be sublinear. Then $f$ is convex.
Let $a, b \in \mathbb{R}$ and $0 \leqq \lambda \leqq 1$. It follows from subadditivity and positive homogeneity of $f$ that

$$
f(\lambda a+(1-\lambda) b) \leqq f(\lambda a)+f((1-\lambda) b))=\lambda f(a)+(1-\lambda) f(b) .
$$

According to prop.12.2 this implies convexity of $f$.
Proposition 12.3 (Convex vs concave-up). Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be concave-up. Then $f$ is convex.
PROOF: Assume to the contrary that $f$ is (differentiable and) concave-up and that $f\left(x_{0}\right)>S\left(x_{0}\right)$ for some $a \leqq x_{0} \leqq b$.

It follows that $a<x_{0}<b$ because $f(a)=S(a)$ and $f(b)=S(b)$.

Let $S: x \mapsto S(x)$ be the line through the points $(a, f(a))$ and $(b, f(b))$. and let $m$ be the slope of $S$, i.e.,

$$
m=\frac{S(b)-S(a)}{b-a}
$$

It then follows that

$$
\begin{equation*}
m=\frac{S(b)-S\left(x_{0}\right)}{b-x_{0}}>\frac{S(b)-f\left(x_{0}\right)}{b-x_{0}}=\frac{f(b)-f\left(x_{0}\right)}{b-x_{0}}=f^{\prime}(\xi) \tag{12.14}
\end{equation*}
$$

for some $x_{0}<\xi<b$ (according to the mean value theorem for derivatives). Further

$$
\begin{equation*}
m=\frac{S\left(x_{0}\right)-S(a)}{x_{0}-a}<\frac{f\left(x_{0}\right)-S(a)}{x_{0}-a}=\frac{f\left(x_{0}\right)-f(a)}{x_{0}-a}=f^{\prime}(\eta) \tag{12.15}
\end{equation*}
$$

for some $a<\eta<x_{0}$ (according to the mean value theorem for derivatives).
Because $f$ is concave up we have

$$
f^{\prime}(a) \leqq f^{\prime}(\eta) \leqq f^{\prime}\left(x_{0}\right) \leqq f^{\prime}(\xi) \leqq f^{\prime}(b)
$$

From (12.14) and (12.15) we obtain

$$
m<f^{\prime}(\eta) \leqq f^{\prime}\left(x_{0}\right) \leqq f^{\prime}(\xi)<m
$$

and we have reached a contradiction.
Proposition 12.4 (Sublinear functions are convex). Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be sublinear. Then $f$ is concave-up.
PROOF: Let $0 \leqq \lambda \leqq 1$ and $x, y \in \mathbb{R}$. Then

$$
\begin{equation*}
p(\lambda x+(1-\lambda) y) \leqq p(\lambda x)+p((1-\lambda) y)=\lambda p(x)+(1-\lambda) p(y) \tag{12.16}
\end{equation*}
$$

and it follows that $f$ is concave-up.

### 12.7 Exercises for Ch .12

Exercise 12.1. It was stated in example 12.1 on p .259 that $(\mathbb{R}, \geq)$ is a linearly ordered set. Prove it. (Prove also that it is a POset)
Exercise 12.2. Let $A \subseteq \mathbb{R}^{2}$ and $\vec{x}_{1}=\left(x_{1}, y_{1}\right), \vec{x}_{2}=\left(x_{2}, y_{2}\right) \in A$. Let

$$
\vec{x}_{1} \preceq \vec{x}_{2} \Leftrightarrow x_{1} \leqq x_{2} \text { and } y_{1} \leqq y_{2} .
$$

a. Prove that " $\preceq$ " defines a partial order on $A$.
b. Prove that no maximal elements exist if $A=\mathbb{R}^{2}$.
c. What are the maximal elements of $A=\left\{\vec{x} \in \mathbb{R}^{2}:\|\vec{x}\|_{2}=1\right\}$ ?

Exercise 12.3. Prove lemma 12.3 on p .262 : Let $V$ be a vector space and let $F$ be a (linear) subspace of $V$ such that $F \varsubsetneqq V$. Let $f: F \rightarrow \mathbb{R}$ be linear. Let $a \in F^{\complement}$ and $\gamma \in \mathbb{R}$.

Let $\tilde{F}:=\operatorname{span}(F \uplus\{a\})$. Then $a$ and $\gamma$ define a linear extension $\tilde{f}$ of $f$ to $\tilde{F}$ as follows:

$$
\tilde{f}(x+\alpha a):=f(x)+\alpha \gamma
$$

## 13 Algebraic Structures ( $\star$ )

This chapter is at its very beginnings. It has been created because it is mentioned in one of the first lectures that the axiomatically defined set $\mathbb{Z}$ of the first chapter of [1] B/G forms a group.

Note that this chapter is starred and hence optional.

### 13.1 Semigroups and Groups ( $\star$ )

Definition 13.1 (Semigroups and monoids). Given is a nonempty set $S$ with a binary operation $\diamond$,
i.e. an "assignment rule" $(s, t) \mapsto s \diamond t$ which assigns to any two elements $s, t \in S$ a third element $u:=s \diamond t \in S .{ }^{131}(G, \diamond)$ is called a semigroup if the operation $\diamond$ satisfies
associativity: $(s \diamond t) \diamond u=s \diamond(t \diamond u)$ for all $s, t, u \in S$.
A semigroup for which there exists in addition a neutral element with respect to the operation $(s, t) \mapsto s \diamond t$, i.e., some $e \in S$ such that

$$
\begin{equation*}
s \diamond e=e \diamond s \text { for all } s \in S \tag{13.2}
\end{equation*}
$$

is called a monoid.
Example 13.1. $(\mathbb{Z},+)$ (the integers with addition) and $(\mathbb{Z}, \cdot)$ (the integers with multiplication) are monoids: Both + and $\cdot$ are associative and addition has zero, multiplication has 1 as neutral element.

This is also true for the other number systems: $(\mathbb{N},+)$ and $(\mathbb{N}, \cdot)$ (natural numbers), $(\mathbb{Q},+)$ and $(\mathbb{Q}, \cdot)$ (rational numbers), $(\mathbb{R},+)$ and $(\mathbb{R}, \cdot)$ (real numbers), $(\mathbb{C},+)$ and $(\mathbb{C}, \cdot)$ (complex numbers) all are monoids.

Example 13.2. You need to know from linear algebra or ch.ssec-general-vector-spaces on p. 159 about general vector spaces to understand this example:

If $V$ is a vector space with addition " + " and scalar multiplication "." then $(V,+)$ is a monoid but $(V, \cdot)$ is not. (Why not?)

Example 13.3. You need to know function composition to understand this example.
Let $A$ be a nonempty set and let
$S:=\{f: f$ is a function $A \rightarrow A\} .{ }^{132}$
Then the operation $(f, g) \mapsto g \circ f$ which assigns to any two functions $f$ and $g$ the function ${ }^{133}$ $x \mapsto g \circ f(x):=g(f(x))$ is associative, i.e., $S$ is a semigroup. Moreover $S$ is a monoid because the identity function $i d_{A}: x \mapsto x$ for all $x \in A$ which does nothing with its arguments satisfies (13.2).

PROOF of associativity: For any three functions $f, g, h \in S$ and any $x \in A$ it is true that

$$
((h \circ g) \circ f)(x)=(h \circ g)(f(x))=h(g(f(x))=h((g \circ f)(x))=h \circ(g \circ f)(x)
$$

[^96]. In other words, both the left side $x \mapsto((h \circ g) \circ f)(x)$ and the right side $h \circ(g \circ f)(x)$ are, for each argument $x \in A$, equal to $h(g(f(x))$. This then means that those two functions coincide and we have proved associativity.

Proof that $i d_{A}$ is a neutral element: We have

$$
\left(i d_{A} \circ f\right)(x)=i d_{A}(f(x))=f(x)=f\left(i d_{A}(x)\right)=\left(f \circ i d_{A}\right)(x)
$$

for all $x \in A$. It follows that the three assignments $x \mapsto\left(i d_{A} \circ f\right)(x), x\left(f \circ i d_{A}\right)(x)$ and $x \mapsto f(x)$ coincide for all $x$, i.e., they all represent the same function $x \mapsto f(x)$. This proves (13.2) and we have proved the existence of a neutral element.

Theorem 13.1 (Uniqueness of the neutral element in monoids). Let ( $S, \diamond$ ) be a monoid and let $e, e^{\prime} \in S$ such that both

$$
\begin{gather*}
s \diamond e=e \diamond s=s  \tag{13.3}\\
s \diamond e^{\prime}=e^{\prime} \diamond s=s \tag{13.4}
\end{gather*}
$$

for all $s \in S$. Then $e=e^{\prime}$.
PROOF: We have

$$
e \stackrel{(13.4)}{=} e^{\prime} \diamond e \stackrel{(13.3)}{=} e^{\prime} .
$$

Here we applied (13.4) with $s=e$ and then (13.3) with $s=e^{\prime}$.
Definition 13.2 (Groups and Abelian groups). Let $(G, \diamond)$ be a monoid with neutral element $e$ which satisfies the following: For each $g \in G$ there exists some $g^{\prime} \in G$ such that

$$
\begin{equation*}
g \diamond g^{\prime}=g^{\prime} \diamond g=e \text { for all } g \in G . \tag{13.5}
\end{equation*}
$$

Then $(G, \diamond)$ is called a group.
Assume moreover that the operation $\diamond$ satisfies
commutativity: $g \diamond h=h \diamond g$ for all $g, h \in G$.
Then $G$ is called a commutative group or abelian group. ${ }^{134}$
Theorem 13.2 (Uniqueness of the inverse in groups). Let $(G, \diamond)$ be a group and let $g \in G$. Assume that there exists besides $g^{\prime}$ another $g^{\prime \prime} \in G$ which satisfies (13.5). Then $g^{\prime \prime}=g^{\prime}$.

PROOF: We have

$$
\begin{equation*}
g^{\prime \prime} \stackrel{(13.2)}{=} e \diamond g^{\prime \prime} \stackrel{(13.5)}{=}\left(g^{\prime} \diamond g\right) \diamond g^{\prime \prime} \stackrel{a s s o c}{=} g^{\prime} \diamond\left(g \diamond g^{\prime \prime}\right) \stackrel{(13.5)}{=} g^{\prime} \diamond e \stackrel{(13.2)}{=} g^{\prime} \tag{13.7}
\end{equation*}
$$

and this proves uniqueness.
Definition 13.3 (inverse element $g^{-1}$ ). From now on we are allowed to write $g^{-1}$ for the unique element of $G$ that is associated with the given $g \in G$ by means of the formula (13.5). We call $g^{-1}$ the inverse element of $g$.

[^97]Example 13.4. a. $(\mathbb{Z},+)$ (the integers with addition) is an abelian group: We have already seen that $(\mathbb{Z},+)$ is a monoid.

Inverse element to $k \in \mathbb{Z}$ is $-k$ because $k+(-k)=(-k)+k=0$ for all $k$.
and this group is abelian because $m+k=k+m$ for all $k, m \in \mathbb{Z}$.
and $(\mathbb{Z}, \cdot)$ (the integers with multiplication) are monoids: Both + and $\cdot$ are associative and addition has zero, multiplication has 1 as neutral element.
b. $(\mathbb{Z}, \cdot)$ (the integers with multiplication) is not a group: Let $k=5$. Then $k \in \mathbb{Z}$ but $1 / 5$, the only number $m$ such that $5 \cdot m=m \cdot 5=1$ ( 1 is the neutral element with respect to ".") is a fraction and does not belong to $\mathbb{Z}$.

Example 13.5. Being a group is a lot more specific than just being a semigroup or monoid. Not all types of numbers form groups for addition and/or multiplication:

Natural numbers: $\quad$ Neither $(\mathbb{N},+)$ nor $(\mathbb{N}, \cdot)$ are groups: $(\mathbb{N},+)$ does not even have a neutral element, $(\mathbb{N}, \cdot)$ has 1 as a neutral element but there is no multiplicative inverse for, say, 5 because $1 / 5 \notin \mathbb{N}$.
Integers: $\quad$ We have seen in example 13.4 that $(\mathbb{Z},+)$ is an abelian group but $(\mathbb{Z},+)$ is not a group.
Rational numbers: $\quad(\mathbb{Q},+)$ is an abelian group but $(\mathbb{Q}, \cdot)$ is not a group because the number 0 does not have a multiplicative inverse: There is no number $x$ such that $0 \cdot x=1$. But note that the set $\mathbb{Q}^{\star}$ of all non-zero rational numbers is an abelian group.
Real numbers: $\quad(\mathbb{R},+)$ is an abelian group but $(\mathbb{R}, \cdot)$ is not a group for the same reason as $(\mathbb{Q}, \cdot)$. Again, the set $\mathbb{R}^{\star}$ of all non-zero real numbers is an abelian group.
Complex numbers: $\quad(\mathbb{C},+)$ is an abelian group but $(\mathbb{C}, \cdot)$ is not a group for the same reason as $(\mathbb{Q}, \cdot)$. However the set $\mathbb{C}^{\star}$ of all non-zero complex numbers is an abelian group.

Example 13.6. As in example 13.3 (function composition) let $A$ be a nonempty set and let $S:=\{f$ : $f$ is a function $A \rightarrow A\}$ with the operation $(f, g) \mapsto g \circ f$ defined as $g \circ f(x)=g(f(x))$. We have seen that $S$ is a monoid but $S$ is not a group because not every $f \in S$ has an inverse function $f^{-1}$ which satisfies $f \circ f^{-1}(x)=f^{-1}(f(x))=x$ for all $x \in A$.

Counterexample: Let $A:=\mathbb{R}$. Some functions will have an inverse. For example $f(x)=x-7$ has inverse $f^{-1}(x)=x+7$.

But $f(x)=x^{2}$ does not have an inverse: $g(x):=\sqrt{x}$ will not work because $g(f(-2))=g(4)=2$, not -2 as required, and $h(x):=-\sqrt{x}$ will not work because $h(f(2))=h(4)=-2$, not 2 as required.

Example 13.7. Let $(G, \diamond)$ and $(H, \bullet)$ be defined as follows:

$$
\begin{align*}
G & :=\left\{g \in \mathbb{R}: g=e^{x} \text { for some } x \in \mathbb{R}\right\}, \quad e^{x} \diamond e^{y}:=e^{x} \cdot e^{y}=e^{x}+y  \tag{13.8}\\
H & :=\{h \in \mathbb{R}: h=\ln (x) \text { for some } u \in] 0, \infty\}, \quad \ln u \bullet \ln v:=\ln u+\ln v=\ln (x y) . \tag{13.9}
\end{align*}
$$

a. Both $(G, \diamond)$ and $(H, \bullet)$ are abelian groups $G$ has neutral element 1 and $H$ has neutral element 0. (Exercise: Prove it. What are the inverses?)
b. Let the functions $\varphi$ and $\psi$ be defined as follows:

$$
\begin{array}{ll}
\varphi:(G, \diamond) \rightarrow(H, \bullet), & \varphi(g):=\ln g \\
\psi:(H, \bullet) \rightarrow(G, \diamond), & \psi(h):=e^{h} \tag{13.11}
\end{array}
$$

Then $\varphi$ and $\psi$ satisfy the following:

$$
\begin{align*}
\varphi\left(g_{1} \diamond g_{2}\right) & =\varphi\left(g_{1}\right) \bullet \varphi\left(g_{2}\right), \quad \varphi(1)=0, \quad \varphi\left(g^{-1}\right)=\varphi(g)^{-1}  \tag{13.12}\\
\psi\left(h_{1} \bullet h_{2}\right) & =\psi\left(h_{1}\right) \diamond \psi\left(h_{2}\right), \quad \psi(0)=1, \quad \psi\left(h^{-1}\right)=\psi(h)^{-1} \tag{13.13}
\end{align*}
$$

Further, the functions $\varphi$ and $\psi$ are inverse to each other, i.e.,

$$
\begin{equation*}
\psi(\varphi(g))=g \quad \text { and } \quad \varphi(\psi(h))=h \tag{13.14}
\end{equation*}
$$

for all $g \in G$ and $h \in H$.
If you talk about $\varphi$ and $\psi$ as "the functions" and $\diamond$ and $\bullet$ as "the operations" you might state the results (13.12) and (13.12) as follows:

The functions are structure compatible with the operations on their domains and codomains:

> It does not matter whether you first apply the operation to two items in the domain and then apply the function to the result or whether you first map those two items into the codomain and then apply the operation to the two function values. Further, the inverse of the function value is the function value of the inverse and the function maps the neutral element to the neutral element.

Such functions between two groups have a special name:
Definition 13.4 (Homomorphisms and isomorphisms). Let $(G, \diamond)$ and $(H, \bullet)$ be two groups with neutral elements $e_{G}$ and $e_{H}$ and let us write $g^{-1}$ and $h^{-1}$ for the inverses (in the sense of def. 13.3.

Let $\varphi:(G, \diamond) \rightarrow(H, \bullet)$ be a function which satisfies the following:

$$
\begin{equation*}
\varphi\left(g_{1} \diamond g_{2}\right)=\varphi\left(g_{1}\right) \bullet \varphi\left(g_{2}\right), \quad \varphi\left(e_{G}\right)=e_{H}, \quad \varphi\left(g^{-1}\right)=\varphi(g)^{-1} \tag{13.15}
\end{equation*}
$$

Then we call $\varphi$ a homomorphism, more specifically, a group homomorphism, from the group $(G, \diamond)$ to the group $(H, \bullet)$.

Let $\psi:(H, \bullet) \rightarrow(G, \diamond)$ be a group homomorphism from $(H, \bullet)$ to $(G, \diamond)$ such that $\varphi$ and $\psi$ are inverse to each other. ${ }^{135}$

We call such a bijective homomorphism an isomorphism.

[^98]
## 14 Construction of the Number Systems

### 14.1 The Peano Axioms ( $\star$ )

## Note that this chapter is starred, hence optional.

Definition 14.1 (Set of non-negative integers). We define the set $\mathbb{N}_{0}$ (the non-negative integers) axiomatically as follows:
Ax. $1 \quad$ There is an element " 0 " contained in $\mathbb{N}_{0}$.
Ax. $2 \quad$ There is a function $\sigma: \mathbb{N}_{0} \rightarrow \mathbb{N}_{0}$ such that
Ax.2.1 $\sigma$ is injective,
Ax.2.2 $0 \notin \sigma\left(\mathbb{N}_{0}\right)$ (range of $f$ ),
Ax.2.3 Induction axiom: Let $U \subseteq \sigma\left(\mathbb{N}_{0}\right)$ such that a. $0 \in U$, $\mathbf{b}$. If $n \in U$ then $\sigma(n) \in U$. Then $U=\mathbb{N}_{0}$.
We define $\mathbb{N}:=\mathbb{N}_{0} \backslash\{0\}$.
Definition 14.2 (Iterative function composition). Let $X \neq \emptyset$ and $f: X \rightarrow X$. We now use the induction axiom above to define $f^{n}$ for an arbitrary function $f: X \rightarrow X$.
a. $f^{0}:=\operatorname{id}_{X}: x \mapsto x, \mathbf{b} . f^{1}:=f, \mathbf{c} . f^{2}:=f \circ f$ (function composition), c. $f^{\sigma(n)}:=f \circ f^{n}$.

Proposition 14.1. $f^{n}$ is defined for all $n \in \mathbb{N}_{0}$.
PROOF: Let $U:=\left\{k \in \mathbb{N}_{0}: f^{k}\right.$ is defined $\}$. Then $0 \in U$ as $f^{0}=i d_{A}$ and if $k \in U$, i.e., $f^{k}$ is defined then $f^{\sigma(k)}=f \circ f^{k}$ also is defined, i.e., $\sigma(k) \in U$. It follows from $A x .2 .3$ that $U=\mathbb{N}_{0}$.
Remark $14.1(\sigma(\cdot)$ as successor function). Of course the meaning of $\sigma(n)$ will be that of $n+1$ :

$$
0 \stackrel{\sigma}{\mapsto} 1 \stackrel{\sigma}{\mapsto} 2 \stackrel{\sigma}{\mapsto} 3 \stackrel{\sigma}{\mapsto} \ldots
$$

Definition 14.3 (Addition and multiplication on $\mathbb{N}_{0}$ ). Let $m, n \in \mathbb{N}_{0}$. Let

$$
\begin{align*}
m+n & :=\sigma^{n}(m),  \tag{14.1}\\
m \cdot n & :=\left(\sigma^{m}\right)^{n}(0) \tag{14.2}
\end{align*}
$$

Note that we know the meaning of $\left(\sigma^{m}\right)^{n}: f:=\sigma^{m}$ is a function $A \rightarrow A$ and we have established in prop.14.1 the meaning of $f^{n}$, i.e., $\left(\sigma^{m}\right)^{n}$.
Proposition 14.2. Addition and multiplication satisfy all rules of arithmetic we learned in high school such as

Here 1 is defined as $1=\sigma(0)$.

PROOF: Drudge work.
Definition 14.4 (Order relation $m<n$ on $\mathbb{N}_{0}$ ). Let $m, n \in \mathbb{N}_{0}$.
a. We say $m$ is less than $n$ and we write $m<n$ if there exists $x \in \mathbb{N}$ such that $n=m+x$.
b. We say $m$ is less or equal than $n$ and we write $m \leqq n$ if $m<n$ or $m=n$.
c. We say $m$ is greater than $n$ and we write $m>n$ if $n<m$. We say $m$ is greater or equal than $n$ and we write $m \geqq n$ if $n \leqq m$.

Proposition 14.3. " $<$ " and " $\leqq$ " satisfy all the usual rules we learned in high school such as
Trichotomy of the order relation: Let $m, n \in \mathbb{N}_{0}$. Then exactly one of the following is true:

$$
m<n, \quad m=n, \quad m>n .
$$

PROOF: Drudge work.

### 14.2 Constructing the Integers from $\mathbb{N}_{0}$

For the following look at B/G project 6.9 in ch.6.1 and B/G prop.6.25 in ch.6.3.
Definition 14.5 (Integers as equivalence classes). We define the following equivalence relation $\left(m_{1}, n_{1}\right) \sim\left(m_{2}, n_{2}\right)$ on the cartesian product $\mathbb{N}_{0} \times \mathbb{N}_{0}$ :

$$
\begin{equation*}
\left(m_{1}, n_{1}\right) \sim\left(m_{2}, n_{2}\right) \Leftrightarrow m_{1}+n_{2}=n_{1}+m_{2} \tag{14.10}
\end{equation*}
$$

We write $\mathbb{Z}:=\left\{[(m, n)]: m, n \in \mathbb{N}_{0}\right\}$. In other words, $\mathbb{Z}$ is the set of all equivalence classes with respect to the equivalence relation (14.10).

We "embed" $\mathbb{N}_{0}$ into $\mathbb{Z}$ with the following injective function $e: \mathbb{N}_{0} \rightarrow \mathbb{Z}: \mathrm{e}(\mathrm{m}):=[(\mathrm{m}, 0)]$.
From this point forward we do not distinguish between $\mathbb{N}_{0}$ and its image $e\left(\mathbb{N}_{0}\right) \subseteq \mathbb{Z}$ and we do not distinguish between $\mathbb{N}$ and its image $e(\mathbb{N}) \subseteq \mathbb{Z}$. In particular we do not distinguish between the two zeroes 0 and $[(0,0)]$ and between the two ones 1 and $[(1,0)]$.

Finally we write $-n$ for the integer $[(0, n)]$.
With those abbreviation we then obtain
Proposition 14.4 (Trichotomy of the integers). Let $z \in \mathbb{Z}$. Then exactly one of the following is true:
Either $\boldsymbol{a} . ~ z \in \mathbb{N}$, i.e., $z=[(m, 0)]$ for some $m \in \mathbb{N}$ or $\boldsymbol{b} . ~-z \in \mathbb{N}$, i.e., $z=[(0, n)]$ for some $n \in \mathbb{N}$ or c. $z=0$.

PROOF: Drudge work.
Remark 14.2. a. The intuition that guided the above definition is that the pairs $(4,0),(7,3),(130,126)$ all define the same integer 4 and the pairs $(0,4),(3,7),(126,130)$ all define the same integer -4 .
b. If it had been possible to define subtraction $m-n$ for all $m, n \in \mathbb{N}_{0}$ then (14.10) could be rewritten as

$$
\left(m_{1}, n_{1}\right) \sim\left(m_{2}, n_{2}\right) \Leftrightarrow m_{1}-n_{1}=m_{2}-n_{2} .
$$

Looking at the equivalent pairs $(4,0),(7,3),(130,126)$ we get $4-0=7-3=130-126=4$ and for $(0,4),(3,7),(126,130)$ we get $0-4=3-7=126-130=-4$.

Definition 14.6 (Addition, multiplication and subtraction on $\mathbb{Z})$. Let $\left[\left(m_{1}, n_{1}\right)\right]$ and $\left[\left(m_{2}, n_{2}\right)\right] \in \mathbb{Z}$. We define

$$
\begin{align*}
-\left[\left(m_{1}, n_{1}\right)\right] & :=\left[n_{1}, m_{1}\right],  \tag{14.11}\\
{\left[\left(m_{1}, n_{1}\right)\right]+\left[\left(m_{2}, n_{2}\right)\right] } & :=\left[\left(m_{1}+m_{2}, n_{1}+n_{2}\right)\right]  \tag{14.12}\\
{\left[\left(m_{1}, n_{1}\right)\right] \cdot\left[\left(m_{2}, n_{2}\right)\right] } & :=\left[\left(m_{1} m_{2}+n_{1} n_{2}, m_{1} n_{2}+n_{1} m_{2}\right)\right] \tag{14.13}
\end{align*}
$$

We write $\left[\left(m_{1}, n_{1}\right)\right]-\left[\left(m_{2}, n_{2}\right)\right]$ (" $\left[\left(m_{1}, n_{1}\right)\right]$ minus $\left[\left(m_{2}, n_{2}\right)\right]$ ") as an abbreviation for $\left[\left(m_{1}, n_{1}\right)\right]+\left(-\left[\left(m_{2}, n_{2}\right)\right]\right)$.

We write $\left[\left(m_{1}, n_{1}\right)\right]<\left[\left(m_{2}, n_{2}\right)\right]$ if $\left[\left(m_{2}, n_{2}\right)\right]-\left[\left(m_{1}, n_{1}\right)\right] \in \mathbb{N}$, i.e., if there is $k \in \mathbb{N}$ such that $\left[\left(m_{2}, n_{2}\right)\right]-\left[\left(m_{1}, n_{1}\right)\right]=[(k, 0)]$. We then say that $\left[\left(m_{1}, n_{1}\right)\right]$ is less than $\left[\left(m_{2}, n_{2}\right)\right]$.

We write $\left[\left(m_{1}, n_{1}\right)\right] \leqq\left[\left(m_{2}, n_{2}\right)\right]$ if $\left[\left(m_{1}, n_{1}\right)\right]<\left[\left(m_{2}, n_{2}\right)\right]$ or if $\left[\left(m_{1}, n_{1}\right)\right]=\left[\left(m_{2}, n_{2}\right)\right]$ and we then say that $\left[\left(m_{1}, n_{1}\right)\right]$ is less than or equal to $\left[\left(m_{2}, n_{2}\right)\right]$.

We write $\left[\left(m_{1}, n_{1}\right)\right]>\left[\left(m_{2}, n_{2}\right)\right]$ if $\left[\left(m_{2}, n_{2}\right)\right]<\left[\left(m_{1}, n_{1}\right)\right]$ and we then say that $\left[\left(m_{1}, n_{1}\right)\right]$ is greater than $\left[\left(m_{2}, n_{2}\right)\right]$.

We write $\left[\left(m_{1}, n_{1}\right)\right] \geqq\left[\left(m_{2}, n_{2}\right)\right]$ if $\left[\left(m_{2}, n_{2}\right)\right] \leqq\left[\left(m_{1}, n_{1}\right)\right]$ and we then say that $\left[\left(m_{1}, n_{1}\right)\right]$ is greater than or equal to $\left[\left(m_{2}, n_{2}\right)\right]$.

We write $\mathbb{Z}_{\geqq 0}$ for the set of all integers $z$ such that $z \geqq 0$ and $\mathbb{Z}_{\neq 0}$ for the set of all integers $z$ such that $z \neq 0$. You should convince yourself that $\mathbb{Z}_{\geqq 0}=\mathbb{N}_{0}$.

It turns out that all three operations are "well defined" in the sense that the resulting equivalence classes on the right of each of the three equations above do not depend on the choice of representatives in the classes on the left. Further we have

Proposition 14.5. Let $m, n \in \mathbb{N}_{0}$. Then

$$
\begin{array}{r}
{[(m, n)]+[(0,0)]=[(0,0)]+[(m, n)]=[(m, n)],} \\
(-[(m, n)])+[(m, n)]=[(m, n)]+(-[(m, n)])=[0,0] \\
{[(m, n)] \cdot[(1,0)]=[(1,0)] \cdot[(m, n)]=[(m, n)],} \tag{14.16}
\end{array}
$$

i.e., $[(0,0)]$ becomes the neutral element with respect to addition, $[(1,0)]$ becomes the neutral element with respect to multiplication and $-[(m, n)]$ becomes the additive inverse of $[(m, n)]$.

PROOF: Drudge work.
Remark 14.3. Again, if it had been possible to define subtraction $m-n$ for all $m, n \in \mathbb{N}_{0}$ then it would be easier to see why addition and multiplication have been defined as you see it in def.14.6:

Addition is defined such that $\left(m_{1}-n_{1}\right)+\left(m_{2}-n_{2}\right)=\left(m_{1}+m_{2}\right)-\left(n_{1}+n_{2}\right)$ and multiplication: $\left(m_{1}-n_{1}\right) \cdot\left(m_{2}-n_{2}\right)=\left(m_{1} m_{2}+\left(-n_{1}\right)\left(-n_{2}\right)\right)-\left(m_{1} n_{2}+n_{1} m_{2}\right)$.

### 14.3 Constructing the Rational Numbers from $\mathbb{Z}$

For the following look again at $B / G$ project 6.9 in ch.6.1 and $B / G$ prop.6.25 in ch.6.3.
Definition 14.7 (Fractions as equivalence classes). We define the following equivalence relation $(p, q) \sim(r, s)$ on the cartesian product $\mathbb{Z} \times \mathbb{Z}_{\neq 0}$ :

$$
\begin{equation*}
(p, q) \sim(r, s) \Leftrightarrow p \cdot s=q \cdot r \tag{14.17}
\end{equation*}
$$

We write $\mathbb{Q}:=\{[(p, q)]: p, q \in \mathbb{Z}$ and $q \neq 0\}$. In other words, $\mathbb{Q}$ is the set of all equivalence classes with respect to the equivalence relation (14.17).

We "embed" $\mathbb{Z}$ into $\mathbb{Q}$ with the injective function $e: \mathbb{Z} \rightarrow \mathbb{Q}$ defined as $e(z):=[(z, 1)]$.
Remark 14.4. a. The intuition that guided the above definition is that the pairs $(12,4),(-21,-7),(105,35)$ all define the same fraction $3 / 1$ and the pairs $(4,-12),(-7,21),(-35,105$ all define the same fraction $-1 / 3$.
b. If it had been possible to define division $p / q$ for all $p, q \in \mathbb{Z}$ for which $q \neq 0$ then (14.17) could be rewritten as

$$
(p, q) \sim(r, s) \Leftrightarrow p / q=r / s
$$

Looking at the equivalent pairs $(12,4),(-21,-7),(105,35)$ we get $12 / 4=(-21) /(-7)=105 / 35=$ 3 and for $(4,-12),(-7,21),(-35,105)$ we get $4 /(-12)=(-7) / 21=(-35) / 105=-1 / 3$.
c. It is easy to see that $(p, q) \sim(r, s)$ if and only if there is (rational) $\alpha \neq 0$ such that $r=\alpha p$ and $s=\alpha q$. A formal proof is just drudgework.

Definition 14.8 (Addition, multiplication, subtraction and division in $\mathbb{Q}$ ). Let $\left[\left(p_{1}, q_{1}\right)\right]$ and $\left[\left(p_{2}, q_{2}\right)\right] \in \mathbb{Q}$. We define

$$
\begin{align*}
-\left[\left(p_{1}, q_{1}\right)\right] & :=\left[\left(-p_{1}, q_{1}\right)\right],  \tag{14.18}\\
{\left[\left(p_{1}, q_{1}\right)\right]+\left[\left(p_{2}, q_{2}\right)\right] } & :=\left[\left(p_{1} q_{2}+q_{1} p_{2}, n_{1} n_{2}\right)\right]  \tag{14.19}\\
{\left[\left(p_{1}, q_{1}\right)\right]-\left[\left(p_{2}, q_{2}\right)\right] } & :=\left[\left(p_{1}, q_{1}\right)\right]+\left(-\left[\left(p_{2}, q_{2}\right)\right]\right)  \tag{14.20}\\
{\left[\left(p_{1}, q_{1}\right)\right] \cdot\left[\left(p_{2}, q_{2}\right)\right] } & :=\left[\left(p_{1} p_{2}, q_{1} q_{2}\right)\right]  \tag{14.21}\\
{\left[\left(p_{1}, q_{1}\right)\right]^{-1} } & :=[(1,1)] /\left[\left(p_{1}, q_{1}\right)\right]:=\left[\left(q_{1}, p_{1}\right)\right]\left(\text { if } p_{1} \neq 0\right),  \tag{14.22}\\
{\left[\left(p_{1}, q_{1}\right)\right] /\left[\left(p_{2}, q_{2}\right)\right] } & :=\left[\left(p_{1} q_{2}, q_{1} p_{2}\right)\right]=\left[\left(p_{1}, q_{1}\right)\right] \cdot\left[\left(p_{2}, q_{2}\right)\right]^{-1}\left(\text { if } p_{2} \neq 0\right) \tag{14.23}
\end{align*}
$$

It turns out that operations above are "well defined" in the sense that the resulting equivalence classes on the right of each of the three equations above do not depend on the choice of representatives in the classes on the left. 136

Further we have
Proposition 14.6 (Trichotomy of the rationals). Let $x \in \mathbb{Q}$. Then exactly one of the following is true:
Either $\boldsymbol{a}$. $x>0$, i.e., $x=[(p, q)]$ for some $p, q \in \mathbb{N}$ or $\boldsymbol{b} . ~-x>0$, i.e., $x=[(-p, q)]$ for some $p, q \in \mathbb{N}$ or $c$. $x=0$.

PROOF: Drudge work.

### 14.4 Constructing the Real Numbers via Dedekind Cuts

The material presented here, including the notation, follows [10] Rudin, Walter: Principles of Mathematical Analysis.

Note that in this section small greek letters denote sets of rational numbers!
The idea behind real numbers as intervals of rational numbers with no lower bounds, called Dedekind cuts, is as follows:

[^99]Given a real number $x$ you can associate with it the $\operatorname{set}\{q \in \mathbb{Q}: q<x\}$ which we call the cut or Dedekind cut associated with $x$ The mapping

$$
\begin{equation*}
\Phi: x \mapsto \Phi(x):=\{q \in \mathbb{Q}: q<x\} \tag{14.24}
\end{equation*}
$$

is injective because if $x, y \in \mathbb{R}$ such that $x \neq y$, say, $x<y$, then we have $\{q \in \mathbb{Q}: q<x\} \subsetneq\{q \in \mathbb{Q}: q<y\}$ because there are (infinitely many) rational numbers in the open interval $] x, y$ [and we get surjectivity of $\Phi$ for free if we take as codomain the set of all cuts. Because $\Phi$ is bijective we can "identify" any real number with its cut. We now go in reverse: we start with a definition of cuts which does not reference the real number $x$, i.e., we define them just in terms of rational numbers and define addition, multiplication and the other usual operations on those cuts and show that those cuts have all properties of the real numbers as they were axiomatically defined in $B / G$ ch.8, including the completeness axiom which states that each subset $A$ of $\mathbb{R}$ with upper bounds has a least upper bound $\sup (A)$, i.e., a minimum in the set of all its upper bounds.

Definition 14.9 (Dedekind cuts). (Rudin def.1.4)
We call a subset $\alpha \subseteq \mathbb{Q}$ a cut or Dedekind cut if it satisfies the following:
a. $\alpha \neq \emptyset$ and $\alpha^{\complement} \neq \emptyset$
b. Let $p, q \in \mathbb{Q}$ such that $p \in \alpha$ and $q<p$. Then $q \in \alpha$.
c. $\alpha$ does not have a max: $\forall p \in \alpha \exists q \in \alpha$ such that $p<q$.

Given a cut $\alpha$, let $p \in \alpha$ and $q \in \alpha^{\complement}$. We call $p$ a lower number of the cut $\alpha$ and we call $q$ an upper number of $\alpha$.
Theorem 14.1. (Rudin thm.1.5)
Let $\alpha \subseteq \mathbb{Q}$ be a cut. Let $p \in \alpha, q \in \alpha^{\complement}$. Then $p<q$.
Assume to the contrary that $q \leqq p$. Then we either have $p=q$ which means that either both $p, q$ belong to $\alpha$ or both belong to its complement, a contradiction to our assumption. Or we have $q<p$. It then follows from $p \in \alpha$ and def.14.9.b that $q \in \alpha$, contrary to our assumption.
Theorem 14.2. (Rudin thm.1.6)

$$
\text { Let } r \in \mathbb{Q} \text {. Let } r^{\star}:=\{p \in \mathbb{Q}: p<r\} \text {. Then } r^{\star} \text { is a cut and } r=\min \left(\left(r^{\star}\right)^{\complement}\right) \text {. }
$$

PROOF: In the following let $p, q, r \in \mathbb{Q}$.
PROOF of def.14.9.a: $r-1<r \Rightarrow r-1 \in r^{\star} \Rightarrow r^{\star} \neq \emptyset$. Further, $r \in\left(r^{\star}\right)^{\complement} \Rightarrow\left(r^{\star}\right)^{\complement} \neq \emptyset$.
PROOF of def.14.9.b: Let $q<p$ and $p \in r^{\star}$. Then also $q \in r^{\star}=\left\{p^{\prime} \in \mathbb{Q}: p^{\prime}<r\right\}$.
PROOF of def.14.9.c: Let $p \in r^{\star}$. Then $p<(p+r) / 2<r$, hence $(p+r) / 2 \in r^{\star}$ and $r$ cannot be the max of $r^{\star}$.

Definition 14.10 (Rational cuts). Let $r \in \mathbb{Q}$. The cut $r^{\star}=\{p \in \mathbb{Q}: p<r\}$ from the previous theorem is called the rational cut associated with $r$.

Remark 14.5. If we define intervals in $\mathbb{Q}$ in the usual way for $p, q \in \mathbb{Q}$ :

$$
] p, q[:=\{r \in \mathbb{Q}: p<r<q\}, \quad[p, q]:=\{r \in \mathbb{Q}: p \leqq r \leqq q\}, \text { etc. }
$$

then rational cuts $r^{\star}(r \in \mathbb{Q})$ are those for which $\left.r^{\star}=\right]-\infty, r\left[\right.$ and $\left(r^{\star}\right)^{\complement}=[r, \infty[$ whereas for nonrational cuts $\alpha$ we cannot specify the "thingy" that should take the role of $r$. It would be the $\sup (\alpha)$ if we already had defined the set of all real numbers and we could understand $\alpha$ as a subset of those real numbers.

Definition 14.11 (Ordering Dedekind cuts). (Rudin def.1.9) Let $\alpha, \beta$ be two cuts.
We say $\alpha<\beta$ if $\alpha \subsetneq \beta$ (strict subset) and we say $\alpha \leqq \beta$ if $\alpha<\beta$ or $\alpha=\beta$, i.e., $\alpha \subseteq \beta$.
Proposition 14.7 (Trichotomy of the cuts). (Rudin thm.1.10)
Let $\alpha, \beta$ be two cuts. Then either $\alpha<\beta$ or $\alpha>\beta$ or $\alpha=\beta$.
PROOF: We only need to show that if $\alpha \nsubseteq \beta$ then $\beta \subsetneq \alpha$.
So let $\alpha \nsubseteq \beta$. Then $\alpha \backslash \beta$ is not empty and there exists $q \in \alpha \backslash \beta$.
But then $q>b$ for all $b \in \beta$. Also, if $a \in \mathbb{Q}$ and $a \leqq q$ then $a \in \alpha$ (we applied def.14.9.b twice.)
As $b<q$ for all $b \in \beta$ it follows that $\beta \subseteq \alpha$. We saw earlier that $\alpha \backslash \beta \neq \emptyset$ and this proves that $\beta \neq \alpha$, i.e., $\beta \subsetneq \alpha$.

Theorem 14.3 (Addition of two cuts). (Rudin thm.1.12) Let $\alpha, \beta$ be two cuts and let

$$
\alpha+\beta:=\{a+b: a \in \alpha, b \in \beta\} .
$$

Then the set of all cuts is an abelian group with this operation. In other words, $+i$ commutative and associative with a neutral element (which turns out to be $0^{\star}$, the rational cut corresponding to $0 \in \mathbb{Q}$ ) and a suitably defined cut $-\alpha$ for a given cut $\alpha$ which satisfies $\alpha+(-\alpha)=(-\alpha)+\alpha=0^{\star}$

Having defined negatives $-\alpha$ for all cuts we then also can define their absolute values

$$
|\alpha|:= \begin{cases}\alpha & \text { if } \alpha \geqq 0^{\star}, \\ -\alpha & \text { if } \alpha<0^{\star} .\end{cases}
$$

PROOF: Not given here.
Theorem 14.4 (Multiplication of two cuts). Let $\alpha \geqq 0^{\star}, \beta \geqq 0^{\star}$ be two non-negative cuts. Let

$$
\alpha \cdot \beta:= \begin{cases}\{q \in \mathbb{Q}: q<0\} \cup\{a b: a \in \alpha, b \in \beta\} & \text { if } \alpha \geqq 0^{\star}, \beta \geqq 0^{\star}, \\ -|\alpha| \cdot|\beta| & \text { if } \alpha<0^{\star}, \beta \geqq 0^{\star} \text { or } \alpha \geqq 0^{\star}, \beta<0^{\star}, \\ |\alpha| \cdot|\beta| & \text { if } \alpha<0^{\star}, \beta<0^{\star} .\end{cases}
$$

Then the set $\alpha \cdot \beta$ is a cut, called the product of $\alpha$ and $\beta$.
It can be proved that for each cut $\alpha \neq 0^{\star}$ there is a cut $\alpha^{-1}$ uniquely defined by the equation $\alpha \cdot \alpha^{-1}=1^{\star}$.
Theorem 14.5 (The set of all cuts forms a field). Let Let $\mathbb{R}$ be the set of all cuts. Then $\mathbb{R}$ satisfies axioms $8.1-8.5$ of $B / G$ :

Addition and multiplication are both commutative and associative and the law of distributivity $\alpha \cdot(\beta+\gamma)=\alpha \cdot \beta+\alpha \cdot \gamma$ holds.

The cut $0^{\star}$ is the neutral element for addtition and the cut $1^{\star}$ is the neutral element for multiplication. $-\alpha$ is the additive inverse of any cut $\alpha$ and $\alpha^{-1}$ is the multiplicative inverse of $\alpha \neq 0^{\star}$.
Further the set $\mathbb{R}_{>0}:=\left\{\alpha \in \mathbb{R}: \alpha>0^{\star}\right\}$ satisfies $B / G$ axiom 8.26.
PROOF: It follows from prop. 14.7 on $p .278$ that $\mathbb{R}_{>0}$ satisfies $B / G$ axiom 8.26. Proofs of the other properties of $\mathbb{R}$ are not given here.

In the remainder of this section we will see that the completeness axiom B/G ax.8.52 (every subset of $\mathbb{R}$ with upper bounds has a supremum) is a consequence from the properties of cuts and there is no need to state it as an axiom.

Theorem 14.6. (Rudin thm.1.29) Let $\alpha, \beta \in \mathbb{R}$ and let $\alpha<\beta$. Then there exists $q \in \mathbb{Q}$ such that $\alpha<q^{\star}<\beta$
PROOF: Any $q \in \beta \backslash \alpha$ will do.
Theorem 14.7. (Rudin thm.1.30) Let $\alpha \in \mathbb{R}, p \in \mathbb{Q}$. Then $p \in \alpha \Leftrightarrow p^{\star}<\alpha$, i.e., $p^{\star} \subsetneq \alpha$
PROOF of $\Leftarrow)$ : Let $p \in \alpha$. it follows for any $q \in p^{\star}$ that $q<p \in \alpha$, hence $q \in \alpha$, hence $p^{\star} \subseteq \alpha$. As $p \notin p^{\star}=\left\{p^{\prime} \in \mathbb{Q}: p^{\prime}<p\right\}$ but $p \in \alpha$ we have strict inclusion $p^{\star} \subsetneq \alpha$.
PROOF of $\Rightarrow)$ : As $p^{\star} \subsetneq \alpha$ there exists $q \in \alpha \backslash p^{\star}$. As $q \geqq p$ and $q \in \alpha$ we obtain $p \in \alpha$ from def.14.9.b.
Theorem 14.8 (Dedekind's Theorem). (Rudin thm.1.32) Let $\mathbb{R}=A \biguplus B$ a partitioning of $\mathbb{R}$ such that
a. $A \neq \emptyset$ and $B \neq \emptyset$
b. $\alpha \in A, \beta \in B \Rightarrow \alpha<\beta$ (i.e., $\alpha \subsetneq \beta$ ).

Then there exists a unique cut $\gamma \in \mathbb{R}$ such that if $\alpha \in A$ then $\alpha \leqq \gamma$ and if $\beta \in B$ then $\gamma \leqq \beta$.
PROOF: We first prove uniqueness and afterwards the existence of $\gamma$.
PROOF of uniqueness: Assume there is $\gamma^{\prime \prime} \in \mathbb{R}$ which satisfies $\alpha \leqq \gamma^{\prime \prime}$ for all $\alpha \in A$ and $\gamma^{\prime \prime} \leqq \beta$ for all $\beta \in B$.

We may assume that $\gamma<\gamma^{\prime \prime}$. It follows from thm. 14.6 on $p .279$ that there is $\gamma^{\prime} \in \mathbb{R}$ (matter of fact, a rational cut) such that $\gamma<\gamma^{\prime}<\gamma^{\prime \prime}$. But $\gamma<\gamma^{\prime}$ implies that $\gamma^{\prime} \in B$ and $\gamma^{\prime}<\gamma^{\prime \prime}$ implies that $\gamma^{\prime} \in A=B^{\complement}$. We have reached a contradiction and conclude that $\gamma$ must be unique.
PROOF of existence of $\gamma$ : Let $\gamma:=\bigcup[\alpha: \alpha \in A]$.
Step 1: We now show that $\gamma$ is a cut.
We first show that def.14.9. $a$ is satisfied. As $B \neq \emptyset$ there is some $\beta \in B$. As $\beta^{\complement} \neq \emptyset$ there is some $q \in \beta^{\complement}$. It follows from $\alpha \subseteq \beta$ for all $\alpha \subseteq \gamma=\bigcup[\alpha: \alpha \in A]$ that $\gamma \subseteq \beta$, hence $\gamma^{\complement} \supseteq \beta^{\complement}$. It follows from $q \in \beta^{\complement}$ that $q \in \gamma^{\complement}$, hence $\gamma^{\complement} \neq \emptyset$. Further, it follows from $A \neq \emptyset$ that $\gamma \neq \emptyset$. We conclude that def.14.9.a is satisfied.

Next we show the validity of def.14.9.b. Let $p \in \gamma$, i.e., $p \in \alpha_{0}$ for some $\alpha_{0} \in A$. Let $q<p$. Then $q \in \alpha_{0} \subseteq \bigcup[\alpha: \alpha \in A]$, i.e., $q \in \gamma$. We conclude that def.14.9.b is satisfied.

Now we show the validity of def.14.9.c. Let $p \in \gamma$, i.e., $p \in \alpha_{0}$ for some $\alpha_{0} \in A$. As the cut $\alpha_{0}$ does not have a maximum there exists some $q \in \alpha_{0}$ such that $q>p$. As $\alpha_{0} \subseteq \gamma$, hence $q \in \gamma$ We have seen that any $p \in \gamma$ is strictly dominated by some $q \in \gamma$. It follows that $\gamma$ does not have a max and this shows that def.14.9.c is satisfied. We conclude that $\gamma$ is a cut and step 1 of the proof for existence is completed.

Step 2: It remains to show that $\alpha \leqq \gamma \leqq \beta$ for all $\alpha \in A$ and $\beta \in B$. It is trivial that $\alpha \leqq \gamma$ for all $\alpha \in A$ because $\gamma:=\bigcup[\alpha: \alpha \in A]$.

To show that $\gamma \leqq \beta$ for all $\beta \in B$ we prove that the opposite statement that

$$
\begin{equation*}
\gamma>\beta \text {, i.e., } \gamma \backslash \beta \neq \emptyset \text { for some cut } \beta \in B \tag{14.25}
\end{equation*}
$$

will lead to a contradiction. As $q \in \gamma$ there is some $\alpha_{0} \in A$ such that $q \in \alpha_{0}$. Actually, $q \in \alpha_{0} \backslash \beta$ because $q \notin \beta$. But then $\alpha_{0} \nless \beta$ even though $\alpha_{0} \in A$ and $\beta \in B$, contrary to the assumptions about the partitioning $A \biguplus B$ of $\mathbb{R}$.

Corollary 14.1. Let $\mathbb{R}=A \uplus B$ be a partitioning of $\mathbb{R}$ such that
a. $A \neq \emptyset$ and $B \neq \emptyset$
b. $\alpha \in A, \beta \in B \Rightarrow \alpha<\beta$ (i.e., $\alpha \subsetneq \beta$ ).

Then either $\max (A)(=$ l.u.b. $(A))$ exists or $\min (B)(=$ g.l.b. $(B))$ exists.
PROOF: According to thm.14.8 there exists $\gamma \in \mathbb{R}$ such that if $\alpha \in A$ then $\alpha \leqq \gamma$ and if $\beta \in B$ then $\gamma \leqq \beta$. Clearly $\gamma$ is an upper bound of $\alpha$ and a lower bound of $\beta$. It follows that if $\gamma \in A$ then $\max (A)=\gamma$ and if $\gamma \notin A$, i.e., $\gamma \in B$, then $\min (B)=\gamma$.

Theorem 14.9 (Completeness theorem for $\mathbb{R}$ ). (Rudin thm.1.36)
Let $\emptyset \neq E \subset \mathbb{R}$ and assume that $E$ is bounded above. Then $E$ has a least upper bound which we denote by $\sup (E)$ or l.u.b.(E).

PROOF: Let $B$ be the set of all upper bounds for $E$, i.e., $b \in B$ if and only if $b \geqq x$ for all $x \in E$. Then $B$ is not empty by assumption. Let $A:=B^{\complement}=\{\alpha \in \mathbb{R}: \alpha<x$ i.e., $\alpha \subsetneq x$ for some $x \in E\}$. In other words, $\alpha \in A$ if and only if $\alpha$ is not an upper bound of $E$.
$A$ is not empty either: As $E \neq \emptyset$ there is some $x \in E$. Let $\alpha:=x-1$. Cleary $x \leqq \alpha$ is not true for all $x \subseteq E$. It follows that $\alpha$ is not an upper bound of $E$, hence $\alpha \in A$, hence $A$ is not empty.

Moreover we have $\alpha<\beta$ for all $\alpha \in A$ and $\beta \in B$. Because for any $\alpha \in A$ there is some $x \in E$ such that $\alpha<x$ and we have $x \leqq \beta$ for all upper bounds $\beta$, i.e., for all $\beta \in B$.

It follows that the sets $A$ and $B$ form a partition which satisfies the requirements of Dedekind's Theorem (thm.14.8). Hence there exists $\gamma \in \mathbb{R}$ such that $\alpha \leqq \gamma \leqq \beta$ for all $\alpha \in A$ and $\beta \in B$.

We now show that that the assumption $\gamma \in A$ leads to a contradiction. As $\gamma$ is not an upper bound of A there exists $x \in E$ such that $\gamma<x$. According to thm.14.6 on $p .279$ there exists $\gamma^{\prime} \in \mathbb{R}$ such that $\gamma<\gamma^{\prime}<x$. It follows that $\gamma^{\prime} \notin B$, i.e., $\gamma^{\prime} \in A$, in contradiction to the fact that $\gamma \geqq$ a for all $a \in A$.

It follows that $\gamma \notin A$, i.e., $\gamma \in B$ and we conclude from cor. 14.1 that $\gamma=\min (B)$, i.e., $\gamma=\sup (E)$.

### 14.5 Constructing the Real Numbers via Cauchy Sequences

This chapter was created after discussions with Nguyen-Phan Tam about teaching the Math 330 course: she plans to construct the real numbers from the rationals by means of equivalence classes of Cauchy sequences in $\mathbb{Q}$.

In the following we always assume that $i, j, k, m, n \in \mathbb{N}, \varepsilon, p, q, r, s, p_{n}, p_{i, j}, \cdots \in \mathbb{Q}$, $x, y, z, x_{n}, x_{i, j}, \cdots \in \mathbb{R}$.
a. def. convergence in $\mathbb{Q}: \lim _{n \rightarrow \infty} q_{n}=q \Leftrightarrow \forall$ pos. $\varepsilon \in \mathbb{Q} \exists N \in \mathbb{Q}$ such that if $n \geqq N$ then $\left|q_{n}-q\right|<\varepsilon$.
b. def. Cauchy seqs. in $\mathbb{Q}:\left(q_{n}\right)_{n}$ is Cauchy $\Leftrightarrow \forall$ pos. $\varepsilon \in \mathbb{Q} \exists N \in \mathbb{Q}$ such that if $i, j \geqq$ $N$ then $\left|q_{i}-q_{j}\right|<\varepsilon$.
c. Let $\mathscr{C}:=\{$ all Cauchy sequences in $\mathbb{Q}\}$. For $\left(q_{n}\right)_{n},\left(r_{n}\right)_{n}$ we define $\left(q_{n}\right)_{n} \sim\left(r_{n}\right)_{n}$ iff $\lim _{n \rightarrow \infty}\left(r_{n}-q_{n}\right)=0$.
d. Let $q \in \mathbb{Q}$ and $q_{n}:=q \forall n$. Write $q$ for $\left[\left(q_{n}\right)_{n}\right]$.
e. Let $\mathbb{R}:=\mathscr{C} / \sim$. Show that for $\left[\left(p_{n}\right)\right],\left[\left(q_{n}\right)\right] \in \mathscr{C}$ the operations $\left(\left[\left(p_{n}\right)_{n}\right],\left[\left(q_{n}\right)_{n}\right]\right) \mapsto\left[\left(p_{n}+q_{n}\right)_{n}\right]$ and $\left(\left[\left(p_{n}\right)_{n}\right],\left[\left(q_{n}\right)_{n}\right]\right) \mapsto\left[\left(p_{n} \cdot q_{n}\right)_{n}\right]$ are well defined (do not depend on the particular members chosen from the equivalence classes).
f. Let $\left[\left(p_{n}\right)_{n}\right] \neq 0$ (i.e., $\lim _{n} p_{n} \neq 0$ ), i.e., we may assume $p_{n} \neq 0$ for all $n$. Show $-\left[\left(q_{n}\right)_{n}\right]:=$ $\left[\left(-q_{n}\right)_{n}\right]$ and $\left[\left(p_{n}\right)_{n}\right]^{-1}:=\left[\left(1 / p_{n}\right)_{n}\right]$ are additive and multiplicative inverses
g1. Define $\left[\left(p_{n}\right)_{n}\right]<\left[\left(q_{n}\right)_{n}\right]$ iff $\exists \varepsilon>0$ and $N \in \mathbb{N}$ such that $q_{n}-p_{n} \geqq \varepsilon \forall n \geqq N$.
g2. Define $\left[\left(p_{n}\right)_{n}\right] \leqq\left[\left(q_{n}\right)_{n}\right]$ iff $\forall \varepsilon>0$ exists $N \in \mathbb{N}$ such that $q_{n}-p_{n} \geqq-\varepsilon \forall n \geqq N$.
g3. show that $\left[\left(p_{n}\right)_{n}\right]<\left[\left(q_{n}\right)_{n}\right]$ iff $\left[\left(p_{n}\right)_{n}\right] \leqq\left[\left(q_{n}\right)_{n}\right]$ and $\left[\left(p_{n}\right)_{n}\right] \neq\left[\left(q_{n}\right)_{n}\right]$.
h. Show that $(\mathbb{R},+, \cdot,<)$ satisfies the axioms of $B / G$ ch. 8 with the exception of the completeness axiom.
Easy to see this specific item: If $\left[\left(p_{n}\right)_{n}\right]>0$ then there is $\left[\left(q_{n}\right)_{n}\right]>0$ such that $\left[\left(q_{n}\right)_{n}\right]<\left[\left(p_{n}\right)_{n}\right]$ : choose $\varepsilon>0$ as in $\mathbf{g 1}$ (remember: $\varepsilon \in \mathbb{Q}$ ) and set $q_{n}:=\varepsilon / 2$.
i. Embed $\mathbb{Q}$ into $\mathbb{R}: q \mapsto \bar{q}:=[(q, q, q, \ldots)]$.
j. Define limits and Cauchy sequences in $\mathbb{R} j u s t$ as in $\boldsymbol{a}$ and $\boldsymbol{b}$.
k. Let $\left(q_{n}\right)_{n}$ be Cauchy in $\mathbb{Q}$. Prove that $\overline{q_{n}} \rightarrow\left[\left(q_{j}\right)_{j}\right]$
l. Let $x_{n} \in \mathbb{R}$ such that $\left(x_{n}\right)_{n}$ is Cauchy in $\mathbb{R}$. With a density argument we find $q_{n} \in \mathbb{Q}$ such that $x_{n} \leqq \bar{q}_{n} \leqq x_{n}+1 / n$. Now show that 1. $\left(q_{n}\right)_{n}$ is Cauchy and then 2. $\lim _{n} x_{n}=\left[\left(q_{n}\right)_{n}\right]$.
m. Prove completeness according to $B / G$ : If nonempty $A \subseteq \mathbb{R}$ is bounded above then its set of upper bounds $U$ has a min: Let $Q_{n}:=\{i / j: i, j \in \mathbb{Z}$ and $j \leqq n\}$. Let $U_{n}:=U \cap Q_{n}$. Let $u_{n}:=\min \left(U_{n}\right)$ (exists because $n \cdot U_{n} \subset Z$ is bounded below and has a min. Easy to see that $u_{n}$ is Cauchy (in $\mathbb{Q}$ and, because distance $\left(u_{n}, A\right) \leqq 1 / n,\left[\left(u_{n}\right)_{n}\right]$ is the least upper bound of $A$.

Proofs for $\boldsymbol{k}$ and $\boldsymbol{l}$ in particular and an entire section on constructing $\mathbb{R}$ from $\mathbb{Q}$ by means of equivalence classes of Cauchy sequences can be found in [7] Haaser/Sullivan: Real Analysis.

## 15 Measure Theory ( $\star$ )

Note that this entire section is starred, hence optional.

## Introduction:

The following are the best known examples of measures ( $a_{j}, b_{j} \in \mathbb{R}$ ):

$$
\begin{aligned}
& \text { Length: } \quad \lambda^{1}\left(\left[a_{1}, b_{1}\right]\right):=b_{1}-a_{1}, \\
& \text { Area : } \quad \lambda^{2}\left(\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right]\right):=\left(b_{1}-a_{1}\right)\left(b_{2}-a_{2}\right), \\
& \text { Volume : } \quad \lambda^{3}\left(\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \times\left[a_{3}, b_{3}\right]\right):=\left(b_{1}-a_{1}\right)\left(b_{2}-a_{2}\right)\left(b_{3}-a_{3}\right) .
\end{aligned}
$$

Then there also are probability measures: $P\{$ a die shows a 1 or a 6$\}=1 / 3$.
We will explore in this chapter some of the basic properties of measures.

### 15.1 Basic Definitions

Definition 15.1 (Extended real valued functions).

$$
\overline{\mathbb{R}}_{+}:=\mathbb{R}_{+} \cup\{+\infty\}=\{x \in \mathbb{R}: x \geqq 0\} \cup\{+\infty\}
$$

is the set of all non-negative real numbers augmented by the element $\infty$.
A mapping $F$ whose codomain is a subset of

$$
\overline{\mathbb{R}}:=\mathbb{R} \cup\{\infty\} \cup\{-\infty\}
$$

is called an extended real valued function.
There are many issues with functions that allow some arguments to have infinite value (hint: if $F(x)=\infty$ and $F(y)=\infty$, what is $F(x)-F(y)$ ?)

We only list the following rule which might come unexpected to you:

$$
\begin{equation*}
0 \cdot \pm \infty= \pm \infty \cdot 0=0 \tag{15.1}
\end{equation*}
$$

Definition 15.2 (Abstract measures). Let $\Omega$ be a nonempty set and let $\mathfrak{F}$ be a set that contains some, but not necessarily all, subsets of $\Omega . \mathfrak{F}$ is called a $\sigma$-algebra or $\sigma$-field for $\Omega$ if it satisfies the following:

$$
\begin{array}{rll}
\emptyset \in \mathfrak{F} & \text { and } & \Omega \in \mathfrak{F} \\
A \in \mathfrak{F} & \Longrightarrow & C A \in \mathfrak{F} \\
\left(A_{n}\right)_{n \in \mathbb{N}} \in \mathfrak{F} & \Longrightarrow & \bigcup_{n \in \mathbb{N}} A_{n} \in \mathfrak{F} \text { and } \bigcap_{n \in \mathbb{N}} A_{n} \in \mathfrak{F} \tag{15.2c}
\end{array}
$$

The pair $(\Omega, \mathfrak{F})$ is called a measurable space. Note that $\mathfrak{F}$ is a set whose elements themselves are sets! The elements of $\mathfrak{F}$ are called measurable sets.

A measure on $\mathfrak{F}$ is an extended real valued function

$$
\mu(\cdot): \mathfrak{F} \rightarrow \overline{\mathbb{R}}_{+} \quad A \mapsto \mu(A)
$$

with the following properties:

$$
\begin{align*}
& \mu(\emptyset)=0  \tag{15.3a}\\
& A, B \in \mathfrak{F} \text { and } A \subseteq B \quad \Longrightarrow \quad \mu(A) \leqq \mu(B)  \tag{15.3b}\\
& \left(A_{n}\right)_{n \in \mathbb{N}} \in \mathfrak{F} \text { mutually disjoint } \quad \Longrightarrow \quad \mu\left(\biguplus_{n \in \mathbb{N}} A_{n}\right)=\sum_{n \in \mathbb{N}} \mu\left(A_{n}\right) \quad(\sigma \text {-additivity) }
\end{align*}
$$

where mutual disjointness means that $A_{i} \cap A_{j}=0$ for any $i, j \in \mathbb{N}$ such that $i \neq j$ (see def.2.4 on p.14). The triplet $(\Omega, \mathfrak{F}, \mu)$ is called a measure space

We call $\mu$ a finite measure on $\mathfrak{F}$ if $\mu(\Omega)<\infty$.
A measure space can support many different measures.
If $\mu(\Omega)=1$ then $\mu(\cdot)$ is called a probability measure. Traditionally, mathematicians write $P(A)$ rather than $\mu(A)$ for probability measures and the elements of $\mathfrak{F}$ (the measurable subsets) are thought of as events for which $P(A)$ is interpreted as the probability with which the event $A$ might happen.

Example 15.1 (Lebesgue measure). The most important measures we encounter in real life are those that measure the length of sets in one dimension, the area of sets in two dimensions and the volume of sets in three dimensions. Given intervals $[a, b] \in \mathbb{R}$, rectangles $\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \in \mathbb{R}^{2}$, boxes or quads $\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \times\left[a_{3}, b_{3}\right] \in \mathbb{R}^{3}$ and $n$-dimensional parallelepipeds $\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \times \cdots \times$ $\left[a_{n}, b_{n}\right] \in \mathbb{R}^{n}$, we define

$$
\begin{align*}
& \lambda^{1}([a, b]):=b-a, \\
& \lambda^{2}\left(\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right]\right):=\left(b_{1}-a_{1}\right)\left(b_{2}-a_{2}\right), \\
& \lambda^{3}\left(\left[a_{1}, b_{1}\right] \times\left[a_{2}, b_{2}\right] \times\left[a_{3}, b_{3}\right]\right):=\left(b_{1}-a_{1}\right)\left(b_{2}-a_{2}\right)\left(b_{3}-a_{3}\right),  \tag{15.4}\\
& \lambda^{n}\left(\left[a_{1}, b_{1}\right] \times \cdots \times\left[a_{n}, b_{n}\right]\right):=\left(b_{1}-a_{1}\right)\left(b_{2}-a_{2}\right) \ldots\left(b_{n}-a_{n}\right)
\end{align*}
$$

It can be shown that any measure that is defined on all parallelepipeds in $\mathbb{R}^{n}$ can be uniquely extended to a measure on the $\sigma$-algebra $\mathcal{B}^{n}$ generated by those parallelepipeds ${ }^{137} \lambda^{n}$ is called $n$ dimensional Lebesgue measure

Note that Lebesgue measure is not finite.
Example 15.2. You can easily verify that the following set function defines a measure on an arbitrary nonempty set $\Omega$ with an arbitrary $\sigma$-field $\mathfrak{F}$.

$$
\mu(\emptyset):=0 ; \quad \mu(A):=\infty \text { if } A \neq \emptyset
$$

Keep this example in mind if you contemplate infinity of measures.
Remark 15.1 (Finite disjoint unions). The $\sigma$-additivity of measures is what makes working with them such a pleasure in many ways. You can now express it as follows: Given any mutually disjoint sequence of measurable sets, the measure of the disjoint union is the sum of the measures. The last property (15.2c) for $\sigma$-algebras is required for exactly that reason: you cannot take advantage of the $\sigma$-additivity of a measure $\mu$ if its domain does not contain countable unions and intersections of all its constituents.

[^100]Note that if we have only finitely many sets then " $\sigma$-additivity" which stands for "additivity of countably many" becomes simple additivity. We obtain the following by setting $A_{N+1}=A_{N+2}=\ldots=0$ :

$$
\begin{align*}
& A_{1}, A_{2}, \ldots, A_{N} \in \mathfrak{F} \text { mutually disjoint }  \tag{15.5}\\
& \quad \Rightarrow \mu\left(A_{1} \uplus A_{2} \uplus \ldots \uplus A_{N}\right)=\mu\left(A_{1}\right)+\mu\left(A_{2}\right)+\ldots+\mu\left(A_{N}\right) \quad \text { (additivity). }
\end{align*}
$$

In the case of only two disjoint measurable sets $A$ and $B$ the above simply becomes

$$
\mu(A \uplus B)=\mu(A)+\mu(B) .
$$

In many circumstances you have a set function on a $\sigma$-algebra which behaves like a measure but you can only prove that it is additive instead of $\sigma$-additive. You should not be surprised that there is a special name for those "generalized measures":

Definition 15.3 (Contents as additive measures). Let $\Omega$ be a nonempty set and let $\mathfrak{F}$ be a $\sigma$-algebra for $\Omega$.

A content on $\mathfrak{F}$ is a real valued function $\quad m(\cdot): \mathfrak{F} \rightarrow \mathbb{R}, \quad A \mapsto m(A)$ which satisfies

$$
\begin{array}{ll}
m(\emptyset)=0 \quad \text { (positivity) }  \tag{15.6a}\\
A, B \in \mathfrak{F} \text { and } A \subseteq B \Rightarrow m(A) \leqq m(B) & \text { (monotony) } \\
A_{1}, A_{2}, \ldots, A_{N} \in \mathfrak{F} \text { mutually disjoint } \Rightarrow & m\left(\biguplus_{n=1}^{N} A_{n}\right)=\sum_{n=1}^{N} m\left(A_{n}\right) \quad \text { (additivity). }
\end{array}
$$

Note that $\mu(\Omega)<\infty$ for a content $\mu$. After this digression on contents let us go back to measures.
Proposition 15.1 (Simple properties of measures). $A$ : Let $A, B, \in \mathfrak{F}$ and let $\mu$ be a measure on $\mathfrak{F}$. Then

$$
\begin{align*}
& \mu(A) \geqq 0 \quad \text { for all } A \in \mathfrak{F}  \tag{15.7a}\\
& A \subseteq B \Rightarrow \mu(B)=\mu(A)+\mu(B \backslash A)  \tag{15.7b}\\
& \mu(A \cup B)+\mu(A \cap B)=\mu(A)+\mu(B) \tag{15.7c}
\end{align*}
$$

If $\mu$ is finite then we can write

$$
\begin{align*}
& A \subseteq B \Rightarrow \mu(B \backslash A)=\mu(B)-\mu(A)  \tag{15.8a}\\
& \mu(A \cup B)=\mu(A)+\mu(B)-\mu(A \cap B) \tag{15.8b}
\end{align*}
$$

PROOF of $\boldsymbol{A}$ : The first property follows from the fact that $\mu(\emptyset)=0, \emptyset \subseteq A$ for all $A \in \mathfrak{F}$ and (15.3b). To prove the second property, observe that $B=A \uplus(B \backslash A)$.

Proving the third property is more complicated because neither $A$ nor $B$ may be a subset of the other. We first note that because $A \backslash B \subseteq A, B \backslash A \subseteq A$ and $A \cap B \subseteq A, \mu(A \cup B)=\infty$ can only be true if $\mu(A)=\infty$ or $\mu(B)=\infty$. In this case (15.7c) is obviously true. Hence we may assume that $\mu(A \cup B)<\infty$. We have

$$
\begin{align*}
& A \cup B=(A \cap B) \uplus(B \backslash A) \uplus(A \backslash B)  \tag{15.9a}\\
& A \cup B=A \uplus(B \backslash A)=B \uplus(A \backslash B) \tag{15.9b}
\end{align*}
$$

It follows from (15.9a) that

$$
\begin{equation*}
\mu(A \cup B)=\mu(A \cap B)+\mu(B \backslash A)+\mu(A \backslash B) \tag{15.10}
\end{equation*}
$$

It follows from (15.9b) that

$$
\begin{equation*}
2 \cdot \mu(A \cup B)=\mu(A)+\mu(B \backslash A)+\mu(B)+\mu(A \backslash B) \tag{15.11}
\end{equation*}
$$

We subtract the left and right sides of (15.10) from those of (15.11) and obtain

$$
\begin{aligned}
\mu(A \cup B) & =\mu(A)+\mu(B \backslash A)+\mu(B)+\mu(A \backslash B)-\mu(A \cap B)-\mu(B \backslash A)-\mu(A \backslash B) \\
& =\mu(A)+\mu(B)-\mu(A \cap B)
\end{aligned}
$$

and the third property is proved.
Proposition 15.2 (Minimal sigma-algebras). Let $\Omega$ be a nonempty set.
A: The intersection of arbitrarily many $\sigma$-algebras is a $\sigma$-algebra.
$B$ : Let $\mathfrak{K}$ be a set which contains subsets of $\Omega$. It is not assumed that $\mathfrak{K}$ is a $\sigma$-algebra. Then there exists a $\sigma$-algebra which contains $\mathfrak{K}$ and is minimal in the sense that it is contained in any other $\sigma$-algebra that also contains $\mathfrak{K}$. We name this $\sigma$-algebra $\sigma(\mathfrak{K})$ because it clearly depends on $\mathfrak{K}$. It is constructed as follows:

$$
\begin{equation*}
\sigma(\mathfrak{K})=\bigcap\{\mathfrak{A}: \mathfrak{A} \supseteq \mathfrak{K} \text { and } \mathfrak{A} \text { is a } \sigma \text {-algebra for } \Omega\} \text {. } \tag{15.12}
\end{equation*}
$$

PROOF of $A$ :
We must prove (15.2a), (15.2b) and (15.2c). Let $\left(\mathfrak{A}_{\alpha}\right)_{\alpha}$ be an arbitrary family of $\sigma$-algebras for $\Omega$. Let

$$
\mathfrak{A}:=\bigcap_{\alpha} \mathfrak{A}_{\alpha} .
$$

$\emptyset$ and $\Omega$ belong to each $\sigma$-algebra according to (15.2a). It follows that they both belong to the intersection $\bigcap_{\alpha} \mathfrak{A}_{\alpha}$, i.e., $\mathfrak{A}$ satisfies (15.2a). Let $A \in \mathfrak{A}$. Then $A \in \mathfrak{A}_{\alpha}$ for each $\alpha$. CA belongs to each $\sigma$-algebra according to (15.2b). It follows that $\complement A \in \bigcap_{\alpha} \mathfrak{A}_{\alpha}$, i.e., $\mathfrak{A}$ satisfies (15.2b). Finally, let $A_{n} \in \mathfrak{A}$ for all $n \in \mathbb{N}$.
Then $A_{n} \in \mathfrak{A}_{\alpha}$ for all $n \in \mathbb{N}$ and for each $\alpha, \cup_{n \in \mathbb{N}} A_{n}$ and $\bigcap_{n \in \mathbb{N}} A_{n}$ both belong to each $\sigma$-algebra according to (15.2c). It follows that they both belong to the intersection $\bigcap_{\alpha} \mathfrak{A}_{\alpha}$, i.e., $\mathfrak{A}$ satisfies (15.2c). It follows that $\mathfrak{A}$ is a $\sigma$-algebra.
PROOF of $B$ :
First of all, we know that $\sigma(\mathfrak{K})$ is an intersection of $\sigma$-algebras and, according to part $\boldsymbol{A}$ of this proposition, really is a $\sigma$-algebra. We now prove that $\sigma(\mathfrak{K})$ contains $\mathfrak{K}$ and is the minimal $\sigma$-algebra with that property. First let us prove that $\sigma(\mathfrak{K}) \supseteq \mathfrak{K}$. But that is obvious because it is the intersection of sets all of which contain $\mathfrak{K}$. On the other hand, $\sigma(\mathfrak{K})$ is the intersection of all $\sigma$-algebras that contain $\mathfrak{K}$, so it is impossible for any other $\sigma$-algebra to both be a strict subset of $\sigma(\mathfrak{K})$ and also contain $\mathfrak{K}$.

### 15.2 Sequences of Sets - limsup and liminf

Assumption 15.1 (Existence of a universal set). We assume the existence of a set $X$ which contains all sets $A_{n}, B_{n}, C_{n}$ that are used here in sequences.
Definition 15.4 (Monotone set sequences). A sequence $A_{k}$ of arbitrary subsets of $X$ is called

$$
\begin{align*}
\text { non-decreasing } & \text { if } A_{1} \subseteq A_{2} \subseteq \ldots  \tag{15.13a}\\
\text { non-increasing } & \text { if } A_{1} \supseteq A_{2} \supseteq \ldots  \tag{15.13b}\\
\text { strictly increasing } & \text { if } A_{1} \subsetneq A_{2} \subsetneq \ldots  \tag{15.13c}\\
\text { strictly decreasing } & \text { if } A_{1} \supsetneq A_{2} \supsetneq \ldots \tag{15.13d}
\end{align*}
$$

Each one of those sequences is called a monotone set sequence.

## Might as well define limits of monotone sequences of sets. It's certainly intuitive enough:

Definition 15.5 (Limits of monotone set sequences). Given are sets $A_{n}, B_{n} \subseteq X(n \in \mathbb{N})$. Assume that

$$
\begin{array}{ll}
A_{1} \subseteq A_{2} \subseteq A_{3} \subseteq \ldots & \text { and let } A:=\bigcup_{k \in \mathbb{N}} A_{k} \\
B_{1} \supseteq B_{2} \supseteq B_{3} \supseteq \ldots & \text { and let } B:=\bigcap_{k \in \mathbb{N}} B_{k}
\end{array}
$$

We say that $A$ is the limit of the sequence $\left(A_{j}\right)_{j \in \mathbb{N}}$ and $B$ is the limit of the sequence $\left(B_{j}\right)_{j \in \mathbb{N}}$ and we write

$$
\begin{array}{lll}
A=\lim _{n \rightarrow \infty} A_{n} & \text { or } \quad A_{n} \nearrow A \text { for } n \rightarrow \infty, \\
B=\lim _{n \rightarrow \infty} B_{n} & \text { or } \quad B_{n} \searrow B \text { for } n \rightarrow \infty . \tag{15.14b}
\end{array}
$$

The above are not terribly useful definitions. What does it matter whether we write $A=\lim _{n \rightarrow \infty} A_{n}$ or $A=\cup_{k \in \mathbb{N}} A_{k}$ ? Things would be very different if we went further and defined limits of sequences of sets. Doing so is at the very beginning of a branch of Mathematics called Measure Theory and its (slightly) more applied version, Abstract Probability Theory.

Definition 15.6 (lim inf and lim sup of set sequences). Given are sets $A_{n}, B_{n} \subseteq X(n \in \mathbb{N})$. Let

$$
\begin{array}{ll}
\liminf _{n \rightarrow \infty} A_{n}:=\bigcup_{n \in \mathbb{N}} \bigcap_{k \geqq n} A_{k} & \text { (limit inferior) } \\
\limsup _{n \rightarrow \infty} A_{n}:=\bigcap_{n \in \mathbb{N}} \bigcup_{k \geqq n} A_{k} & \text { (limit superior) } \tag{15.16}
\end{array}
$$

In general those two will not coincide. But if they do then we define

$$
\begin{equation*}
\lim _{n \rightarrow \infty} A_{n}:=\liminf _{n \rightarrow \infty} A_{n}=\limsup _{n \rightarrow \infty} A_{n} \tag{15.17}
\end{equation*}
$$

We call $\lim _{n \rightarrow \infty} A_{n}$ the limit of the sequence $\left(A_{n}\right)$ and we write

$$
A_{n} \rightarrow A \quad \text { for } n \rightarrow \infty
$$

The following comments should make matters easier to understand if you abbreviate
Lemma 15.1 (lim inf and lim sup as monotone limits). Given are sets $A_{n}, B_{n} \subseteq X(n \in \mathbb{N})$. Let

$$
\begin{align*}
A_{\star n} & :=\bigcap_{k \geqq n} A_{k} & \text { Then } & A_{\star n} \nearrow \liminf _{n \rightarrow \infty} A_{n}  \tag{15.18}\\
A_{n}^{\star} & :=\bigcup_{k \geqq n} A_{k} & \text { Then } & A^{\star}{ }_{n} \searrow \limsup _{n \rightarrow \infty} A_{n} \tag{15.19}
\end{align*}
$$

PROOF: Let $m, n \in \mathbb{N}$ such that $m<n$. Then

$$
\begin{aligned}
& A_{\star m}=\bigcap_{k=m}^{n-1} A_{k} \cap \bigcap_{k \geqq n} A_{k}=\bigcap_{k=m}^{n-1} A_{k} \cap A_{\star n} \subseteq A_{\star n} \\
& A^{\star} m=\bigcup_{k=m}^{n-1} A_{k} \cup \bigcup_{k \geqq n} A_{k}=\bigcup_{k=m}^{n-1} A_{k} \cup A^{\star}{ }_{n} \supseteq A^{\star}{ }_{n}
\end{aligned}
$$

This proves that $A_{\star n}$ is non-decreasing and $A * \star_{n}$ is non-increasing. By the very definition of the limit of a monotone sequence of sets it is true that

$$
\begin{aligned}
& \lim _{n \rightarrow \infty} A_{\star n}=\bigcup_{n \in \mathbb{N}} A_{\star n}=\liminf _{n \rightarrow \infty} A_{n} \\
& \lim _{n \rightarrow \infty} A^{\star}{ }_{n}=\bigcap_{n \in \mathbb{N}} A^{\star}{ }_{n}=\limsup _{n \rightarrow \infty} A_{n}
\end{aligned}
$$

### 15.3 Conditional Expectations as Generalized Averages

EMPTY!!

## 16 Appendix: Addenda to Beck/Geoghegan's "The Art of Proof"

This chapter contains extensions of material found in [1] Beck/Geoghegan, the book which is meant to be read in conjunction with these lecture notes. Some of this material is referenced in earlier chapters.

### 16.1 AoP Ch.1: Integers

### 16.1.1 Limited Division in $\mathbb{Z}$

Proposition 16.1. Let $m, n \in \mathbb{Z}$ such that $m \mid n$ and $m \neq 0$, i.e., there exists $j \in \mathbb{Z}$ be such that $n=j \cdot m$. Then $j$ is uniquely determined.

PROOF: Assume that there is $j^{\prime} \in \mathbb{Z}$ such that also $n=j^{\prime} \cdot m$. Then

$$
n=j^{\prime} \cdot m=j \cdot m, \quad \text { hence }\left(j^{\prime}-j\right) m=0 .
$$

$B / G$ Proposition 1.26 states that if $a, b \in \mathbb{Z}$ and if $a \cdot b=0$, then $a=0$ or $b=0$.
It follows from $m \neq 0$ that $j^{\prime}-j=0$ and hence $j^{\prime}=j$.
The above result implies that the integer $j$ is unique determined by $m$ and $n$ and hence allows $u$ s to make the following definition.

Definition 16.1 (Quotients). Let $m, n \in \mathbb{Z}$ such that $m \mid n$ and $m, n \neq 0$. Let $j \in \mathbb{Z}$ be the unique integer for which $n=j \cdot m$. We write either of

$$
\begin{equation*}
\frac{n}{m}, \quad n / m, \quad n \div m \tag{16.1}
\end{equation*}
$$

instead of $j$ and we call $n$ the dividend or numerator, $m$ the divisor or denominator, and $j$ the quotient of the expression $n / m$.

Note 16.1. Note that the assignment $(m, n) \mapsto n / m$ is not a "binary operation" on $\mathbb{Z}$ as is the case for $(m, n) \mapsto m+n$ and $(m, n) \mapsto m \cdot n$. The reason: $m+n$ and $m \cdot n$ are defined for any two $m, n \in \mathbb{Z}$ whereas $n / m$ is only defined for those $m, n \in \mathbb{Z}$ which satisfy the condition $m \mid n$.

### 16.2 AoP Ch.2: Natural Numbers and Induction

### 16.2.1 AoP Ch.2.2 (Ordering the Integers)

The material given here complements ch.2.2 (Ordering the Integers) of [1] Beck/Geoghegan.
Remark 16.1. B/G Axioms 1.1 - 1.5 together with axiom 2.1 do not suffice to characterize the integers; a symbol different from $\mathbb{Z}$ might have been more appropriate. These axioms are also satisfied by the set $\mathbb{Q}$ of all rational numbers (fractions), provided one interprets the set $\mathbb{N}$ described by axiom 2.1 as the set of all (strictly) positive fractions. They are also satisfied by the set $\mathbb{R}$ of all real numbers (decimals), provided one interprets $\mathbb{N}$ as the set of all (strictly) positive decimals. Only addition of the induction axiom ( $\mathrm{B} / \mathrm{G}$ axiom 2.15 ) excludes $\mathbb{Q}$ and $\mathbb{R}$.

The next two propositions belong after $B / G$ prop. $2.3(1 \in \mathbb{N})$.
Proposition 16.2. If $k \in \mathbb{Z}$ then $k+1>k$.
Proof: Left as an exercise.
Corollary 16.1. $1>0$.
PROOF: This follows from $1=1+0$ and prop.16.2.1, applied to $k:=0$.
The next proposition belongs after B/G prop.2.4 ("§" is transitive).
Proposition 16.3. Let $m, n, p \in \mathbb{Z}$ such that

$$
\begin{equation*}
m \leqq n, \quad n \leqq p \tag{16.2}
\end{equation*}
$$

Then $m \leqq p$.
PROOF: There are four cases.

1. $m<n, n<p$ : It follows from B/G prop. 2.4 (transitivity of " $<$ ") that $m<p$, in particular, $m \leqq p$.
2. $m<n, n=p$ : It follows that $m<p$. This implies $m \leqq p$.
3. $m=n, n<p$ : It follows again that $m<p$. This implies $m \leqq p$.
4. $m=n, n=p$ : It follows that $m=p$. This implies $m \leqq p$.

### 16.2.2 AoP Ch.2.3 (Induction)

The material given here complements ch.2.3 (Induction) of [1] Beck/Geoghegan.
Proposition 16.4 (B/G Prop.2.20). If $k \in \mathbb{N}$ then

$$
\begin{equation*}
k \geqq 1 \tag{16.3}
\end{equation*}
$$

PROOF: The proof is done by induction on $k$.
Base case: For $k=1$ (16.3) reads $1 \geqq 1$. This is true because " $=$ " implies " $\geqq "$.
Induction step: Let $k \in \mathbb{N}$ and assume that $k \geqq 1$. We are done if we can prove under this assumption that

$$
\begin{equation*}
k+1 \geqq 1 \tag{16.4}
\end{equation*}
$$

According to prop.16.2, $k+1>k$, hence $k+1 \geqq k$. Further $k \geqq 1$ (our induction assumption.) It follows from the transitivity of " $\leqq$ " (prop.16.3) that $k+1 \geqq 1$, i.e., (16.4) is valid.
Proposition 16.5 (B/G Prop.2.21). There exists no $x \in \mathbb{Z}$ such that $0<x<1$.
PROOF: Assume to the contrary that there is some $x \in \mathbb{Z}$ such that $0<x$, i.e., $x>0$, and $x<1$.
It follows from $x>0$ that $x \in\{k \in \mathbb{Z}: k>0\}$. But according to $B / G$ prop.2.13 this set equals $\mathbb{N}$. So $x \in \mathbb{N}$. It now follows from $B / G$ prop. 2.20 that $x \geqq 1$, contrary to the assumption that $x<1$.

Corollary 16.2 (B/G Cor.2.22). Let $n \in \mathbb{Z}$. There exists no $x \in \mathbb{Z}$ such that $n<x<n+1$.
PROOF: Assume to the contrary that such $x$ exists. According to B/G prop.2.7(i) an inequality between integers is preserved if one adds the same item to both of its sides. It follows that

$$
n-n<x-n<(n+1)-n, \quad \text { i.e., } 0<x-n<1 \text {. }
$$

The integer $x^{\prime}:=x-n$ satisfies $0<x^{\prime}<1$, a contradiction to prop.16.5.
Proposition 16.6 (sharpening of $\mathrm{B} / \mathrm{G}$ Prop.2.13). $\mathbb{N}=\{k \in \mathbb{Z}: k \geqq 1\}$.
PROOF: This follows from $1>0$ (cor. 16.1 on p.289) and prop. 16.5 above.

### 16.2.3 Bounded Sets in $\mathbb{Z}$

The material given here complements ch.2.4 (The Well-Ordering Principle) of [1] Beck/Geoghegan.
Definition 16.2 (Upper and lower bounds, maxima and minima in $\mathbb{Z}$ ). Let $A \subseteq \mathbb{Z}$. ${ }^{138}$ Let $l, u \in \mathbb{Z}$. We call $l$ a lower bound of $A$ if $l \leqq a$ for all $a \in A$.

We call $u$ an upper bound of $A$ if $u \geqq a$ for all $a \in A$.
We call $A$ bounded above if this set has an upper bound and we call $A$ bounded below if $A$ has a lower bound. We call $A$ bounded if $A$ is both bounded above and bounded below.

A minimum (min) of $A$ is a lower bound $l$ of $A$ such that $l \in A$. A maximum (max) of $A$ is an upper bound $u$ of $A$ such that $u \in A$.

The next proposition states that $\min$ and max are unique if they exist. This makes it possible to write $\min (A)$ or $\min A$ for the minimum of $A$ and $\max (A)$ or $\max A$ for the maximum of $A$.

Proposition 16.7. Let $A \subseteq \mathbb{Z}$. If $A$ has a maximum then it is unique. If $A$ has a minimum then it is unique.
PROOF: Left as an exercise.

Definition 16.3 (Translation and dilation of sets of integers). Let $A \subseteq \mathbb{Z}$. and $\alpha, b \in \mathbb{Z}$. We define

$$
\begin{equation*}
\lambda A+b:=\{\lambda a+b: a \in A\} . \tag{16.5}
\end{equation*}
$$

In particular, for $\lambda= \pm 1$, we obtain

$$
\begin{align*}
A+b & =\{a+b: a \in A\},  \tag{16.6}\\
-A & =\{-a: a \in A\} . \tag{16.7}
\end{align*}
$$

Remark 16.2. Note that the above makes sense for any algebraic structure ${ }^{139}$ with binary operations " + " and/or ".", e.g., groups or vector spaces (in the latter case " $(\lambda, a) \mapsto \lambda a$ " would be the scalar product between a real number (scalar) $\lambda$ and a vector $a$.

[^101]Theorem 16.1 (Generalization of the Well-Ordering Principle).
a. Let $A$ be a nonempty subset of $\mathbb{Z}$ which is bounded below. Then $A$ possesses a minimum in $\mathbb{Z}$.
b. Let $B$ be a nonempty subset of $\mathbb{Z}$ which is bounded above. Then $B$ possesses a maximum in $\mathbb{Z}$.
c. Let $C$ be a nonempty bounded subset of $\mathbb{Z}$. Then $C$ possesses both minimum and maximum in $\mathbb{Z}$.

Proof (outline):
a. If $A$ has 1 as a lower bound then $A \subseteq \mathbb{N}$ and the theorem simply is the Well-Ordering Principle ( $B / G$ theorem 2.32). Next we just assume that $A$ is bounded below. Let $a_{\star}$ be a lower bound of $A$. Let $A^{\prime}:=A-a_{\star}+1$. Then $a^{\prime} \geq 1$ for all $a^{\prime} \in A^{\prime}$, i.e., $A^{\prime} \subseteq \mathbb{N}$. and it follows from the Well-Ordering Principle that the minimum $\min \left(A^{\prime}\right)$ of $A^{\prime}$ exists.

It is easy to see from $\min \left(A^{\prime}\right) \in A^{\prime}$ that then $m:=\min \left(A^{\prime}\right)+a_{\star}-1 \in A$ and that $m$ is a lower bound of $A$ because $a_{\star}$ is a lower bound of $A$. It follows that $m=\min (A)$.
b. We assume that $B$ is bounded above. Let $b^{\star}$ be an upper bound of $B$. Let $B^{\prime}:=-B$. Then $B^{\prime}$ has $-b^{\star}$ as a lower bound and it follows from the already proven part a that the minimum $\min \left(B^{\prime}\right)$ of $B^{\prime}$ exists. Let $m:=-\min \left(B^{\prime}\right)$. It follows from $\min \left(B^{\prime}\right) \in B^{\prime}$ that $m \in-B^{\prime}=-(-B)=B$ and it follows from $\min \left(B^{\prime}\right) \leqq b^{\prime}$ for all $b^{\prime} \in B^{\prime}$ that $m \geqq b$ for all $b \in B$. But then $m$ must be the maximum of $A$.
$\boldsymbol{c}$. is a trivial consequence of $\boldsymbol{a}$ and $\boldsymbol{b}$
Example 16.1 (The Well-Ordering Principle is not true in $\mathbb{Q}$ and $\mathbb{R}$ ).
a. $\mathbb{R}$ : The set $A:=\left\{x \in \mathbb{R}: x^{2}<2\right\}$ is bounded in $\mathbb{R}($ by $\pm 2)$ but has neither min (would have to be $-\sqrt{2} \notin A$ ) nor $\max$ (would have to be $+\sqrt{2} \notin A$ ).
but: $-\sqrt{2}$ is a special lower bound: it is the greatest lower bound $\inf (A)$ of $A$ and $\sqrt{2}$ is a special upper bound: it is the least (smallest) upper bound $\sup (A)$ of $A$.
b. $\mathbb{Q}$ : The set $B:=\left\{x \in \mathbb{Q}: x^{2}<2\right\}=A \cap \mathbb{Q}$ is bounded in $\mathbb{Q}$ (by $\pm 2$ ) and also has neither min nor max for the same reasons as $A$.
Further: $-\sqrt{2}$ is not a lower bound of $B$ and $\sqrt{2}$ is not an upper bound of $B$ because those numbers are not in our "universe" $\mathbb{Q}$. The set $B$ has neither min, max, inf, sup!

Theorem 16.2 ( $\mathrm{B} / \mathrm{G}$ thm.10.1: $\mathbb{N}$ is unbounded in $\mathbb{R}$ ). For any $x \in \mathbb{R}$ there exists $n \in \mathbb{N}$ such that $n>r$, i.e., there are no upper bounds for $\mathbb{N}$ in $\mathbb{R}$.

PROOF: Assume to the contrary that there exists an upper bound of $\mathbb{N}$. According to thm16.1 (Generalization of the Well-Ordering Principle) on p. $291 \mathbb{N}$ has a maximum. Let $u^{\star}:=\max (\mathbb{N})$. Then $u^{\star}+1$ belongs to $\mathbb{N}$ as the sum of two natural numbers. It follows from $u^{\star}+1>u^{\star}$ that $u^{\star}$ is not the largest element of $\mathbb{N}$ and we have reached a contradiction.

### 16.2.4 Exercises for Ch.16.2

Exercise 16.1. Use everything up to AND including B/G prop.2.3, to prove prop. 16.2 on p. 289 of this document: If $k \in \mathbb{Z}$ then $k+1>k$.

### 16.3 AoP Ch.3: Some Points of Logic

Here are some references for items discussed in ch. 3 of the $B / G$ text to where they appear in ch. 3 (Logic) on p. 23 of this document.
a. Universal quantifier $\forall$, existental quantifier $\exists$, unique existental quantifier $\exists$ !: ch.3.6.1 on p.50.
b. " $(\forall x)(\forall y)$ has the same meaning as $(\forall x$ and $y)$ ": part a of def.3.19 (Doubly quantified expressions) in ch.3.6.2 on p.52.
We recommend that you look at prop.3.1 and the note which precedes it. The latter is reproduced here:

1. The order in which the qualifiers are applied is important.
$\forall x \exists y$ generally does not mean the same as $\exists y \forall x$.
2. Interchanging variable names in the qualifiers is not $O K$.
$\forall x \exists y$ generally does not mean the same as $\forall y \exists x$.
c. Skip part c if you have no knowledge of logic beyond what is in ch. 3 of B/G.

If you have had some training in logic you may have learned to express "if $P$ (is true) then $Q$ (is true)" as " $P \rightarrow Q$ " rather than " $P \Rightarrow Q$ ". There is a difference: Proving a statement of the form "if $P$ then $Q$ " means to show that the statements $P$ and $Q$ are related in a fashion that makes it "logically impossible", ${ }^{140}$ for $P$ to be true and $Q$ to be false. " $P \Rightarrow Q$ " has the combination ( $P$ is true, $Q$ is false ) marked as irrelevant (logically impossible); the remaining three combinations give the same outcome as for $P \rightarrow Q$, i.e., they all evaluate to true.
Example: If $P:$ " $x=100$ " and $Q:$ " $x>10$ " is correct because it is not possible that $P$ is true and $Q$ is false.
d. The equivalent forms of "if $P$ then $Q$ " such as " $Q$ whenever $P$ " are listed in ch.3.3.5 (Arrow and Implication Operators) on p.37. There you also find the definition of the converse and the contrapositive of an implication.

B/G ch.3.3: Negations:
a. The negation of " $A$ and $B$ " is "not $A$ or not $B$ " and the negation of " $A$ or $B$ " is "not $A$ and not B". This is known as "De Morgan's law" for statements (see thm.3.3 on p.42). Do you see the connection to De Morgan's law for sets? (see thm.5.1 on p.97).
b. Skip part b if you have no knowledge of logic beyond what is in ch. 3 of B/G.
$B / G$ states that the negation of " $P \Rightarrow Q^{\prime}$ " is " $P$ and not $Q$ ". This should be stated more appropriately as follows: The negation of " $P \Rightarrow Q$ " is "Pand not $Q$ ". The reason for this equivalence is that " $P \Rightarrow Q$ " is defined as "not $P$ or $Q$ ": we can apply De Morgan's laws to obtain the negation.
e. In the above observe the "binding power" or preference of the logical operators: The negation "not $P$ " of a statement $P$ has higher preference than " $P$ and $Q$ " and " $P$ or $Q$ ":
The meaning of "not $P$ or $Q$ " is " $(\operatorname{not} P)$ or $Q$ ", not "not $(P$ or $Q)$."
e. $x x x$

### 16.4 AoP Ch.4: Recursion

Some of the recursive definitions in chapter 4 of AoP, such as $\sum x_{j}$ and $\prod x_{j}$ were defined to start at $j=1$. One cannot deduce from those definitions what it means to start at some index other than 1 . We fix this shortcoming in this chapter.

[^102]
### 16.4.1 Generalization of Some Recursive Definitions

Definition 16.4 (Extension of the definition $\sum_{j=1}^{n} x_{j}$ ). We modify the definition of " $\Sigma$ " (B/G p.34, 35) as follows: Let $k \in \mathbb{Z}$ and let $\left(x_{j}\right)_{j=k}^{\infty}$ be a sequence of integers.

For each $n \in \mathbb{Z}$ such that $k \leq n$, we define an integer called $\sum_{j=k}^{n} x_{j}$ as follows. ${ }^{141}$

$$
\begin{equation*}
\text { (i) } \quad \sum_{j=k}^{k} x_{j}=x_{k}, \quad \text { (ii) } \quad \sum_{j=k}^{n+1} x_{j}=\sum_{j=k}^{n} x_{j}+x_{n+1} \text {. } \tag{16.8}
\end{equation*}
$$

We call $\sum_{j=k}^{n} x_{j}$ the sum of the integers $x_{k}, x_{k+1}, \ldots, x_{n-1}, x_{n}$.
Definition 16.5 (Extension of the definition $\prod_{j=1}^{n} x_{j}$ ). We modify the definition of " $\Pi$ " (B/G p.34, 35) as follows: Let $k \in \mathbb{Z}$ and let $\left(x_{j}\right)_{j=k}^{\infty}$ be a sequence of integers.

For each $n \in \mathbb{Z}$ such that $k \leq n$, we define an integer called $\prod_{j=k}^{n} x_{j}$ as follows. ${ }^{142}$

$$
\begin{equation*}
\text { (i) } \quad \prod_{j=k}^{k} x_{j}=x_{k}, \quad \text { (ii) } \quad \prod_{j=k}^{n+1} x_{j}=\prod_{j=k}^{n} x_{j} \cdot x_{n+1} \tag{16.9}
\end{equation*}
$$

We call $\prod_{j=k}^{n} x_{j}$ the prod of the integers $x_{k}, x_{k+1}, \ldots, x_{n-1}, x_{n}$.

### 16.5 AoP Ch.5: Underlying Notions in Set Theory

There are no addenda at this point in time.

### 16.6 AoP Ch.6: Equivalence Relations and Modular Arithmetic

### 16.6.1 Equivalence Relations

Remark 16.3. Note that B/G defines in this chapter the absolute value for integers as usual: (see def.2.11 on p.18)

$$
|m|= \begin{cases}m & \text { if } m \geqq 0 \\ -m & \text { if } m<0\end{cases}
$$

[^103]
### 16.6.2 The Division Algorithm

The following simple proposition about the absolute value of the difference of two non-negative numbers is very useful in chapters 6 and 7 of the B/G text.

Proposition 16.8. Let $m, n \in \mathbb{Z}_{\geqq}$. Then

$$
\begin{align*}
|n-m| & \leqq \max (m, n) \text {, i.e., }  \tag{16.10}\\
-\max (m, n) & \leqq n-m \leqq \max (m, n) . \tag{16.11}
\end{align*}
$$

Proof: Left as an exercise.

Corollary 16.3. Let $x, y, n \in \mathbb{Z}$ such that $0 \leqq x, y<n$. Then

$$
\begin{equation*}
-n<y-x<n . \tag{16.12}
\end{equation*}
$$

PROOF: Obvious.

### 16.6.3 The Integers Modulo $n$

In this chapter we assume that $n \in \mathbb{N}$ is fixed and that for two integers $S a S$ and $b$,

$$
a \sim b \quad \text { means } \quad a \equiv b \quad \bmod n, \quad \text { if } n \mid(a-b) .
$$

We also define $[a]^{0}:=[1]$. The following cannot be found in the $B / G$ text.
Proposition $16.9($ Arithmetic $\bmod n)$. Let $m_{1}, m_{2}, \ldots m_{k} \in \mathbb{Z}$. Then

$$
\begin{align*}
{\left[m_{1}+m_{2}+\cdots+m_{k}\right] } & =\left[m_{1}\right] \oplus\left[m_{2}\right] \oplus \cdots \oplus\left[m_{k}\right],  \tag{16.13}\\
{\left[m_{1} \cdot m_{2} \cdots m_{k}\right] } & =\left[m_{1}\right] \odot\left[m_{2}\right] \odot \cdots \odot\left[m_{k}\right],  \tag{16.14}\\
{\left[\sum_{j=1}^{k} a_{j} x^{j}\right] } & =\sum_{j=1}^{k}\left[a_{j}\right] \odot[x]^{j} . \tag{16.15}
\end{align*}
$$

PROOF: We only give the proof of (16.14). It is done by induction on the number of factors $k$. The proof of (16.13) is similar and (16.15) is a simple consequence of the two first equations.

Basis: The proof is obvious for $k=1$. We note that (16.14) is true for two factors (B/G prop.6.25 and the definition of " $a \odot b$ " that follows).

Induction assumption: We assume that (16.14) holds for some $k \in \mathbb{N}$. We then obtain for $k+1$ that

$$
\begin{aligned}
{\left[m_{1} \cdot m_{2} \cdots m_{k+1}\right] } & =\left[\left(m_{1} \cdot m_{2} \cdots m_{k}\right) \cdot m_{k+1}\right] \\
& =\left(\left[m_{1} \cdot m_{2} \cdots m_{k}\right]\right) \odot\left[m_{k+1}\right] \quad\left(B / G \text { def. of " } a \odot b^{\prime \prime}\right) \\
& =\left(\left[m_{1}\right] \odot\left[m_{2}\right] \odot \cdots \odot\left[m_{k}\right]\right) \odot\left[m_{k+1}\right] . \quad \text { (Induction assumption (16.14)) }
\end{aligned}
$$

### 16.6.4 Prime Numbers

Definition 16.6 (Prime numbers). Let $p \in \mathbb{N}, p \geqq 2$. $p$ is a prime number or $p$ is prime if $q \in \mathbb{Z}$ and $q \mid p$ implies that $q= \pm 1$ or $q= \pm p$.

We note that 1 is not prime.
Lemma 16.1. Let $p$ be prime and let $m \in \mathbb{N}$. If $p \nmid n$ ( $p$ does not divide $n$ ) then $\operatorname{gcd}(p, n)=1$.
PROOF: Let $g:=\operatorname{gcd}(p, n)$. Then $g \in \mathbb{N}, g \mid n$ and $g \mid p(B / G$ prop.6.29(i)). Only 1 and $p$ divide $p$, so either $g=p$ or $g=1$.

It follows from $p \nmid n$ and $g \mid n$ that $g \neq p$, hence $g=1$.
B/G thm 6.32: Every integer $\geqq 2$ can be factored uniquely (i.e. up to permutation) into primes.
The proof of this theorem requires the generalization of Euclid's lemma to more than two factors:
Proposition 16.10 (Euclid's Lemma for more than two factors). Let $p$ be prime and $m_{1}, m_{2}, \ldots, m_{k} \in$ $\mathbb{N}$. If $p \mid\left(m_{1} m_{2} \cdots m_{k}\right)$ then $p \mid m_{j}$ for some $1 \leqq j \leqq k$.

PROOF: Done by strong induction on the number of factors $k$.
Basis: There is nothing prove for $k=1$ and B/G prop. 6.31 (Euclid's lemma for two factors) shows the validity for $k=2$.

Induction assumption: We assume that if $p$ divides a product $n=n_{1} n_{2} \cdots n_{j}$ ) of less than $k$ factors then $p \mid n_{i}$ for some $1 \leqq i \leqq j$.

To prove that $p \mid m_{i}$ for some $1 \leqq i \leqq k$ we write $m_{1} m_{2} \cdots m_{k}=\left(m_{1} m_{2} \cdots m_{k-1}\right) m_{k}$. It follows from $B / G$ prop. 6.31 that $p \mid m_{k}$ or $p \mid\left(m_{1} m_{2} \cdots m_{k-1}\right)$. If $p$ divides $m_{k}$ then we are done. Otherwise we apply the induction assumption to the product $m_{1} m_{2} \cdots m_{k-1}$ of less than $k$ factors and obtain that there is some $1 \leqq i<k$ such that $p$ divides $m_{i}$.

Theorem 16.3 (Uniqueness of prime factorizations). Every integer $\geqq 2$ can be factored uniquely (i.e. up to permutation) into primes.

PROOF: by induction on $n$.
Base case: $n=2$ : 2 is its own and obviously unique prime factorization.
Induction assumption: assume that if $2 \leqq j<n$ then $j$ has a unique PF (up to reordering).
We now show that $n$ has a unique PF (up to reordering).
Case 1: $n$ is prime: then $n$ is the only and hence unique PF of itself.
Case 2: Else let $n=p_{1} p_{2} \cdots p_{k}=q_{1} q_{2} \cdots q_{l}$ be two PFs for $n$. $p_{1} \mid q_{1} q_{2} \cdots q_{l}$, hence $p_{1} \mid q_{j_{0}}$ for some $j_{j_{0}}$ by the generalized form of Euclid's lemma.

But then $p_{1}=q_{j_{0}}$ because $p_{1} \neq 1$ and $q_{j_{0}}$ is the only integer bigger than 1 that divides the prime $q_{j_{0}}$.
Let us reorder the $q_{j}$ in such a way that $j_{0}=1$. Then $n_{2}:=n / p_{1}=p_{2} p_{3} \cdots p_{k}=q_{2} q_{3} \cdots q_{l}$ is an integer less than $n$ and it follows from the induction assumption that $q_{2} \cdots q_{l}$ are just a reordering of $p_{2} \cdots p_{k}$.

Proposition 16.11 (Corollary to Fermat's Little Theorem). Let $p$ be prime and let $m \in \mathbb{N}$ such that $p \nmid m$. Then

$$
m^{p-1} \equiv 1 \quad \bmod p .
$$

PROOF: It follows from Fermat's Little Theorem (B/G thm.6.35) that $p \mid\left(m^{p}-m\right)$, i.e., $p \mid m\left(m^{p-1}-1\right)$. It follows from Euclid's lemma and $p \nmid m$ that $p \mid\left(m^{p-1}-1\right)$, i.e., $m^{p-1} \equiv 1 \bmod p$.

### 16.7 AoP Ch.7: Arithmetic in Base Ten

### 16.7.1 Base-Ten Representation of Integers

Proposition 16.12 (B/G Prop.7.5). Let $n, q, r \in \mathbb{N}$ such that $n=10 q+r$ and $0 \leqq r<10$ (the unique division algorithm decomposition of $n$ for division by 10.) Then $\nu(n)=\nu(q)+1$.

PROOF: We abbreviate $k:=\nu(q)$. The proof will be done as follows.
a. We prove by contradiction that $n<10^{k+1}$.
b. We prove by contradiction that $10^{k} \leqq n$.
c. The proof is done because we obtain from $\boldsymbol{a}$ and $\boldsymbol{b}$ that $10^{k} \leqq n<10^{k+1}$ and it follows from B/G prop.7.3 that $\nu(n)=k+1=\nu(q)+1$.
PROOF of $\boldsymbol{a}$.
If we had instead that $n \geqq 10^{k+1}$ then

$$
10 q=n-r \geqq 10^{k+1}-r>10^{k+1}-10=10\left(10^{k}-1\right)
$$

It follows that $q>10^{k}-1$, hence $q \geqq 10^{k}$. This contradicts $k=\nu(q)$.
PROOF of $\boldsymbol{b}$.
Assume to the contrary that $n<10^{k}$. Then

$$
10 q=n-r<10^{k}-r \leqq 10^{k} ; \text { hence } q<10^{k-1}
$$

which contradicts our assumption that $k=\nu(q)$.
Proposition 16.13 (B/G prop.7.9: Uniqueness of base-ten representation for positive integers). Let $n \in \mathbb{N}$ and assume that

$$
\begin{equation*}
n=\sum_{j=0}^{p} x_{j} 10^{j}=\sum_{j=0}^{q} y_{j} 10^{j} \tag{16.16}
\end{equation*}
$$

where $p, q \in \mathbb{Z}_{\geqq}$, each $x_{i}$ and each $y_{i}$ is a digit, $x_{p} \neq 0$ and $y_{q} \neq 0$, then $p=q$ and $x_{i}=y_{i}$ for all $i$.
PROOF: By allowing leading digits to be zero, we may assume $p=q$. Let $j^{\star}$ be the largest index $j$ such that $x_{j} \neq y_{j}$. We obtain

$$
\begin{equation*}
0=\sum_{j=0}^{p} x_{j} 10^{j}-\sum_{j=0}^{q} y_{j} 10^{j}=\sum_{j=0}^{j^{\star}}\left(x_{j}-y_{j}\right) 10^{j} . \tag{16.17}
\end{equation*}
$$

It follows from cor. 16.3 on p. 294 and B/G prop. 7.8 that

$$
\begin{equation*}
\left|\sum_{j=0}^{j^{\star}-1}\left(x_{j}-y_{j}\right) 10^{j}\right| \leqq \sum_{j=0}^{j^{\star}-1} 9 \cdot 10^{j}=10^{j^{\star}}-1 . \tag{16.18}
\end{equation*}
$$

One of $x_{j^{\star}}, y_{j^{\star}}$ is larger than the other. We may assume $x_{j^{\star}}>y_{j^{\star}}$ and hence $x_{j^{\star}}-y_{j^{\star}} \geqq 1$. We obtain

$$
\begin{aligned}
0 & =\sum_{j=0}^{j^{\star}}\left(x_{j}-y_{j}\right) 10^{j}=\left(x_{j^{\star}}-y_{j^{\star}}\right) 10^{j^{\star}}+\sum_{j=0}^{j^{\star}-1}\left(x_{j}-y_{j}\right) 10^{j} \\
& \geqq\left(x_{j^{\star}}-y_{j^{\star}}\right) 10^{j^{\star}}-\left(10^{j^{\star}}-1\right) \geqq 1 \cdot 10^{j^{\star}}-10^{j^{\star}}+1=1 .
\end{aligned}
$$

We have reached a contradiction.

Proposition 16.14 (B/G Prop.7.11). Let $n:=\sum_{j=0}^{\nu(n)-1} x_{j} 10^{j}$, where each $x_{j}$ is a digit; then

$$
\begin{equation*}
n=x_{0}+x_{1}+\cdots+x_{\nu(n)-1} \quad \bmod 3 \tag{16.19}
\end{equation*}
$$

PROOF: We define $a \sim b$ if $3 \mid b-a$. It follows from $3 \mid 9$ that $[10]=[9] \oplus[1]=[0] \oplus[1]=[1]$, hence $[10]^{j}=[1]^{j}=\left[1^{j}\right]=[1]$. But then

$$
\begin{aligned}
{[n] } & =\left[\sum_{j=0}^{\nu(n)-1} x_{j} 10^{j}\right] \\
& =\sum_{j=0}^{\nu(n)-1}\left[x_{j}\right] \odot[10]^{j} \\
& =\sum_{j=0}^{\nu(n)-1}\left[x_{j}\right] \odot[1] \\
& =\sum_{j=0}^{\nu(n)-1}\left[x_{j}\right] \\
& =\left[x_{0}+x_{1}+\cdots+x_{\nu(n)-1}\right] .
\end{aligned}
$$

It follows that $n=x_{0}+x_{1}+\cdots+x_{\nu(n)-1} \bmod 3$

### 16.7.2 The Addition Algorithm for Two Nonnegative Numbers (Base 10)

We give a simpler version of the addition algorithm than the one found in ch.7.2 of B/G.
Remark 16.4 (Addition subroutine). Given are

$$
x:=\sum_{n=0}^{K} x_{n} \cdot 10^{n}, \quad y:=\sum_{n=0}^{K} y_{n} \cdot 10^{n}
$$

in base-10 representation, i.e., $x_{n}, y_{n}$ are digits $0,1,2 \ldots, 9$. We may choose the same ending index $K$ for both $x$ and $y$ by "filling up" the number with less digits with leading zeros.

Here is the pseudocode for a subroutine, $\operatorname{Add}\left(n, x_{n}, y_{n}, z_{n}\right)$, whose task it is to compute the digits $z_{n}$ for the sum $z:=\sum_{n=0}^{K+1} z_{n} \cdot 10^{n}:=x+y$.

```
Subroutine Add( }n,\mp@subsup{x}{n}{},\mp@subsup{y}{n}{},\mp@subsup{z}{n}{})
//*
//* Inputs: n, \mp@subsup{x}{n}{},\mp@subsup{y}{n}{}
//* Output: }\mp@subsup{z}{n}{
//*
If }n=0\mathrm{ then {
    i-1 := 0;
```

$$
\begin{aligned}
& \} \\
& i_{n}:= \begin{cases}0 & \text { if } x_{n}+y_{n}+i_{n-1}<10 \\
1 & \text { if } x_{n}+y_{n}+i_{n-1} \geqq 10\end{cases} \\
& z_{n}:=\left(x_{n}+y_{n}+i_{n-1}\right)-i_{n} \cdot 10
\end{aligned}
$$

end-of-Subroutine

Note that the "sum digit" $z_{n}$ and the "carry" $i_{n}$ are associated with the "Euclidian division algorithm decomposition"

$$
x_{n}+y_{n}+i_{n-1}=10 \cdot q+r
$$

for the integer $x_{n}+y_{n}+i_{n-1}$ as follows:

$$
i_{n}=q \quad \text { and } \quad z_{n}=r .
$$

### 16.8 AoP Ch.8: Real Numbers

$B / G$ ch. 8 axiomatically defines the set $\mathbb{R}$ of all real numbers. It was stated in ch. 2.2 (Numbers) on $p .16$ that real numbers are the same as decimal numbers. A proof of this will be given in B/G ch. 12 (Decimal Expansions).

### 16.8.1 Axioms

## Remark 16.5.

Note that B/G axioms 8.1 - 8.5 and axiom 8.26 do not suffice to determine what we think of as the real numbers because the set $\mathbb{Q}$ also satisfies each one of them. Only the addition of axiom 8.52 (completeness axiom) will accomplish this.

We encountered a similar situation in the first two chapters of $B / G$ with the set $\mathbb{Z}$ of all integers where axioms $1.1-1.5$ and 2.1 also are valid for $\mathbb{Q}$ and axiom 2.15 (induction axiom) was needed to completely determine the set $\mathbb{Z}$.

The following proposition belongs between B/G prop.8.7:
Let $x, y, z \in \mathbb{R}$ and $x \neq 0$. If $x y=x z$ then $y=z ;$
and $B / G$ prop.8.8:

$$
\text { Let } x, y, z \in \mathbb{R} \text {. Then }(x+y) z=x y+x z \text {. }
$$

Proposition 16.15 (Uniqueness of multiplicative inverse). Let $a \in \mathbb{R}$ and $a \neq 0$. Then the equation $a x=1$ has $a^{-1}$ as a unique solution.

PROOF: Let $x:=a^{-1}$. Then $a x=a a^{-1}=1$, i.e., $a^{-1}$ is a solution.
Uniqueness: Let $x^{\prime} \in \mathbb{R}$ such that $a x^{\prime}=1$. It follows from $a x=1=a x^{\prime}$. and $B / G$ prop.8.7 that $x=x^{\prime}$.

The following proposition belongs at the end of B/G ch.8.1.
Proposition 16.16 (Uniqueness of $a x=b$ ). Let $a, b \in \mathbb{R}$ and $a \neq 0$. Then the equation $a x=b$ has $a$ solution $x \in \mathbb{R}$ and this solution is unique.

PROOF: Existence of a solution: Let $x:=a^{-1} b$. Then $a x=a\left(a^{-1} b\right)=\left(a a^{-1}\right) b=b, i . e ., x$ is a solution.
Uniqueness: Let $x^{\prime} \in \mathbb{R}$ such that $a x^{\prime}=b$. It follows from $0=b-b=(a x)-\left(a x^{\prime}\right)=a\left(x-x^{\prime}\right)$. Hence, from B/G prop.8.7, we obtain $x+\left(-x^{\prime}\right)=0$. It follows from the uniqueness of the additive inverse that $x=x^{\prime}$.

### 16.9 AoP Ch.9: Embedding $\mathbb{Z}$ in $\mathbb{R}$

No material so far for addenda to $B / G$ ch. 9 .

### 16.10 AoP Ch.10: Limits and Other Consequences of Completeness

Proposition 16.17 (Generalization of $\mathrm{B} / \mathrm{G}$ prop.10.5). Let $x, y \in \mathbb{R}_{>=0}$. Then
a. $x \leqq y$ if and only if $x^{2} \leqq y^{2}$,
b. $x=y$ if and only if $x^{2}=y^{2}$,
c. $x<y$ if and only if $x^{2}<y^{2}$.

PROOF of $\boldsymbol{a}, \Rightarrow)$ : Let $0 \leqq x \leqq y$. It follows from $B / G$ prop.8.37(iv) that this implies both $x \cdot x \leqq x \cdot y$ and $x \cdot y \leqq y \cdot y$. Transitivity of " $\leqq$ " ${ }^{143}$ then yields $x^{2} \leqq y^{2}$.
PROOF of $\boldsymbol{a}, \Leftarrow)$ : Let $x^{2} \leqq y^{2}$. Then $(y-x)(y+x)=y^{2}-x^{2} \geqq 0$.
Case 1: $x_{2}=y^{2}$, i.e., $(y-x)(y+x)=0$. It follows from B/G prop.8.23 that $y-x=0$ or $y+x=0$. $y-x=0$ implies $x=y$ and $y+x=0$ implies $x=y=0$ because neither $x$ nor $y$ is allowed to be negative. Either way we have $x=y$, hence $x \leqq y$ and we are done.
Case 2: $x^{2}<y^{2}$, i.e., $(y-x)(y+x)>0$. We have seen in case 1 that this necessitates $x \neq y$, and we have seen in the proof of $\boldsymbol{a}, \Rightarrow)$ that $y \leqq x$ implies $y^{2} \leqq x^{2}$, so $x^{2}<y^{2}$ implies $0 \leqq x<y$.

But then $x+y \geqq y>0$. We conclude from B/G prop.8.36 that $y-x>0$, hence $x<y$.
PROOF of $\boldsymbol{b}$ : It is obvious that $x=y$ implies $x^{2}=y^{2}$. We have seen in the proof of $\left.\boldsymbol{a}, \Leftarrow\right)$, case 1 , that $y^{2}=x^{2}$ implies $x=y$.
PROOF of $\boldsymbol{c}$ : This is a consequence of $\boldsymbol{a}$ and $\boldsymbol{b}$.
Proposition 16.18 (Alternate proof of $\mathrm{B} / \mathrm{G}$ thm.10.25). Let $A:=\left\{a \in \mathbb{R}_{>0}: a^{2}<2\right\}$. Then $r:=$ $\sup (A)$ exists and $r^{2}=2$.

PROOF: We write $A_{\text {upp } 6}$ for the set of upper bounds of $A$ (see def.8.2 on p.128).
First, we prove the existence of $r$.
It follows from $\frac{14}{10}>0$ and $\left(\frac{14}{10}\right)^{2}<2$ that $A \neq \emptyset$. It follows from $\frac{15}{10}>0$ and $\left(\frac{15}{10}\right)^{2}>2$ and $a^{2}<2$ for all $a \in A$ that $a^{2}<\left(\frac{15}{10}\right)^{2}$.

Hence, from prop.16.17 (Generalization of B/G prop.10.5), we obtain $a<\frac{15}{10}$ for all $a \in \mathbb{R}$ and we conclude that $\frac{15}{10}$ is an upper bound of $A$.

This not only proves that $r=\sup (A)$ exists, but also that $\frac{14}{10}<r \leqq \frac{15}{10}$
Let $x_{n}:=r-\frac{1}{n}$ and $y_{n}:=r+\frac{1}{n}(n \in \mathbb{N})$.
Second, we prove that $x_{n}^{2} \leqq 2$ for all $n$.

[^104]It follows from $r>\frac{14}{10}$ and $n \geqq 1$ that $x_{n}>0$. It further follows from $x_{n}<r$ that $x_{n} \notin A_{\text {upp }}$, i.e., $x_{n}<a$ for some $a \in A$, hence $x_{n}^{2}<a^{2}<2$. Third, we prove that $y_{n}^{2} \geqq 2$ for all $n$.

It follows from $r+\frac{1}{n}>r$ that $y_{n}$ is an upper bound of $A$ which is different from $r$, its least upper bound. This rules out the possibility that $y_{n} \in A$, because $r \in A_{\text {upp }}$ implies $r+\frac{1}{n} \leqq r$, a contradiction. So we have $y_{n} \notin A$, hence, as $y_{n}>0, y_{n}^{2} \geqq 2$.
We conclude the proof as follows. It follows from the rules of arithmetic for limits (prop.8.10) that

$$
\lim _{n \rightarrow \infty} y_{n}^{2}-\lim _{n \rightarrow \infty} x_{n}^{2}=\lim _{n \rightarrow \infty} \frac{4 r}{n}=4 r \cdot \lim _{n \rightarrow \infty} \frac{1}{n}
$$

and the latter expression is zero because $\lim _{n \rightarrow \infty} \frac{1}{n}=0$ (see example 8.4.a on $p$.133). So $y_{n}^{2}$ and $x_{n}^{2}$ have the same limit. We established in parts 2 and 3 of the proof that $x_{n}^{2} \leqq 2 \leqq y_{n}$. It follows from prop.8.2 on p.138 the untagged proposition that follows prop.8.11 on $p .138$ that

$$
\lim _{n \rightarrow \infty} x_{n}^{2} \leqq \lim _{n \rightarrow \infty} 2=2 \leqq \lim _{n \rightarrow \infty} y_{n}^{2}=\lim _{n \rightarrow \infty} x_{n}^{2}
$$

hence $\lim _{n \rightarrow \infty} x_{n}^{2}=2$. We use again the rules of arithmetic for limits and $\lim _{n \rightarrow \infty} \frac{1}{n}=0$ to conclude

$$
\lim _{n \rightarrow \infty} x_{n}=r+\lim _{n \rightarrow \infty} \frac{1}{n}=r, \text { hence, } 2=\lim _{n \rightarrow \infty} x_{n}^{2}=\left(\lim _{n \rightarrow \infty} x_{n}\right)^{2}=r^{2}
$$

### 16.11 AoP Ch.11: Rational and Irrational Numbers

No material so far for addenda to B/G ch.11.

### 16.12 AoP Ch.12: Decimal Expansions

## Notation Alert:

For all of this chapter let $\mathscr{D}:=\{0,1,2,3,4,5,6,7,8,9\}$ denote the set of digits.
Much of the following proposition is part of B/G prop.12.4 and prop.12.5 on p.116.
Proposition 16.19 (Geometric series for decimals).

$$
\begin{array}{ll}
\text { a. } & 0 \leqq 9 \sum_{j=n}^{\infty} 10^{-j}=\frac{1}{10^{n-1}}, \\
\text { b. } & \sum_{j=n}^{\infty} d_{j} 10^{-j} \leqq \frac{1}{10^{n-1}}, \\
\text { c. } & \sum_{j=n}^{\infty} d_{j} 10^{-j}=\frac{1}{10^{n-1}} \text { iff } d_{j}=9 \forall j \geqq n .
\end{array}
$$

PROOF of $a$ and $b$ :

$$
\sum_{j=n}^{\infty} d_{j} 10^{-j} \leq 9 \cdot \sum_{j=n}^{\infty} 10^{-j}=9 \cdot 10^{-n} \sum_{k=0}^{\infty} 10^{-k}=\frac{9}{10^{n}} \frac{1}{9 / 10}=10^{n-1}
$$

PROOF of c: For $k \geqq n$ let $s_{k}=\sum_{j=n}^{k} d_{j} 10^{-j}, \quad s=\sum_{j=n}^{\infty} d_{j} 10^{-j}, \quad s_{k}^{\prime}=9 \sum_{j=n}^{k} 10^{-j}, \quad s^{\prime}=9 \sum_{j=n}^{\infty} 10^{-j}$. If there is $j_{0} \geqq n$ such that $d_{j_{0}}<9$ then $s_{k}^{\prime}-s_{k} \geqq 10^{-j_{0}}$ for all $k \geqq n$. Hence

$$
s^{\prime}-s=\lim k \rightarrow \infty\left(s_{k}^{\prime}-s_{k}\right) \geqq 10^{-j_{0}}
$$

because the difference of the limits is the limit of the differences and according to ?? on p.??. This proves the "only if" direction. The "if" direction is immediate from part a.

Lemma 16.2 (Uniqueness of Decimal Expansions). Given is a decimal expansion $x=$ $m+\sum_{j=1}^{\infty} d_{j} 10^{-j}$. Then
a. $m \leqq x \leqq m+1$.
b. $x=m$ iff $d_{j}=0$ for all $j \in \mathbb{N}$.
c. $x=m+1$ iff $d_{j}=9$ for all $j \in \mathbb{N}$.

PROOF:
Let $s:=\sum_{j=1}^{\infty} d_{j} 10^{-j}$.
Part a is immediate from prop.16.19.a.
Part b is clear because if $d_{n}>0$ for some $n$ then $s \geqq s_{n}>d_{n} 10^{-n}>0$.
Part c follows from prop.16.19.c.
The above lemma is a generalization of the following proposition.
Proposition 16.20 (B/G prop.12.7). Let $m \in \mathbb{Z}$ and $d_{j} \in \mathscr{D}$ for all $j \in \mathbb{N}$ ) such that $m+\sum_{j=1}^{\infty} \frac{d_{j}}{10^{j}}=1$.
Then either $m=1$ and $d_{j}=0$ for all $j \in \mathbb{N}$ or $m=0$ and $d_{j}=9$ for all $j \in \mathbb{N}$. In other words, $1.00000 \ldots$ and 0.999999... are the only two decimal representations of the number 1.

PROOF: See [1] B/G ch.12.2 (Decimals).
Theorem 16.4 ( $\mathrm{B} / \mathrm{G}$ thm.12.8). Let $x \in \mathbb{R}_{\geqq 0}$ have two different decimal representations

$$
\begin{equation*}
x=d_{0}+\sum_{j=1}^{\infty} \frac{d_{j}}{10^{j}}=e_{0}+\sum_{j=1}^{\infty} \frac{e_{j}}{10^{j}} \tag{16.20}
\end{equation*}
$$

where $d_{0}, e_{0} \in \mathbb{Z}_{\geqq 0}$ and $d_{j}, e_{j} \in \mathscr{D}$ for all $j \in \mathbb{N}$.
Let $K$ be the smallest subscript such that $d_{K} \neq e_{K}$. If $d_{K}<e_{K}$ then $e_{K}=d_{K}+1, e_{j}=0$, and $d_{j}=1$ for all $j>K$.

PROOF: It follows from $d_{j}=e_{j}$ for $0 \leqq j<K$ that

$$
\begin{align*}
& \qquad 10^{K} \sum_{j=K}^{\infty} \frac{d_{j}}{10^{j}}=10^{K} \sum_{j=K}^{\infty} \frac{e_{j}}{10^{j}} \text {, i.e., } d_{K}+\sum_{j=1}^{\infty} \frac{d_{K+j}}{10^{j}}=e_{K}+\sum_{j=1}^{\infty} \frac{e_{K+j}}{10^{j}},  \tag{16.21}\\
& \text { hence } \sum_{j=1}^{\infty} \frac{d_{K+j}}{10^{j}}=\left(e_{K}-d_{K}\right)+\sum_{j=1}^{\infty} \frac{e_{K+j}}{10^{j}} .  \tag{16.22}\\
& \text { Thus } \sum_{j=1}^{\infty} \frac{d_{K+j}}{10^{j}} \leqq \sum_{j=1}^{\infty} \frac{9}{10^{j}}=1 \leqq\left(e_{K}-d_{K}\right)+\sum_{j=1}^{\infty} \frac{e_{K+j}}{10^{j}}=\sum_{j=1}^{\infty} \frac{d_{K+j}}{10^{j}} . \tag{16.23}
\end{align*}
$$

Here the first equality of (16.23) follows from prop. 16.20 (with $m=0$ ). We obtain the second inequality of (16.23) from $\left(e_{K}-d_{K}\right) \geqq 1$ and $e_{K+j} \geqq 0$. The last equality of (16.23) follows from (16.22).
(16.23) can only be valid if all inequalities are equalities. In particular we obtain

$$
1=\left(e_{K}-d_{K}\right)+\sum_{j=1}^{\infty} \frac{e_{K+j}}{10^{j}}=\sum_{j=1}^{\infty} \frac{d_{K+j}}{10^{j}}
$$

According to prop. 16.20 it follows, together with $e_{K}-d_{K} \geqq 1$, that $e_{K}-d_{K}=1, e_{K+j}=0$ for all $j \geqq 1$, and $d_{K+j}=9$ for all $j \geqq 1$.

This proves the theorem.
Remark 16.6. Note that if $x \in \mathbb{R}$ has a finite decimal expansion $x=m+\sum_{j=1}^{K} \frac{d_{j}}{10^{j}}$, i.e., one in which only finitely many digits $d_{j}$ are not zero, then $x$ is rational as the finite sum of rational terms $m$ and $d_{j} 10^{-j}$.

Corollary 16.4. If $x \in \mathbb{R}$ has different decimal expansions then $x \in \mathbb{Q}$.
PROOF: It follows from thm.16.4 above that $x$ has exactly two expansions, one of which has finitely many terms. It then follows from rem.16.6 that $x$ is rational.

The following is copied from B/G for convenience:
Definition 16.7 (Repeating Decimals). A non-negative decimal

$$
x=\left(m, d_{1}, d_{2}, \ldots\right)=m+\sum_{j=1}^{\infty} d_{j} 10^{-j} \quad\left(d_{j} \in\{0,1,2 \ldots, 9\}\right)
$$

is repeating if there are natural numbers $N$ and $p$ such that

$$
d_{N+n+k p}=d_{N+n} \forall 0 \leqq n<p, k \in \mathbb{N} .
$$

Note that the above INCLUDES $p=1$ and $d_{N}=0$ (finite expansion!).
Proposition 16.21 (B/G Prop.12.11, p.119). Every repeating decimal represents a rational number.

PROOF: Let $x=\left(m, d_{1}, d_{2}, \ldots\right)=m+\sum_{j=1}^{\infty}$ be a repeating decimal, i.e., there exist $N, p \in \mathbb{N}$ such that for all $0 \leqq n<p$ and for all $k \in \mathbb{N}$,

$$
d_{N+n+k p}=d_{N+n} .
$$

Then

$$
\begin{aligned}
x & =m+\sum_{j=1}^{N-1} d_{j} 10^{-j}+\sum_{j=N}^{\infty} d_{j} 10^{-j} \\
& =m+\sum_{j=1}^{N-1} d_{j} 10^{-j}+\sum_{k=0}^{\infty} \sum_{n=0}^{p-1} \frac{d_{N+n}}{10^{N+n+k p}} \\
& =m+\sum_{j=1}^{N-1} d_{j} 10^{-j}+\sum_{n=0}^{p-1} d_{N+n} \sum_{k=0}^{\infty} \frac{1}{10^{N+n+k p}} \\
& =m+\sum_{j=1}^{N-1} d_{j} 10^{-j}+\sum_{n=0}^{p-1} \frac{d_{N+n}}{10^{N+n}} \sum_{k=0}^{\infty} \frac{1}{10^{k p}} \\
& =m+\sum_{j=1}^{N-1} d_{j} 10^{-j}+\sum_{n=0}^{p-1} \frac{d_{N+n}}{10^{N+n}} \frac{1}{1-10^{-p}}
\end{aligned}
$$

is a finite sum of rational numbers and therefore rational.
Note 16.2 (Decimal expansions of real numbers).
Let $x \in \mathbb{R}$.
a. $x$ has at most two different decimal expansions.
b. If $x$ has two expansions then one is all zeroes except for finitely many digits and the other is all nines except for finitely many digits.
c. If $x$ has more than one expansion then $x$ is rational.
d. $x$ is a repeating decimal if and only if $x \in \mathbb{Q}$.

### 16.13 AoP Ch.13: Cardinality

Proposition 16.22 (B/G prop.13.2). Let $X, Y \neq \emptyset$ and let $f: X \xrightarrow{\sim} Y$ be a bijection. Let $a \in X$ and let

$$
\begin{equation*}
\tilde{f}: X \backslash\{a\} \rightarrow Y \backslash\{f(a)\} ; \quad x \mapsto f(x) \tag{16.24}
\end{equation*}
$$

be the restriction $\left.f\right|_{X \backslash\{a\}}$ of $f$ to $X \backslash\{a\}$ (with $f(a)$ having been removed from the original codomain). Then $\tilde{f}$ is bijective.

PROOF: This is cor.6.4 on p.109.
Proposition 16.23 (B/G prop.13.6.). Every subset of a finite set is finite.
PROOF: This is cor. 6.3 on p.109.

## 17 Other Appendices

### 17.1 Greek Letters

The following section lists all greek letters that are commonly used in mathematical texts. You do not see the entire alphabet here because there are some letters (especially upper case) which look just like our latin alphabet letters. For example: $A=$ Alpha $B=$ Beta. On the other hand there are some lower case letters, namely epsilon, theta, sigma and phi which come in two separate forms. This is not a mistake in the following tables!

| $\alpha$ alpha | $\theta$ theta | $\xi x i$ | $\phi$ phi |
| :---: | :---: | :---: | :---: |
| $\beta$ beta | $\vartheta$ theta | $\pi \quad p i$ | $\varphi$ phi |
| $\gamma$ gamma | $\iota$ iota | $\rho$ rho | $\chi$ chi |
| $\delta$ delta | $\kappa$ kappa | $\varrho$ rho | $\psi$ psi |
| $\epsilon$ epsilon | $\varkappa$ kappa | $\sigma$ sigma | $\omega$ omega |
| $\varepsilon$ epsilon | $\lambda$ lambda | $\varsigma$ sigma |  |
| $\zeta$ zeta | $\mu \quad m u$ | $\tau$ tau |  |
| $\eta$ eta | $\nu \quad n u$ | $v$ upsilon |  |
| $\Gamma$ Gamma | $\Lambda$ Lambda | $\Sigma$ Sigma | $\Psi \quad \mathrm{Psi}$ |
| $\Delta$ Delta | $\Xi X i$ | $\Upsilon$ Upsilon | $\Omega$ Omega |
| $\Theta$ Theta | $\Pi \quad \mathrm{Pi}$ | $\Phi$ Phi |  |

### 17.2 Notation

This appendix on notation has been provided because future additions to this document may use notation which has not been covered in class. It only covers a small portion but provides brief explanations for what is covered.

For a complete list check the list of symbols and the index at the end of this document.
Notations 17.1. a) If two subsets $A$ and $B$ of a space $\Omega$ are disjoint, i.e., $A \cap B=\emptyset$, then we often write $A \biguplus B$ rather than $A \cup B$ or $A+B$. Both $A^{\complement}$ and, occasionally, $\complement A$ denote the complement $\Omega \backslash A$ of $A$.
b) $\mathbb{R}_{>0}$ or $\mathbb{R}^{+}$denotes the interval $] 0,+\infty\left[, \mathbb{R}_{\geqq 0}\right.$ or $\mathbb{R}_{+}$denotes the interval $[0,+\infty[$,
c) The set $\mathbb{N}=\{1,2,3, \cdots\}$ of all natural numbers excludes the number zero. We write $\mathbb{N}_{0}$ or $\mathbb{Z}_{+}$ or $\mathbb{Z}_{\geqq 0}$ for $\mathbb{N} \biguplus\{0\}$. $\mathbb{Z}_{\geqq 0}$ is the $B / G$ notation. It is very unusual but also very intuitive.

Definition 17.1. Let $\left(x_{n}\right)_{n \in \mathbb{N}}$ be a sequence of real numbers. We call that sequence non-decreasing or increasing if $x_{n} \leqq x_{n+1}$ for all $n \in \mathbb{N}$.

We call it strictly increasing if $x_{n}<x_{n+1}$ for all $n \in \mathbb{N}$.
We call it non-increasing or decreasing if $x_{n} \geqq x_{n+1}$ for all $n$.
We call it strictly decreasing if $x_{n}>x_{n+1}$ for all $n \in \mathbb{N}$.

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## List of Symbols

$(V,\|\cdot\|) \quad$ (normed vector space), 171

- A, 131, 290
$1_{A} \quad$ (indicator function of $A$ ), 110
$2^{\Omega}, \mathfrak{P}(\Omega)$ (power set, 16
$A+b, 131,290$
$A \biguplus B$ (disjoint union), 304
$A^{\complement}$ (complement), 304
$A_{\text {Lowb }} \quad$ (lower bounds of $A$ ), 128
$A_{\text {upp }}$ (upper bounds of $A$ ), 128
$F_{0}$ (contradiction statement), 33
$N_{\varepsilon}^{A}(a) \quad$ (Trace of $N_{\varepsilon}^{A}(a)$ in $\left.A\right), 199$
$T_{0}$ (tautology statement), 33
$[a, b[] a, b$,$] \quad (half-open intervals), 18$
$[a, b]$ (closed interval), 18
$[x]_{f} \quad$ (fiber of $f$ over $f(x)$ ), 107
$\Leftrightarrow \quad$ (logical equivalence), 35
$f_{n}(\cdot) \rightarrow f(\cdot) \quad$ (pointwise convergence), 224
$f_{n}(\cdot) \xrightarrow{u c} f(\cdot) \quad$ (uniform convergence), 224
$\Rightarrow$ (implication), 37
$\|\vec{x}\|_{p} \quad\left(p\right.$-norm on $\left.\mathbb{R}^{n}\right), 171$
$\|f\| \quad$ (norm of linear $f$ ), 221
$\|f\|_{L^{2}} \quad\left(L^{2}-\right.$ norm $), 175$
$\|f\|_{L^{p}} \quad\left(L^{p}\right.$-norm on $\left.\mathscr{C}([a, b], \mathbb{R})\right), 175$
$\|x\| \bullet \quad($ Norm for $x \bullet y), 172$
$\mathfrak{P}(\Omega), 2^{\Omega} \quad$ (power set, 16
$\mathscr{U}$ (universe of discourse), 26
$\vec{x}$ vector, 93
$\bar{A} \quad$ (closure of $A$ ), 202, 203
$\bigcap\left[A_{i}: i \in I\right], 95$
$\bigcap_{i \in I} A_{i}, 95$
$\bigcup\left[A_{i}: i \in I\right], 95$
$\bigcup_{i \in I} A_{i}, 95$
$\chi_{A}$ (indicator function of $A$ ), 111
СA (complement), 304
$\exists$ (exists), 50
$\exists$ ! (exists unique), 51
$\forall$ (for all), 50
$\frac{n}{m}$ divisor, 288
$\inf \left(x_{i}\right), \inf \left(x_{i}\right)_{i \in I}, \inf _{i \in I} x_{i} \quad$ (families), 132
$\inf \left(x_{n}\right), \inf \left(x_{n}\right)_{n \in \mathbb{N}}, \inf _{n \in \mathbb{N}} x_{n} \quad$ (sequences), 132
$\inf (A) \quad$ (infimun of $A$ ), 128
$\inf _{x \in A} f(x) \quad$ (infimum of $f(\cdot)$ ), 132
$\inf _{A} f \quad($ infimum of $f(\cdot)), 132$
$\lambda^{1}, \lambda^{2}, \ldots, \lambda^{n}, \quad$ (Lebesgue measure), 283
$\leftrightarrow \quad$ (double arrow logic op.), 35
$\lim _{n \rightarrow \infty} x_{n}, 95,133,192$
$\liminf _{n \rightarrow \infty} A_{n}, 151$
$\liminf _{n \rightarrow \infty} f_{n}, 149$
$\liminf _{n \rightarrow \infty} x_{j} \quad$ (limit inferior), 141
$\limsup _{n \rightarrow \infty} A_{n}, 151$
$\limsup _{n \rightarrow \infty} f_{n}, 149$
$\limsup _{n \rightarrow \infty} x_{j} \quad$ (limit superior), 141
$\mathbb{1}_{A} \quad$ (indicator function of $A$ ), 111
$\mathbb{N}, \mathbb{N}_{0}, 304$
$\mathbb{R}^{+}, \mathbb{R}_{>0}, 304$
$\mathbb{R}_{+}, \mathbb{R}_{\geq 0}, 304$
$\mathbb{R}_{>0}, \mathbb{R}^{+}, 304$
$\mathbb{R}_{\geq 0}, \mathbb{R}_{+}, 304$
$\mathbb{Z}_{+}, \mathbb{Z}_{\geqq 0}, 304$
epi $(f)$ (epigraph), 266
$\chi=x_{n_{\star}}, x_{n_{\star+1}}, \ldots, x_{n^{\star}} \quad$ list, 93
$\max (A), \max A \quad$ ( maximum of $A$ ), 128, 290
$\min (A), \min A \quad($ minimum of $A), 128,290$
$\neg$ (negation), 30
$\partial A \quad$ (boundary of $A$ ), 191, 203
$\sup \left(x_{i}\right), \sup \left(x_{i}\right)_{i \in I}, \sup _{i \in I} x_{i} \quad$ (families), 132
$\sup \left(x_{n}\right), \sup \left(x_{n}\right)_{n \in \mathbb{N}}, \sup _{n \in \mathbb{N}} x_{n} \quad$ (sequences), 132
$\sup (A) \quad($ supremun of $A), 128$
$\sup f(x) \quad$ (supremum of $f(\cdot)), 132$
$\sup _{A} f \quad$ (supremum of $\left.f(\cdot)\right), 132$
$\rightarrow$ (arrow operator), 37
$\checkmark$ (disjunction), 35
$|X| \quad$ (size of a set), 115
$\wedge$ (conjunction), 31
] $a, b[$ (open interval), 18
$f(A)$ (direct image), 103
$f^{-1}(B) \quad$ (indirect image, preimage), 103
$g \circ f \quad$ (function composition), 74
$g \circ f(x) \quad$ (function composition), 74
$m+n \bmod 2 \quad$ (addition mod 2), 112
$n / m$ divisor, 288
$n \div m$ divisor, 288
$r^{\star}$ (rational cut), 277
$x \bullet y$ (inner product), 168
$x \bullet y \quad$ (inner product), 168
$x \diamond y \quad$ binary operation, 93
$x \in X \quad$ (element of a set, 12, 18
$x \notin X \quad$ (not an element of a set, 12
$x^{\bullet}$ unary operation, 93
$x_{n} \rightarrow-\infty, 134$
$x_{n} \rightarrow \infty, 134$
$x_{n} \rightarrow a, 95,133,192$
$\prod_{j=k}^{n} x_{j} \quad$ prod, 293
$\sum_{j=k}^{n} x_{j} \quad$ sum, 293
$(X, d(\cdot, \cdot)) \quad($ metric space $), 184$
$\left(A, d_{A \times A}\right) \quad$ (metric subspace), 198
$(\Omega, \mathfrak{F}) \quad$ (measurable space), 282
$(\Omega, \mathfrak{F}, \mu) \quad$ (measure space), 283
$(x, y)^{T} \quad$ (transpose), 76
$\left(x_{1}, x_{2}, \ldots, x_{N}\right) \quad(N$-tuple), 100
$\left(x_{1}, x_{2}, x_{3}, \ldots, x_{n}\right) \quad$ ( $n$-dimensional vector), 155
$-f(\cdot),-f \quad$ (negative function), 84
$-x \quad$ (negative of $x$ ), 161
$0(\cdot) \quad$ (zero function), 84
$A \times B \quad$ (cartesian product of 2 sets), 67
$A^{\complement}$ (complement of $A$ ), 15
$D_{f} \quad$ (natural domain of $f$ ), 70
$N_{\varepsilon}\left(x_{0}\right) \quad(\varepsilon-n e i g h b o r h o o d), 189$
$X^{I}=\prod_{i \in I} X \quad$ (cartesian product), 101
$X_{1} \times \ldots \times X_{N} \quad$ (cartesian product), 99
$[x]_{\sim},[x] \quad$ (equivalence class), 68
$\Gamma_{f}, \Gamma(f) \quad$ (graph of $f$ ), 73
$\|f\|_{\infty} \quad$ (sup-norm), 170
$\|x\| \quad$ (norm on a vector space), 171
$\mathscr{C}_{\mathscr{B}}(X, \mathbb{R}), 218$
$\mathfrak{B}$ (base of a topology), 194
$\mathfrak{F}$ ( $\sigma$-algebra), 282
$\mathfrak{N}(x) \quad$ (neighborhood system), 197
$\mathfrak{U}_{\|\cdot\|} \quad$ (norm topology), 194
$\mathfrak{U}_{d(\cdot, \cdot)} \quad$ (metric topology), 194
$\vec{x}+\vec{y} \quad$ (vector sum), 156
$\alpha \vec{x} \quad$ (scalar product), 156
$\alpha f$ (scalar product of functions), 84
$\alpha x, \quad \alpha \cdot x \quad$ (scalar product), 161
СA (complement of $A$ ), 15
$\emptyset$ (empty set), 13
$\lambda A+b \quad$ (translation/dilation in $\mathbb{R}$ ), 131
$\lambda A+b, 290$
$\mapsto \quad$ (maps to), 73
$\mathbb{N}$ (natural numbers), 16
$\mathbb{N}_{0}$ (non-negative integers), 18
Q (rational numbers), 16
$\mathbb{R}$ (real numbers), 17
$\mathbb{R}^{\star} \quad$ (non-zero real numbers), 18
$\mathbb{R}^{+} \quad$ (positive real numbers), 18
$\mathbb{R}_{>0}$ (positive real numbers), 18
$\mathbb{R}_{\geq 0} \quad$ (non-negative real numbers), 18
$\mathbb{R}_{\neq 0} \quad$ (non-zero real numbers), 18
$\mathbb{R}_{+} \quad$ (non-negative real numbers), 18
$\mathbb{Z}$ (integers), 16
$\mathbb{Z}_{\geq 0} \quad$ (non-negative integers), 18
$\mathbb{Z}_{+} \quad$ (non-negative integers), 18
$\operatorname{span}(A) \quad$ (linear span), 164
$\mu(\cdot)$ (finite measure), 283
$\mu(\cdot) \quad$ (measure), 282
$\overline{\mathbb{R}}:=\mathbb{R} \cup\{-\infty, \infty\} \quad$ (extended real numbers), 149
$\overline{\mathbb{R}}:=\mathbb{R} \cup\{-\infty, \infty\} \quad$ (extended real numbers), 282
$\overline{\mathbb{R}}_{+}$(non-negative extended), 282
$\pi_{j}(\cdot) \quad$ ( $j$ th coordinate function), 166
$\pi_{i_{1}, i_{2}, \ldots, i_{m}}(\cdot) \quad$ (m-dim projection), 166
$\mathscr{B}(X, \mathbb{R}) \quad$ (bounded real valued functions), 170
$\mathscr{C}(A, \mathbb{R}) \quad$ (cont. real valued functions on $A \subseteq \mathbb{R}$ ), 162
$\mathscr{C}(X, \mathbb{R}) \quad$ (continuous real valued functions on $X$ ), 217
$\mathscr{F}(X, \mathbb{R}) \quad$ (all real valued functions on X ), 162
$\preceq_{A}$ (partial order on subset), 69
$\prod_{i \in I} X_{i} \quad$ (cartesian product), 101
$\therefore$ (therefore), 60
$\varepsilon_{x_{0}} \quad$ (point mass), 166
$d_{A \times A} \quad$ (induced/inherited metric), 198
$d_{\|\cdot\|} \quad$ (metric induced by norm), 185
$f: X \xrightarrow{\sim} Y \quad$ (bijective function), 79
$\left.f\right|_{A} \quad$ (restriction of $f$ ), 83
$f+g \quad$ (sum of functions), 84
$f-g$ (difference of functions), 84
$f / g$ (quotient of functions), 84
$f^{-1}(\cdot) \quad$ (inverse function), 79
$f g, f \cdot g \quad$ (product of functions), 84
$x R y$ (equivalent items), 68
$x \preceq y \quad$ (precedes), 69
$x \sim y \quad$ (equivalent items), 68
$x \succeq y \quad$ (succeeds, 69
$x+y \quad$ (vector sum), 160
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xor (exclusive or), 36
\{\} (empty set), 13
$\left(x_{\imath}\right)_{\imath \in J} \quad$ (family), 86
$\left(A, \mathfrak{U}_{A}\right)$ - topol. subspace, 200
$[n]=\{1,2, \ldots, n\}, 115$
$\mathfrak{U}_{A}$ - induced/inherited topology, 200
$\mathfrak{U}_{A}$ ) - subspace topology, 200
$A \cap B \quad(A$ intersection $B), 14$
$A \backslash B \quad(A$ minus $B), 15$
$A \subset B \quad(A$ is strict subset of $B), 13$
$A \subseteq B \quad(A$ is subset of $B), 13$
$A \subsetneq B \quad(A$ is strict subset of $B), 13$
$A \triangle B \quad$ (symmetric difference of $A$ and $B$ ), 15
$A \uplus B \quad(A$ disjoint union $B), 14$
$B \supset A \quad(B$ is strict superset of $A), 13$
$B \supsetneq A \quad(B$ is strict superset of $A), 13$
$g^{-1}$ group: inverse element, 270
$(X, \mathfrak{U}) \quad$ (topological space), 194
$\left(x_{j}\right) \quad$ (sequence), 88
$\left(x_{j}\right) \quad$ (sequence), 227
$\left(x_{n_{j}}\right) \quad$ (subsequence), 88
$\liminf _{n \rightarrow \infty} A_{n} \quad$ (limit inferior for sets), 286
$\limsup A_{n} \quad$ (limit superior for sets), 286
$\sum_{k=1}^{n \rightarrow \infty} a_{k} \quad$ (series), 228
$\|\vec{v}\|_{2} \quad$ (length or Euclidean norm of $\vec{v}$ ), 157
$\mathfrak{U}$ (topology), 194
$A \cup B \quad(A$ union $B), 14$
$A \supseteq B \quad(A$ is superset of $B), 13$
$A^{o} \quad$ (interior of $A$ ), 195
$A_{n} \nearrow \bigcup_{n} A_{n}, 151$
$A_{n} \searrow \bigcap_{n} A_{n}, 151$
$\operatorname{int}(A) \quad$ (interior of $A), 195$
$x_{n} \nearrow \xi(n \rightarrow \infty), 134$
$x_{n} \searrow \xi(n \rightarrow \infty), 134$
$\operatorname{card}(X)<\operatorname{card}(Y), 115$
$\operatorname{card}(X)=\operatorname{card}(Y), 115$
$\operatorname{card}(X)>\operatorname{card}(Y), 115$
$\operatorname{card}(X) \geqq \operatorname{card}(Y), 115$
$\operatorname{card}(X) \leqq \operatorname{card}(Y), 115$
$F$ (false), 26
g.l.b.(A) (greatest lower bound of A), 128
l.u.b.(A) (least upper bound of $A$ ), 128

L/I (logically impossible), 31
$T$ (true), 26

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[^0]:    ${ }^{1}$ You will learn later in this chapter about injective functions which guarantee that if you do an operation (apply a function) to two different items then the results will also be different. If method c was restricted to only such operations then there would not be a problem. In the two "proofs" that show $-2=2$ we use operations that are not injective: In the first proof the assignment $x \mapsto 0 \cdot x$ throws everything into the same result zero. The second proof employs the assignment $x \mapsto x^{2}$ which maps two numbers $x, y$ that differ by sign only to the same squared value $x^{2}=y^{2}$.

[^1]:    ${ }^{2}$ See remark 2.2 ("Russell's Antinomy") below.

[^2]:    ${ }^{3}$ A formal definition of "if and only if" will be given in def. 3.12 on p .35 where we will also introduce the symbolic notation $P \Leftrightarrow Q$. Informally speaking, a statement is something that is either true or false.

[^3]:    ${ }^{4} \mathbb{R}$ is the set of all real numbers, i.e., the kind of numbers that make up the $x$-axis and $y$-axis in a beginner's calculus course (see remark 2.6 ("Classification of numbers") on p.16).

[^4]:    ${ }^{5}$ Neighborhoods of a point will be discussed in the chapter on the topology of $\mathbb{R}^{n}$ (see (10.6) on p.189) In short, the $\delta-$ neighborhood of $a$ is the set of all points with distance less than $\delta$ from $a$.
    ${ }^{6}$ The classification of numbers in this section is not meant to be mathematically exact. For this consult, e.g., [1] B/G (Beck/Geoghegan).

[^5]:    ${ }^{7}$ The first two chapters of [1] B/G (Beck/Geoghegan) use the "axiomatic" method to develop the mathematical structure of integers and natural numbers and give an exact proof of the induction principle.

[^6]:    ${ }^{8}$ The inverse element for $A$ in the sense of def. 13.3 on p.270. is $A$ itself!
    ${ }^{9}$ Indicator functions will be discussed in ch.6.2 on p. 110 and in ch.8.4 on p.148.

[^7]:    ${ }^{10}$ See ch.4.1 (Cartesian products and relations) on p. 67 for the real thing and examples.

[^8]:    ${ }^{11}$ usually called a proposition in a course on logic but we do not use this term as in mathematics "proposition" means a theorem of lesser importance.

[^9]:    ${ }^{12}$ A statement function is usually called a proposition function in a course on logic. As previously mentioned, we do not use the term "proposition" in this document because in most brances of mathematics it refers to a theorem of lesser importance.

[^10]:    ${ }^{13}$ To indicate that we consider $d$ as fixed but arbitrary and want to interpret " $x+7=16$ and $d$ is a dog" as a statement function of only $x$ as a variable we could have introduced the notation $B(\cdot, d): x \mapsto B(x, d)$. Similarly, to indicate that we consider $x$ as fixed but arbitrary and want to interpret " $x+7=16$ and $d$ is a $\operatorname{dog}^{\prime \prime}$ as a statement function of only $d$ as a variable we could have introduced the notation $B(x, \cdot): d \mapsto B(x, d)$. We choose not to overburden the reader with this additional notation. Rather, this situation can be handled by defining two new predicates $C: x \mapsto C(x):=" x+7=$ 16 and $z$ is a $\operatorname{dog}^{\prime \prime}$ and $D: d \mapsto D(d):=" z+7=16$ and $d$ is a $\operatorname{dog}^{\prime \prime}$ and then state that $z$ is not a variable but a fixed (but unspecified) value.
    ${ }^{14}$ This order is rather unusual in that usually you would discuss biconditional and logical equivalence operators last, but logical equivalence between two statements $A$ and $B$ is what we think of when saying " $A$ if and only if $B$ " and it helps to understand what this phrase means in the context of logic as early as possible.

[^11]:    ${ }^{15}$ The definition of a truth table will be given shortly. See def.3.8 on p.31.

[^12]:    ${ }^{16}$ e.g., if we have a statement function $P: x \mapsto P(x)$ and we look at the statements $P\left(x_{0}\right)$ for which $x_{0}$ belongs to the UoD of $P$ or a certain subset thereof

[^13]:    ${ }^{17}$ It is tough to come up with some decent examples of compound statements if the only operators at your disposal so far are negation and conjunction.

[^14]:    ${ }^{18}$ [4] Bryant, Kirby Course Notes for MAD 2104 calls this operator the equivalence operator but we abstain from that terminology because " $A$ is equivalent to $B$ " has a different meaning and is written $A \Leftrightarrow B$.

[^15]:    ${ }^{19}$ Some documents such as [4] Bryant, Kirby Course Notes for MAD 2104. also use the symbol $\oplus$.

[^16]:    ${ }^{20}$ [4] Bryant, Kirby Course Notes for MAD 2104 calls this operator the implication operator but we abstain from that terminology because " $A$ implies $B$ " has a different meaning and is written $A \Rightarrow B$.

[^17]:    ${ }^{21}$ also called the antecedent
    ${ }^{22}$ Another word for conclusion is consequent .

[^18]:    ${ }^{23}$ A counterexample is the function $f(x)=|x|$ because it is continuous everywhere but not differentiable at $x=0$.

[^19]:    ${ }^{24}$ This is theorem 3.3 (De Morgan's laws for statements).

[^20]:    ${ }^{25}$ Equivalently, " $A(x)$ is true for some $x \in \mathscr{U}$ " or " $A(x)$ is true for at least one $x \in \mathscr{U}$ ".

[^21]:    ${ }^{26}$ Equivalently, " $A(x)$ is true for exactly one $x \in \mathscr{U}$ ".

[^22]:    ${ }^{27}$ See def.10.31 ( $\varepsilon-\delta$ continuity) on p. 212.
    ${ }^{28}$ See def.10.34 (Uniform continuity of functions) on p. 220.

[^23]:    ${ }^{29}$ see remark 3.8 on p. 39 in ch.3.3.5 (Arrow and Implication Operators).
    ${ }^{30}$ We have an elementary counterexample for $n=2: 3^{2}+4^{2}=25=5^{2}$.

[^24]:    ${ }^{31}$ https://en.wikipedia.org/wiki/Euclidean_geometry\#Axioms

[^25]:    ${ }^{32}$ Note here and for the other theorems the use of $A_{2} \Rightarrow B_{2}$ instead of $A_{2} \rightarrow B_{2}$ : We assume that Thm-2 has been proved, i.e., $A_{2} \rightarrow B_{2}$ is a tautology.
    ${ }^{33}$ The expression $x^{j}$ in $\mathbf{c}$ and $\mathbf{d}$ denotes the function $x \mapsto x^{j}$.
    ${ }^{34}$ As is the case for the theorem we want to prove, note here and for Thm- 2 and Thm-3 below the use of $A_{1} \Rightarrow B_{1}$ instead of $A_{1} \rightarrow B_{1}$ : Thm-1 has been proved already, i.e., we know that $A_{1} \rightarrow B_{1}$ is a tautology.

[^26]:    ${ }^{35}$ This is example 1.8.3 in ch. 3 (Methods of Proofs) of [4] Bryant, Kirby Course Notes for MAD 2104.

[^27]:    ${ }^{36} P:=" x /(x-5)>0$ and $x \neq 5$ " if you want to be a stickler for precision
    ${ }^{37}$ The precise definition of a function will be given in section 4.2 on p.70.

[^28]:    ${ }^{38}$ See def. 9.3 on p.158. of the length or Euclidean norm of a vector in $n$-dimensional space.

[^29]:    ${ }^{39}$ Some authors, Dudley among them, do not include reflexivity into the definition of a partial ordering and then distinguish between strict partial orders and reflexive partial orders.

[^30]:    ${ }^{40}$ Real numbers were defined in section "Numbers" on p.16.
    ${ }^{41}$ Matter of fact, only for $y=2$ there exists a single argument $x$ such that $y=f(x)(x=0)$. All other $y$-values in the interval $[0,2]$ are "mapped to" by two different arguments $x= \pm \sqrt{4-y^{2}}$.

[^31]:    ${ }^{42}$ Skip this example on first reading if you do not know about functions of several variables. You will find information about this in chapter 9 ("Vectors and vector spaces") on p. 155.

[^32]:    ${ }^{43}$ We distinguish the target (codomain) $Y$ of $f(\cdot)$ from its image (range) $f(X)$ which is a subset of $Y$.
    ${ }^{44}$ or if the codomain is a subset of the complex numbers, but we won't discuss complex numbers in this document.

[^33]:    ${ }^{46}$ The ability to pick do such a selection regardless of the nature of $X, Y$ and the surjective function $f: X \rightarrow Y$ is not something one can prove. It requires acceptance of the Axiom of Choice. See remark 12.1 on p. 259 in ch. 12 (Applications of Zorn's Lemma).

[^34]:    ${ }^{47}$ Continuity of functions $y=f(x)$ with real numbers $x$ and $y$ will be defined in ch.8.2 (Convergence and Continuity in $\mathbb{R}$ ). See def.8.12 on p.139. Until then use your knowledge from calculus.

[^35]:    ${ }^{48}$ General unions and intersections will be defined in ch.5.1 (More on set operations). See def.5.1 on p.95.

[^36]:    ${ }^{49} \mathbb{Z}$ is a linearly ordered set, also called a totally ordered set set in the terminology of def. 12.1 (Linear orderings) on p. 259 .
    ${ }^{50}$ The precise definition of a series will be given in ch. 10.3 (Function Sequences and Infinite Series) on p.223.

[^37]:    ${ }^{51}$ Be sure to understand that $J \times A=J \times B$ !

[^38]:    ${ }^{52}$ Vectors can be of a more general nature than just being a finite list of numbers. See ch.9.2 (General Vector Spaces) on p. 159 (General Vector Spaces).

[^39]:    ${ }^{53}$ If $N>3$ there are many ways to group the factors of a cartesian product. For $N=4$ there already are 3 times as many possibilities as for $N=3$ :

    $$
    X_{1} \times\left(X_{2} \times X_{3} \times X_{4}\right), \quad\left(X_{1} \times X_{2}\right) \times\left(X_{3} \times X_{4}\right), \quad X_{1} \times\left(X_{2} \times X_{3} \times X_{4}\right)
    $$

    Actually proving that we can group the sets with parentheses any way we like is very tedious and will not be done in this document.

[^40]:    ${ }^{54}$ See def.4.20 (Vectors and finite lists) on p. 93 .

[^41]:    ${ }^{55}$ In other words, if and only if those two families are equal in the sense of def.4.18 on p.87.

[^42]:    ${ }^{56}$ The range $f(X)$ of $f$ (see (4.9) on p. 73 is a special case of a direct image.

[^43]:    ${ }^{57}[x]_{f}$ is called the fiber over $f(x)$ of the function $f$.

[^44]:    ${ }^{58}$ This is B/G [1] prop.13.2.

[^45]:    ${ }^{59}$ Some authors call this characteristic function of $A$ and some choose to write $\chi_{A}$ or $\mathbb{1}_{A}$ instead of $1_{A}$.
    ${ }^{60}$ See remark 5.4 on p.101, ch.5.2 (Cartesian Products of more than Two Sets).

[^46]:    ${ }^{61}$ Understanding this proof is not very important and you will understand the essence of it if you read instead the subsequent remark 7.2.

[^47]:    ${ }^{62}$ Excluding the digit 9 avoids dealing with numbers which have non-unique decimal representations such as $1 / 2=0.5 \overline{0}=0.4 \overline{9}$.

[^48]:    ${ }^{63}$ We had to list the case $X=Y$ separately to capture the case $X=Y=\emptyset$ because functions, hence bijections, are not defined for empty domains or codomains.

[^49]:    64 "Cardinality as a Partial Ordering". The proof of the Cantor-Schröder-Bernstein Theorem is by no means trivial.

[^50]:    ${ }^{65}$ The definitions are given for $\mathbb{Z}$ in def.16.2 on p.290.

[^51]:    ${ }^{66}$ Note that if $A$ is not bounded above then the same holds for the superset $B$ and that if $B$ is empty then $A$ is empty.
    ${ }^{67}$ See also def.16.3 in ch.16.2.3

[^52]:    ${ }^{68}$ (8.22)can also be deduced from (8.21) and the fact that $\inf \{\varphi(u): u \in A\}=-\sup \{-\varphi(v): v \in A\}$.
    ${ }^{69}$ We will define convergence of a sequence of items more general than real numbers in ch.10.1.4 (see def.10.11 (convergence of sequences in metric spaces) on p.192).

[^53]:    ${ }^{70}$ See [1] B/G (Beck/Geoghegan) prop.10.23

[^54]:    ${ }^{71}$ The constant sequence $(x y)$ has limit $x y$ according to part a

[^55]:    ${ }^{72}$ We call such a function sequence continuous in def. 10.30 (Sequence continuity) on p .212 where continuity is generalized to metric spaces.

[^56]:    ${ }^{73}$ See def.9.6 (linear combinations) on p. 164
    ${ }^{74}$ See thm. 10.13 on p. 213 .

[^57]:    ${ }^{75}$ Here is a shorter proof of " $\Leftarrow$ ". The drawback: It makes use of prop.8.13 and prop.8.14. Let $L:=\liminf _{n \rightarrow \infty} x_{n}=\limsup _{n \rightarrow \infty} x_{n}$. Let $\varepsilon>0$. We know from (8.54), p. 142 and (8.57), p. 144 that $L+\varepsilon / 2 \notin \mathscr{U}$ and $L-\varepsilon / 2 \notin \mathscr{L}$ But then there are at most finitely many $n$ for which $x_{n}$ has a distance from $L$ which exceeds $\varepsilon / 2$. Let $N$ be the maximum of those $n$. It follows that $\left|x_{n}-L\right|<\varepsilon$ for all $n>N$, hence $L=\lim _{n \rightarrow \infty} x_{n}$.

[^58]:    ${ }^{76}$ In more advanced texts you will find the following

[^59]:    ${ }^{77}$ See note 8.1 on p. 134 .

[^60]:    ${ }^{78}$ The definition of an abstract vector space will be given in def.9.4 on p.160.
    ${ }^{79}$ See def.5.4 (Cartesian Product of three or more sets) on $p .99$ ) concerning $\mathbb{R}^{n}=\underbrace{\mathbb{R} \times \mathbb{R} \times \cdots \times \mathbb{R}}_{n \text { times }}$. .

[^61]:    ${ }^{80}$ You will learn in the chapter on vector spaces that the vector $\vec{a}=(4,-3)$ is the projection on the $x y$-coordinates $\pi_{1,2}(\cdot): \mathbb{R}^{3} \rightarrow \mathbb{R}^{2} ; \quad(x, y, z) \mapsto(x, y)$ of the vector $\vec{v}=(4,-3,12)$. See Example 9.19) on p.166.

[^62]:    ${ }^{81}$ Note that $\mathscr{F}(X, \mathbb{R})=\mathbb{R}^{X}$ (see remark 5.4, p. 101 which follows def. 5.5 of the Cartesian Product of a family of sets.)
    ${ }^{82}$ Continuity for such functions was discussed in ch.8.2 on p.133.

[^63]:    ${ }^{84}$ also called dot product, e.g., in [3] Brin/Marchesi Linear Algebra, ch.6, Orthogonality.

[^64]:    ${ }^{85}$ see def. 8.6 (bounded functions) on p. 131

[^65]:    ${ }^{86}$ Actually, the proof that $\|\cdot\|_{2}$ satisfies the triangle inequality was given only for dimensions $1,2,3$. It will be proved in this chapter that it is true for all dimensions $n$. See cor.9.1 (Inner products define norms) on p.172.

[^66]:    ${ }^{87}$ If you forgot the meaning of $X \times X$, it's time to review the definition of a cartesian product (def.4.1 in ch.4.1 on p. 67.
    ${ }^{88}$ see def. 10.3 on p. 186 and prop. 10.2 directly thereafter.

[^67]:    ${ }^{89}$ Compare this to def.9.16 (Norm for an inner product) on p.172.

[^68]:    ${ }^{90}$ A shorter proof can be given if the previous proposition is used.

[^69]:    ${ }^{91}$ by the way, this is the exact spot where the proof breaks down if you deal with an infinite intersection of open sets: the minimum would have to be replaced by an infimum and there is no guarantee that it would be strictly larger than zero.

[^70]:    ${ }^{92}$ You could have used $N_{1}+N_{2}$ instead. Do you see why?
    ${ }^{93}$ See prop. 8.6 on p. 135.

[^71]:    ${ }^{94}$ See prop. 8.7 on p. 135.
    ${ }^{95}$ We encountered subsets of $2^{X}$ with special properties previously when looking at rings of sets in def.2.13 (Rings and Algebras of Sets) on p. 21 .

[^72]:    ${ }^{96}$ Source: https:/ / en.wikipedia.org/wiki/Interior_(topology)

[^73]:    ${ }^{97}$ German: Berührungspunkt - see [12] Von Querenburg, p. 21

[^74]:    ${ }^{98} a \in A \subseteq(X, d)$ is called an isolated point of $A$ if there is a neighborhood $U$ of $a$ such that $U \cap A=\{a\}$, i.e., the "punctured neighborhood" $U \backslash\{a\}$ of $a$ has empty intersection with $a$.

[^75]:    ${ }^{99}$ see def.10.25, def. 10.26 and def. 10.10 for the corresponding definitions in metric spaces.

[^76]:    ${ }^{102}$ Remember the completeness axiom for $\mathbb{R}$ (axiom 8.1 on p .128) which states that any subset $A$ of $\mathbb{R}$ which possesses upper bounds has a least upper bound (the supremum $\sup (A)$ ). This axiom was needed to establish the validity of thm. 8.3 (Characterization of limits via limsup and liminf) on p .147 , a theorem which will be used in this chapter to prove the completeness of $\mathbb{R}$ as a metric space.
    ${ }^{103}$ See ch.8.1 (Minima, Maxima, Infima and Suprema).

[^77]:    ${ }^{105}$ Actually, we have proved a very strong form of continuity. Generally speaking, $\delta=\delta\left(\varepsilon, x_{0}\right)$ is tailored not only to the given $\varepsilon$, but also to the particular argument $x_{0}$ at which continuity needs to be verified. We were able to find $\delta$ which does not depend on the argument $x_{0}$ but only on $\varepsilon$. We will learn later that this makes $i d_{X}$ uniformly continuous on its domain $X$. See def.10.34 (Uniform continuity of functions) on p.220.

[^78]:    ${ }^{106}$ See def.9.6 (linear combinations) on p. 164

[^79]:    ${ }^{107}$ Here $f(\cdot)$ is a function, not necessarily continuous, on the unit interval $[0,1]$. The binomial coefficient $\binom{n}{k}$ is defined as $\frac{n!}{k!(n-k)!}$ where $0!=1$ and $n!=1 \cdot 2 \cdot 3 \cdots n$ for $n \in!\mathbb{N}$ (see ch. 4 of [1] B/G Art of Proof)

[^80]:    ${ }^{108}$ see chapter 11.6 (Continuous Functions and Compact Spaces) on p. 255
    ${ }^{109}$ For the special case of $(X, d)=\left(\mathbb{R}, d_{|\cdot|}\right)$ where $d_{|\cdot|}(x, y)=|y-x|$, see [1] Beck/Geoghegan, Appendix A.3, "Uniform continuity".
    ${ }^{110}$ Well, that's not entirely accurate: There is a notion of "uniform spaces" which generalize the concept of a metric but are less general than topological spaces and there is a notion of uniform continuity for those sets.

[^81]:    ${ }^{111}$ Note that we use the same notation $\|\cdot\|$ for both the norm on $V$ and the norm of the linear function $f$. Do not confuse the two!

[^82]:    ${ }^{112}$ I must confess that " $f_{n}(\cdot) \xrightarrow{u c} f(\cdot)$ " is a notation that I coined myself because it is not as tedious as writing " $f_{n}(\cdot) \rightarrow f(\cdot)$ uniformly"

[^83]:    ${ }^{113}$ This is a generalization of [1] B/G (Beck/Geoghegan) prop.12.3, p.115.

[^84]:    ${ }^{114}$ We could also have concluded as follows: $\left|s_{j}-s_{j-1}\right|=1$ implies that the Cauchy formulation of the convergence criteria for series (see (10.92a) on p.228) is not satisfied, hence no convergence of the series.
    ${ }^{115}$ This notation was copied from chapter 13 (Cardinality) of [1] B/G (Beck/Geoghegan). It has nothing to do with equivalence classes!

[^85]:    ${ }^{116}$ see cor. 10.4 on p. 238
    ${ }^{117}$ This was proved by the German mathematician Peter Gustav Lejeune Dirichlet (1805-1859).
    ${ }^{118} \sum a_{k}$ converges according to prop.10.30.

[^86]:    ${ }^{119}$ This was proved by the German mathematician Bernhard Riemann (1826-1866).

[^87]:    ${ }^{120}$ That is as intuitive a grid as you can think of, especially if you look at the 2-dimensional plane or 3-dimensional space.

[^88]:    ${ }^{121}$ See def.4.19 on p.87.

[^89]:    ${ }^{122}$ Note that it is not claimed that there would be infinitely many different points $x_{n_{j}}$, only infinitely many indices $n_{j}$. Indeed, what would you do if the original Cauchy sequence was chosen to be $x_{1}=x_{2}=\cdots=a$ for some $a \in A$ ?

[^90]:    ${ }^{123}$ see def. 5.3 on p. 97
    ${ }^{124} \mathrm{We}$ had to remove the example of the indiscrete topology because this topology does not come from a metric.

[^91]:    ${ }^{125}$ We could have written more concisely $F:=\left\{x_{j}: j \in \mathbb{N}\right\}$ but the above definition was chosen to remind you that $F$ does not contain any duplicates. Note that $F$ can be very small even if there are infinitely many indices $j$ : If $x_{j}=(-1)^{j}$ then $F=\{-1,1\}$ only contains two elements!
    ${ }^{126}$ The number $\lambda=2 \rho$ is called a Lebesgue number of $\left(U_{i}\right)_{i \in I}$. In other words, the Lebesgue number is the diameter of the $\rho$-neighborhoods. Note that if $\lambda$ is a Lebesgue number of an open covering then any $\lambda^{\prime}$ which satisfies $0<\lambda^{\prime}<\lambda$ also is a Lebesgue number. of that same cover.

[^92]:    ${ }^{127}$ See def. 10.34 on p. 220.

[^93]:    ${ }^{128}$ Axiom of Choice: Let $\Omega$ be any nonempty set. Then there is a function $f: 2^{\Omega} \backslash\{\emptyset\} \rightarrow \Omega$ such that $f(A) \in A$ for each $A \in 2^{\Omega} \backslash\{\emptyset\}$, i.e., for each nonempty $A \subseteq \Omega . f$ is referred to as a choice function because it chooses from any nonempty subset $A$ one of its elements.

[^94]:    ${ }^{129}$ see def. 12.2 on p. 259.

[^95]:    ${ }^{130}$ See [7] Haaser/Sullivan: Real Analysis.

[^96]:    ${ }^{131}$ In other words, we have a function $\diamond: S \times S \rightarrow S,(s, t) \mapsto \diamond(s, t):=s \diamond t$ in the sense of def.4.6 (Mappings (functions)) on p. 72 .
    ${ }^{132}$ If this is too abstract for you, choose $A:=\mathbb{R}$, the set of real numbers. Then the elements of $S$ will be functions such as $f(x)=3 x^{2}$ and $g(x)=7 x+5 e^{x}$. for those two specific functions you get $g \circ f(x)=g(f(x))=g\left(3 x^{2}\right)=21 x^{2}+5 e^{3 x^{2}}$.
    ${ }^{133}$ See def.4.7 (Function composition) on p.74.

[^97]:    ${ }^{134}$ named so after the Norwegian mathematician Niels Henrik Abel who lived in the first half of the 19th century and died at age 26.

[^98]:    ${ }^{135}$ They then of course also are bijective functions (see def. 4.9).

[^99]:    ${ }^{136}$ This was shown for multiplication $\left[\left(p_{1}, q_{1}\right)\right] \cdot\left[\left(p_{2}, q_{2}\right)\right]=\left[\left(p_{1} p_{2}, q_{1} q_{2}\right)\right]$ in exercise 4.9 on p.91.

[^100]:    ${ }^{137}$ This is not entirely correct: we must demand that the measure is $\sigma$-finite, i.e., there are measurable sets with finite measure whose union is the entire space. Such is the case for Lebesgue measure: Let $A_{k}:=[-k, k]^{n}$. The union of those sets is $\mathbb{R}^{k}$ and $\lambda^{n}\left(A_{K}\right)=(2 k)^{n}<\infty$.

[^101]:    ${ }^{138}$ The definitions given here also apply to subsets of $\mathbb{Q}$ and $\mathbb{R}$. See def.8.1 on p.128.
    ${ }^{139}$ See the optional chapter 13 (Algebraic structures) on p. 269

[^102]:    ${ }^{140}$ See ch.3.3.2, def.3.9 on p.31.

[^103]:    $141 \sum_{j=k}^{n} x_{j}$ can also be written $x_{k}+x_{k+1}+\cdots+x_{n}$.
    $142 \prod_{j=k}^{n} x_{j}$ can also be written $x_{k} \cdot x_{k+1} \cdots x_{n}$ or $x_{k} x_{k+1} \cdots x_{n}$

[^104]:    ${ }^{143}$ Proved for integers in prop16.3 on p. 289.

