Calculus I
for Computer Science and Statistics Students

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Lecture Notes

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1 Foundations: Mathematical Logic and Set Theory

1.1 Introductory Remarks

The task of mathematics is to establish the truth or falsehood of (formalizable) statements using rigorous logic, and to provide methods for the solution of classes of (e.g. applied) problems, ideally including rigorous logical proofs verifying the validity of the methods (proofs that the method under consideration will, indeed, provide a correct solution).

The topic of this class is calculus, which is short for infinitesimal calculus, usually understood (as it is here) to mean differential and integral calculus of real and complex numbers (more generally, calculus may refer to any method or system of calculation guided by the symbolic manipulation of expressions, we will briefly touch on another example in Sec. 1.2 below). In that sense, calculus is the beginning part of the broader field of (mathematical) analysis, the section of mathematics concerned with the notion of a limit (for us, the most important examples will be limits of sequences (Def. 7.1 below) and limits of functions (Def. 8.17 below)).

Before we can properly define our first limit, however, it still needs some preparatory work. In modern mathematics, the objects under investigation are almost always so-called sets. So one aims at deriving (i.e. proving) true (and interesting and useful) statements about sets from other statements about sets known or assumed to be true. Such a derivation or proof means applying logical rules that guarantee the truth of the derived (i.e. proved) statement.

However, unfortunately, a proper definition of the notion of set is not easy, and neither is an appropriate treatment of logic and proof theory. Here, we will only be able to briefly touch on the bare necessities from logic and set theory needed to proceed to the core matter of this class. We begin with logic in Sec. 1.2, followed by set theory in Sec. 1.3, combining both in Sec. 1.4. The interested student can find an introductory presentation of axiomatic set theory in [Phi16, Sec. A] and he/she should consider taking a separate class on set theory, logic, and proof theory at a later time.

1.2 Propositional Calculus

1.2.1 Statements

Mathematical logic is a large field in its own right and, as indicated above, a thorough introduction is beyond the scope of this class – the interested reader may refer to [EFT07], [Kun12], and references therein. Here, we will just introduce some basic concepts using common English (rather than formal symbolic languages – a concept explained in books like [EFT07]).

As mentioned before, mathematics establishes the truth or falsehood of statements. By a statement or proposition we mean any sentence (any sequence of symbols) that can
reasonably be assigned a truth value, i.e. a value of either true, abbreviated T, or false, abbreviated F. The following example illustrates the difference between statements and sentences that are not statements:

**Example 1.1. (a) Sentences that are statements:**

Every dog is an animal. (T)
Every animal is a dog. (F)
The number 4 is odd. (F)
$2 + 3 = 5$. (T)
$\sqrt{2} < 0$. (F)
$x + 1 > 0$ holds for each natural number $x$. (T)

**Example 1.1. (b) Sentences that are not statements:**

Let’s study calculus!
Who are you?
$3 \cdot 5 + 7$.
$x + 1 > 0$.
All natural numbers are green.

The fourth sentence in Ex. 1.1(b) is not a statement, as it can not be said to be either true or false without any further knowledge on $x$. The fifth sentence in Ex. 1.1(b) is not a statement as it lacks any meaning and can, hence, not be either true or false. It would become a statement if given a definition of what it means for a natural number to be green.

### 1.2.2 Logical Operators

The next step now is to combine statements into new statements using logical operators, where the truth value of the combined statements depends on the truth values of the original statements and on the type of logical operator facilitating the combination.

The simplest logical operator is **negation**, denoted $\neg$. It is actually a so-called unary operator, i.e. it does not combine statements, but is merely applied to one statement. For example, if $A$ stands for the statement “Every dog is an animal.”, then $\neg A$ stands for the statement “Not every dog is an animal.”; and if $B$ stands for the statement “The number 4 is odd.”, then $\neg B$ stands for the statement “The number 4 is not odd.”, which can also be expressed as “The number 4 is even.”

To completely understand the action of a logical operator, one usually writes what is known as a truth table. For negation, the truth table is

<table>
<thead>
<tr>
<th>$A$</th>
<th>$\neg A$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>F</td>
</tr>
<tr>
<td>F</td>
<td>T</td>
</tr>
</tbody>
</table>

(1.1)
that means if the input statement $A$ is true, then the output statement $\neg A$ is false; if
the input statement $A$ is false, then the output statement $\neg A$ is true.

We now proceed to discuss binary logical operators, i.e. logical operators combining
precisely two statements. The following four operators are essential for mathematical
reasoning:

Conjunction: $A$ and $B$, usually denoted $A \land B$.

Disjunction: $A$ or $B$, usually denoted $A \lor B$.

Implication: $A$ implies $B$, usually denoted $A \implies B$.

Equivalence: $A$ is equivalent to $B$, usually denoted $A \iff B$.

Here is the corresponding truth table:

<table>
<thead>
<tr>
<th>$A$</th>
<th>$B$</th>
<th>$A \land B$</th>
<th>$A \lor B$</th>
<th>$A \implies B$</th>
<th>$A \iff B$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
<td>T</td>
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</tbody>
</table>

(1.2)

When first seen, some of the assignments of truth values in (1.2) might not be completely
intuitive, due to the fact that logical operators are often used somewhat differently in
common English. Let us consider each of the four logical operators of (1.2) in sequence:

For the use in subsequent examples, let $A_1, \ldots, A_6$ denote the six statements from Ex.
1.1(a).

Conjunction: Most likely the easiest of the four, basically identical to common language
use: $A \land B$ is true if, and only if, both $A$ and $B$ are true. For example, using Ex. 1.1(a),
$A_1 \land A_4$ is the statement “Every dog is an animal and $2 + 3 = 5.”$, which is true since
both $A_1$ and $A_4$ are true. On the other hand, $A_1 \land A_3$ is the statement “Every dog is
an animal and the number 4 is odd.”, which is false, since $A_3$ is false.

Disjunction: The disjunction $A \lor B$ is true if, and only if, at least one of the statements
$A, B$ is true. Here one already has to be a bit careful – $A \lor B$ defines the inclusive or,
whereas “or” in common English is often understood to mean the exclusive or (which is
false if both input statements are true). For example, using Ex. 1.1(a), $A_1 \lor A_4$ is the statement
“Every dog is an animal or $2 + 3 = 5.”$, which is true since both $A_1$ and $A_4$ are true. The statement
$A_1 \lor A_3$, i.e. “Every dog is an animal or the number 4 is odd.”
is also true, since $A_1$ is true. However, the statement $A_2 \lor A_5$, i.e. “Every animal is a
dog or $\sqrt{2} < 0.$” is false, as both $A_2$ and $A_5$ are false.

As you will have noted in the above examples, logical operators can be applied to
combine statements that have no obvious contents relation. While this might seem
strange, introducing contents-related restrictions is unnecessary as well as undesirable,
since it is often not clear which seemingly unrelated statements might suddenly appear
in a common context in the future. The same occurs when considering implications and
equivalences, where it might seem even more obscure at first.
Implication: Instead of $A$ implies $B$, one also says *if $A$ then $B$*, $A$ is a consequence of $A$, $B$ is concluded or inferred from $A$, $A$ is sufficient for $B$, or $B$ is necessary for $A$. The implication $A \Rightarrow B$ is always true, except if $A$ is true and $B$ is false. At first glance, it might be surprising that $A \Rightarrow B$ is defined to be true for $A$ false and $B$ true, however, there are many examples of incorrect statements implying correct statements. For instance, squaring the (false) equality of integers $-1 = 1$, implies the (true) equality of integers $1 = 1$. However, as with conjunction and disjunction, it is perfectly valid to combine statements without any obvious context relation: For example, using Ex. 1.1(a), the statement $A_1 \Rightarrow A_6$, i.e. “Every dog is an animal implies $x + 1 > 0$ holds for each natural number $x$.” is true, since $A_6$ is true, whereas the statement $A_4 \Rightarrow A_2$, i.e. “$2 + 3 = 5$ implies every animal is a dog.” is false, as $A_4$ is true and $A_2$ is false.

Of course, the implication $A \Rightarrow B$ is not really useful in situations, where the truth values of both $A$ and $B$ are already known. Rather, in a typical application, one tries to establish the truth of $A$ to prove the truth of $B$ (a strategy that will fail if $A$ happens to be false).

**Example 1.2.** Suppose we know Sasha to be a member of a group of children. Then the statement $A$ “Sasha is a girl.” implies the statement $B$ “There is at least one girl in the group.” A priori, we might not know if Sasha is a girl or a boy, but if we can establish Sasha to be a girl, then we also know $B$ to be true. If we find Sasha to be a boy, then we do not know whether $B$ is true or false.

Equivalence: $A \Leftrightarrow B$ means $A$ is true if, and only if, $B$ is true. Once again, using input statements from Ex. 1.1(a), we see that $A_1 \Leftrightarrow A_4$, i.e. “Every dog is an animal is equivalent to $2 + 3 = 5.$”, is true as well as $A_2 \Leftrightarrow A_3$, i.e. “Every animal is a dog is equivalent to the number 4 is odd.”. On the other hand, $A_4 \Leftrightarrow A_5$, i.e. “$2 + 3 = 5$ is equivalent to $\sqrt{2} < 0$, is false.

Analogous to the situation of implications, $A \Leftrightarrow B$ is not really useful if the truth values of both $A$ and $B$ are known a priori, but can be a powerful tool to prove $B$ to be true or false by establishing the truth value of $A$. It is obviously more powerful than the implication as illustrated by the following example (compare with Ex. 1.2):

**Example 1.3.** Suppose we know Sasha is the tallest member of a group of children. Then the statement $A$ “Sasha is a girl.” is equivalent to the statement $B$ “The tallest kid in the group is a girl.” As in Ex. 1.2, if we can establish Sasha to be a girl, then we also know $B$ to be true. However, in contrast to Ex. 1.2, if we find Sasha to be a boy, we know $B$ to be false.

**Remark 1.4.** In computer science, the truth value $T$ is often coded as 1 and the truth value $F$ is often coded as 0.

### 1.2.3 Rules

Note that the expressions in the first row of the truth table (1.2) (e.g. $A \land B$) are not statements in the sense of Sec. 1.2.1, as they contain the *statement variables* (also known
as propositional variables) \( A \) or \( B \). However, the expressions become statements if all statement variables are substituted with actual statements. We will call expressions of this form propositional formulas. Moreover, if a truth value is assigned to each statement variable of a propositional formula, then this uniquely determines the truth value of the formula. In other words, the truth value of the propositional formula can be calculated from the respective truth values of its statement variables – a first justification for the name propositional calculus.

**Example 1.5.** (a) Consider the propositional formula \((A \land B) \lor (\neg B)\). Suppose \( A \) is true and \( B \) is false. The truth value of the formula is obtained according to the following truth table:

\[
\begin{array}{c|c|c|c|c}
A & B & A \land B & \neg B & (A \land B) \lor (\neg B) \\
T & F & F & T & T \\
\end{array}
\]  
(1.3)

(b) The propositional formula \( A \lor (\neg A) \), also known as the law of the excluded middle, has the remarkable property that its truth value is \( T \) for every possible choice of truth values for \( A \):

\[
\begin{array}{c|c|c}
A & \neg A & A \lor (\neg A) \\
T & F & T \\
F & T & T \\
\end{array}
\]  
(1.4)

Formulas with this property are of particular importance.

**Definition 1.6.** A propositional formula is called a tautology or universally true if, and only if, its truth value is \( T \) for all possible assignments of truth values to all the statement variables it contains.

**Notation 1.7.** We write \( \phi(A_1, \ldots, A_n) \) if, and only if, the propositional formula \( \phi \) contains precisely the \( n \) statement variables \( A_1, \ldots, A_n \).

**Definition 1.8.** The propositional formulas \( \phi(A_1, \ldots, A_n) \) and \( \psi(A_1, \ldots, A_n) \) are called equivalent if, and only if, \( \phi(A_1, \ldots, A_n) \iff \psi(A_1, \ldots, A_n) \) is a tautology.

**Lemma 1.9.** The propositional formulas \( \phi(A_1, \ldots, A_n) \) and \( \psi(A_1, \ldots, A_n) \) are equivalent if, and only if, they have the same truth value for all possible assignments of truth values to \( A_1, \ldots, A_n \).

**Proof.** If \( \phi(A_1, \ldots, A_n) \) and \( \psi(A_1, \ldots, A_n) \) are equivalent and \( A_i \) is assigned the truth value \( t_i \), \( i = 1, \ldots, n \), then \( \phi(A_1, \ldots, A_n) \iff \psi(A_1, \ldots, A_n) \) being a tautology implies it has truth value \( T \). From (1.2) we see that either \( \phi(A_1, \ldots, A_n) \) and \( \psi(A_1, \ldots, A_n) \) both have truth value \( T \) or they both have truth value \( F \).

If, on the other hand, we know \( \phi(A_1, \ldots, A_n) \) and \( \psi(A_1, \ldots, A_n) \) have the same truth value for all possible assignments of truth values to \( A_1, \ldots, A_n \), then, given such an assignment, either \( \phi(A_1, \ldots, A_n) \) and \( \psi(A_1, \ldots, A_n) \) both have truth value \( T \) or both have truth value \( F \), i.e. \( \phi(A_1, \ldots, A_n) \iff \psi(A_1, \ldots, A_n) \) has truth value \( T \) in each case, showing it is a tautology. ■
For all logical purposes, two equivalent formulas are exactly the same – it does not matter if one uses one or the other. The following theorem provides some important equivalences of propositional formulas. As too many parentheses tend to make formulas less readable, we first introduce some precedence conventions for logical operators:

Convention 1.10. \( \neg \) takes precedence over \( \land, \lor \), which take precedence over \( \Rightarrow, \Leftrightarrow \).

So, for example,

\[(A \lor \neg B \Rightarrow \neg B \land \neg A) \Leftrightarrow \neg C \land (A \lor \neg D)\]

is the same as

\[
\left( (A \lor \neg B) \Rightarrow ((\neg B) \land (\neg A)) \right) \Leftrightarrow \left( (\neg C) \land (A \lor (\neg D)) \right).
\]

Theorem 1.11. (a) \((A \Rightarrow B) \Leftrightarrow \neg A \lor B\). This means one can actually define implication via negation and disjunction.

(b) \((A \Leftrightarrow B) \Leftrightarrow ((A \Rightarrow B) \land (B \Rightarrow A))\), i.e. \(A\) and \(B\) are equivalent if, and only if, \(A\) is both necessary and sufficient for \(B\). One also calls the implication \(B \Rightarrow A\) the converse of the implication \(A \Rightarrow B\). Thus, \(A\) and \(B\) are equivalent if, and only if, both \(A \Rightarrow B\) and its converse hold true.

(c) Commutativity of Conjunction: \(A \land B \Leftrightarrow B \land A\).

(d) Commutativity of Disjunction: \(A \lor B \Leftrightarrow B \lor A\).

(e) Associativity of Conjunction: \((A \land B) \land C \Leftrightarrow A \land (B \land C)\).

(f) Associativity of Disjunction: \((A \lor B) \lor C \Leftrightarrow A \lor (B \lor C)\).

(g) Distributivity I: \(A \land (B \lor C) \Leftrightarrow (A \land B) \lor (A \land C)\).

(h) Distributivity II: \(A \lor (B \land C) \Leftrightarrow (A \lor B) \land (A \lor C)\).

(i) De Morgan’s Law I: \(\neg(A \land B) \Leftrightarrow \neg A \lor \neg B\).

(j) De Morgan’s Law II: \(\neg(A \lor B) \Leftrightarrow \neg A \land \neg B\).

(k) Double Negative: \(\neg\neg A \Leftrightarrow A\).

(l) Contraposition: \((A \Rightarrow B) \Leftrightarrow (\neg B \Rightarrow \neg A)\).

Proof. Each equivalence is proved by providing a truth table and using Lem. 1.9.

(a):

<table>
<thead>
<tr>
<th>A</th>
<th>B</th>
<th>\neg A</th>
<th>A \Rightarrow B</th>
<th>\neg A \lor B</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
<td>F</td>
<td>T</td>
<td>T</td>
</tr>
<tr>
<td>T</td>
<td>F</td>
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<td>F</td>
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</tbody>
</table>

(b) – (h): Exercise.
(i):

<table>
<thead>
<tr>
<th>$A$</th>
<th>$B$</th>
<th>$\neg A$</th>
<th>$\neg B$</th>
<th>$A \land B$</th>
<th>$\neg (A \land B)$</th>
<th>$\neg A \lor \neg B$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
<td>F</td>
<td>F</td>
<td>T</td>
<td>F</td>
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</table>

(j): Exercise.

(k):

<table>
<thead>
<tr>
<th>$A$</th>
<th>$\neg A$</th>
<th>$\neg \neg A$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>F</td>
<td>T</td>
</tr>
<tr>
<td>F</td>
<td>T</td>
<td>F</td>
</tr>
</tbody>
</table>

(l):

<table>
<thead>
<tr>
<th>$A$</th>
<th>$B$</th>
<th>$\neg A$</th>
<th>$\neg B$</th>
<th>$A \Rightarrow B$</th>
<th>$\neg B \Rightarrow \neg A$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
<td>F</td>
<td>F</td>
<td>T</td>
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</table>

Having checked all the rules completes the proof of the theorem.

The importance of the rules provided by Th. 1.11 lies in their providing proof techniques, i.e. methods for establishing the truth of statements from statements known or assumed to be true. The rules of Th. 1.11 will be used frequently in proofs throughout this class.

In subsequent proofs, we will also frequently use so-called transitivity of implication as well as transitivity of equivalence (we will encounter equivalence again in the context of relations in Sec. 1.3 below). In preparation for the transitivity rules, we generalize implication to propositional formulas:

**Definition 1.12.** In generalization of the implication operator defined in (1.2), we say the propositional formula $\phi(A_1, \ldots, A_n)$ implies the propositional formula $\psi(A_1, \ldots, A_n)$ (denoted $\phi(A_1, \ldots, A_n) \Rightarrow \psi(A_1, \ldots, A_n)$) if, and only if, each assignment of truth values to the $A_1, \ldots, A_n$ that makes $\phi(A_1, \ldots, A_n)$ true, makes $\psi(A_1, \ldots, A_n)$ true as well.

**Theorem 1.13.** (a) Transitivity of Implication: $(A \Rightarrow B) \land (B \Rightarrow C) \Rightarrow (A \Rightarrow C)$.

(b) Transitivity of Equivalence: $(A \Leftrightarrow B) \land (B \Leftrightarrow C) \Rightarrow (A \Leftrightarrow C)$.

**Proof.** According to Def. 1.12, the rules can be verified by providing truth tables that show that, for all possible assignments of truth values to the propositional formulas on the left-hand side of the implications, either the left-hand side is false or both sides are
true. (a):

<table>
<thead>
<tr>
<th>$A$</th>
<th>$B$</th>
<th>$C$</th>
<th>$A \Rightarrow B$</th>
<th>$B \Rightarrow C$</th>
<th>$(A \Rightarrow B) \land (B \Rightarrow C)$</th>
<th>$A \Rightarrow C$</th>
</tr>
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<tr>
<td>T</td>
<td>T</td>
<td>T</td>
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Having checked both rules, the proof is complete.

Definition and Remark 1.14. A proof of the statement $B$ is a finite sequence of statements $A_1, A_2, \ldots, A_n$ such that $A_1$ is true; for $1 \leq i < n$, $A_i$ implies $A_{i+1}$, and $A_n$ implies $B$. If there exists a proof for $B$, then Th. 1.13(a) guarantees that $B$ is true.

1.3 Set Theory

In the previous section, we have had a first glance at statements and corresponding truth values. In the present section, we will move our focus to the objects such statements are about. Reviewing Example 1.1(a), and recalling that this is a mathematics class rather than one in zoology, the first two statements of Example 1.1(a) are less relevant for us than statements 3–6. As in these examples, we will nearly always be interested in statements involving numbers or collections of numbers or collections of such collections etc.

In modern mathematics, the term one usually uses instead of “collection” is “set”. In 1895, Georg Cantor defined a set as “any collection into a whole $M$ of definite and separate objects $m$ of our intuition or our thought”. The objects $m$ are called the elements of the set $M$. As explained in [Phi16, Sec. A], without restrictions and refinements, Cantor’s set theory is not free of contradictions and, thus, not viable to be used in the foundation of mathematics. Axiomatic set theory provides these necessary restrictions and refinements and an introductory treatment can also be found in [Phi16,
Sec. A]. However, it is possible to follow and understand the rest of this class, without having studied axiomatic set theory.

**Notation 1.15.** We write \( m \in M \) for the statement “\( m \) is an element of the set \( M \).”

**Definition 1.16.** The sets \( M \) and \( N \) are equal, denoted \( M = N \), if, and only if, \( M \) and \( N \) have precisely the same elements.

Definition 1.16 means we know everything about a set \( M \) if, and only if, we know all its elements.

**Definition 1.17.** The set with no elements is called the *empty set*; it is denoted by the symbol \( \emptyset \).

**Example 1.18.** For finite sets, we can simply write down all its elements, for example, \( A := \{0\}, B := \{0, 17.5\}, C := \{5, 1, 5, 3\}, D := \{3, 5, 1\}, E := \{2, \sqrt{2}, -2\} \), where the symbolism “:=” is to be read as “is defined to be equal to”.

Note \( C = D \), since both sets contain precisely the same elements. In particular, the order in which the elements are written down plays no role and a set does not change if an element is written down more than once.

If a set has many elements, instead of writing down all its elements, one might use abbreviations such as \( F := \{-4, -2, \ldots, 20, 22, 24\} \), where one has to make sure the meaning of the dots is clear from the context.

**Definition 1.19.** The set \( A \) is called a *subset* of the set \( B \) (denoted \( A \subseteq B \) and also referred to as the *inclusion* of \( A \) in \( B \)) if, and only if, every element of \( A \) is also an element of \( B \) (one sometimes also calls \( B \) a *superset* of \( A \) and writes \( B \supseteq A \)). Please note that \( A = B \) is allowed in the above definition of a subset. If \( A \subseteq B \) and \( A \neq B \), then \( A \) is called a *strict* subset of \( B \), denoted \( A \subsetneq B \).

If \( B \) is a set and \( P(x) \) is a statement about an element \( x \) of \( B \) (i.e., for each \( x \in B \), \( P(x) \) is either true or false), then we can define a subset \( A \) of \( B \) by writing

\[
A := \{x \in B : P(x)\}.
\]

This notation is supposed to mean that the set \( A \) consists precisely of those elements of \( B \) such that \( P(x) \) is true (has the truth value \( T \) in the language of Sec. 1.2).

**Example 1.20.** (a) For each set \( A \), one has \( A \subseteq A \) and \( \emptyset \subseteq A \).

(b) If \( A \subseteq B \), then \( A = \{x \in B : x \in A\} \).

(c) We have \( \{3\} \subseteq \{6.7, 3, 0\} \). Letting \( A := \{-10, -8, \ldots, 8, 10\} \), we have \( \{-2, 0, 2\} = \{x \in A : x^3 \in A\}, \emptyset = \{x \in A : x + 21 \in A\} \).

**Remark 1.21.** As a consequence of Def. 1.16, the sets \( A \) and \( B \) are equal if, and only if, one has both inclusions, namely \( A \subseteq B \) and \( B \subseteq A \). Thus, when proving the equality of sets, one often divides the proof into two parts, first proving one inclusion, then the other.
Definition 1.22. (a) The intersection of the sets $A$ and $B$, denoted $A \cap B$, consists of all elements that are in $A$ and in $B$. The sets $A, B$ are said to be disjoint if, and only if, $A \cap B = \emptyset$.

(b) The union of the sets $A$ and $B$, denoted $A \cup B$, consists of all elements that are in $A$ or in $B$ (as in the logical disjunction in (1.2), the or is meant nonexclusively). If $A$ and $B$ are disjoint, one sometimes writes $A \cup B$ and speaks of the disjoint union of $A$ and $B$.

(c) The difference of the sets $A$ and $B$, denoted $A \setminus B$ (read “$A$ minus $B$” or “$A$ without $B$”), consists of all elements of $A$ that are not elements of $B$, i.e. $A \setminus B := \{ x \in A : x \notin B \}$. If $B$ is a subset of a given set $A$ (sometimes called the universe in this context), then $A \setminus B$ is also called the complement of $B$ with respect to $A$. In that case, one also writes $B^c := A \setminus B$ (note that this notation suppresses the dependence on $A$).

Example 1.23. (a) Examples of Intersections:

\[
\begin{align*}
\{1, 2, 3\} \cap \{3, 4, 5\} &= \{3\}, \\
\{\sqrt{2}\} \cap \{1, 2, \ldots, 10\} &= \emptyset,
\end{align*}
\]

\[
\{−1, 2, −3, 4, 5\} \cap \{−10, −9, \ldots, −1\} \cap \{−1, 7, −3\} = \{−1, −3\}.
\]

(b) Examples of Unions:

\[
\begin{align*}
\{1, 2, 3\} \cup \{3, 4, 5\} &= \{1, 2, 3, 4, 5\}, \\
\{1, 2, 3\} \cup \{4, 5\} &= \{1, 2, 3, 4, 5\}, \\
\{−1, 2, −3, 4, 5\} \cup \{−99, −98, \ldots, −1\} \cup \{−1, 7, −3\} &= \{−99, −98, \ldots, −2, −1, 2, 4, 5, 7\}.
\end{align*}
\]

(c) Examples of Differences:

\[
\begin{align*}
\{1, 2, 3\} \setminus \{3, 4, 5\} &= \{1, 2\}, \\
\{1, 2, 3\} \setminus \{3, 2, 1, \sqrt{5}\} &= \emptyset, \\
\{−10, −9, \ldots, 9, 10\} \setminus \{0\} &= \{−10, −9, \ldots, −1\} \cup \{1, 2, \ldots, 9, 10\}.
\end{align*}
\]

With respect to the universe $\{1, 2, 3, 4, 5\}$, it is

\[
\{1, 2, 3\}^c = \{4, 5\};
\]

with respect to the universe $\{0, 1, \ldots, 20\}$, it is

\[
\{1, 2, 3\}^c = \{0\} \cup \{4, 5, \ldots, 20\}.
\]

As mentioned earlier, it will often be unavoidable to consider sets of sets. Here are first examples: $\{\emptyset, \{0\}, \{0, 1\}\}, \{\{0, 1\}, \{1, 2\}\}$. 
**Definition 1.24.** Given a set $A$, the set of all subsets of $A$ is called the **power set** of $A$, denoted $\mathcal{P}(A)$ (for reasons explained later (cf. Prop. 2.17), the power set is sometimes also denoted as $2^A$).

**Example 1.25.** Examples of Power Sets:

\[
\mathcal{P}(\emptyset) = \{\emptyset\}, \quad (1.9a)
\]
\[
\mathcal{P}\left(\{0\}\right) = \{\emptyset, \{0\}\}, \quad (1.9b)
\]
\[
\mathcal{P}\left(\mathcal{P}(\{0\})\right) = \mathcal{P}\left(\{\emptyset, \{0\}\}\right) = \{\emptyset, \emptyset, \{\emptyset\}, \mathcal{P}(\{0\})\}. \quad (1.9c)
\]

So far, we have restricted our set-theoretic examples to finite sets. However, not surprisingly, many sets of interest to us will be infinite (we will have to postpone a mathematically precise definition of finite and infinite to Sec. 3.2). We will now introduce the most simple infinite set.

**Definition 1.26.** The set $\mathbb{N} := \{1, 2, 3, \ldots\}$ is called the set of **natural numbers**. Moreover, we define $\mathbb{N}_0 := \{0\} \cup \mathbb{N}$.

**Remark 1.27.** Mathematicians tend to desire as few fundamental objects as possible. One of the consequences is the idea to actually define numbers as special sets: $0 := \emptyset$, $1 := \{0\}$, $2 := \{0, 1\}$; in general, define the natural number $n := \{0, 1, \ldots, n-1\} = (n-1) \cup \{n-1\}$.

The following theorem compiles important set-theoretic rules:

**Theorem 1.28.** Let $A, B, C, U$ be sets.

(a) **Commutativity of Intersections:** $A \cap B = B \cap A$.

(b) **Commutativity of Unions:** $A \cup B = B \cup A$.

(c) **Associativity of Intersections:** $(A \cap B) \cap C = A \cap (B \cap C)$.

(d) **Associativity of Unions:** $(A \cup B) \cup C = A \cup (B \cup C)$.

(e) **Distributivity I:** $A \cap (B \cup C) = (A \cap B) \cup (A \cap C)$.

(f) **Distributivity II:** $A \cup (B \cap C) = (A \cup B) \cap (A \cup C)$.

(g) **De Morgan’s Law I:** $U \setminus (A \cap B) = (U \setminus A) \cup (U \setminus B)$.

(h) **De Morgan’s Law II:** $U \setminus (A \cup B) = (U \setminus A) \cap (U \setminus B)$.

(i) **Double Complement:** If $A \subseteq U$, then $U \setminus (U \setminus A) = A$. 
Proof. In each case, the proof results from the corresponding rule of Th. 1.11:
(a): 
\[ x \in A \cap B \iff x \in A \wedge x \in B \]  
Th. 1.11(c) 
\[ \iff x \in B \wedge x \in A \iff x \in B \cap A. \]

(g): Under the general assumption of \( x \in U \), we have the following equivalences:
\[ x \in U \setminus (A \cap B) \iff \neg(x \in A \cap B) \iff \neg(x \in A \wedge x \in B) \]  
Th. 1.11(i) 
\[ \iff \neg(x \in A) \lor \neg(x \in B) \iff x \in (U \setminus A) \cup (U \setminus B). \]

The proofs of the remaining rules are left as an exercise. 

Remark 1.29. The correspondence between Th. 1.11 and Th. 1.28 is no coincidence. One can actually prove that, starting with an equivalence of propositional formulas \( \phi(A_1, \ldots, A_n) \iff \psi(A_1, \ldots, A_n) \), where both formulas contain only the operators \( \wedge, \lor, \neg \), one obtains a set-theoretic rule (stating an equality of sets) by reinterpretting all statement variables \( A_1, \ldots, A_n \) as variables for sets, all subsets of a universe \( U \), and replacing \( \wedge \) by \( \cap \), \( \lor \) by \( \cup \), and \( \neg \) by \( U \setminus \) (if there are no multiple negations, then we do not need the hypothesis that \( A_1, \ldots, A_n \) are subsets of \( U \)). The procedure also works in the opposite direction – one can start with a set-theoretic formula for an equality of sets and translate it into two equivalent propositional formulas.

1.4 Predicate Calculus

Now that we have introduced sets in the previous section, we have to return to the subject of mathematical logic once more. As it turns out, propositional calculus, which we discussed in Sec. 1.2, does not quite suffice to develop the theory of calculus (nor most other mathematical theories). The reason is that we need to consider statements such as

\[
\begin{align*}
x + 1 &> 0 \text{ holds for each natural number } x. \quad (T) \\
\text{All real numbers are positive.} \quad (F) \\
\text{There exists a natural number bigger than 10.} \quad (T) \\
\text{There exists a real number } x \text{ such that } x^2 = -1. \quad (F) \\
\text{For all natural numbers } n, \text{ there exists a natural number bigger than } n. \quad (T)
\end{align*}
\]

That means we are interested in statements involving \textit{universal quantification} via the quantifier “for all” (one also often uses “for each” or “for every” instead), \textit{existential quantification} via the quantifier “there exists”, or both. The quantifier of universal quantification is denoted by \( \forall \) and the quantifier of existential quantification is denoted by \( \exists \). Using these symbols as well as \( \mathbb{N} \) and \( \mathbb{R} \) to denote the sets of natural and real
numbers, respectively, we can restate (1.10) as

\[ \forall x \in \mathbb{N} \quad x + 1 > 0. \quad \text{(T)} \quad (1.11a) \]
\[ \forall x \in \mathbb{R} \quad x > 0. \quad \text{(F)} \quad (1.11b) \]
\[ \exists n \in \mathbb{N} \quad n > 10. \quad \text{(T)} \quad (1.11c) \]
\[ \exists x \in \mathbb{R} \quad x^2 = -1. \quad \text{(F)} \quad (1.11d) \]
\[ \forall n \in \mathbb{N} \quad \exists m \in \mathbb{N} \quad m > n. \quad \text{(T)} \quad (1.11e) \]

**Definition 1.30.** A **universal statement** has the form

\[ \forall x \in A \quad P(x), \quad (1.12a) \]

whereas an **existential statement** has the form

\[ \exists x \in A \quad P(x), \quad (1.12b) \]

In (1.12), \( A \) denotes a set and \( P(x) \) is a sentence involving the variable \( x \), a so-called **predicate** of \( x \), that becomes a statement (i.e. becomes either true or false) if \( x \) is substituted with any concrete element of the set \( A \) (in particular, \( P(x) \) is allowed to contain further quantifiers, but it must not contain any other quantifier involving \( x \) – one says \( x \) must be a **free** variable in \( P(x) \), not bound by any quantifier in \( P(x) \)).

The universal statement (1.12a) has the truth value \( \text{T} \) if, and only if, \( P(x) \) has the truth value \( \text{T} \) for all elements \( x \in A \); the existential statement (1.12b) has the truth value \( \text{T} \) if, and only if, \( P(x) \) has the truth value \( \text{T} \) for **at least one** element \( x \in A \).

**Remark 1.31.** Some people prefer to write \( \bigwedge_{x \in A} \) instead of \( \forall x \in A \) and \( \bigvee_{x \in A} \) instead of \( \exists x \in A \).

Even though this notation has the advantage of emphasizing that the universal statement can be interpreted as a big logical conjunction and the existential statement can be interpreted as a big logical disjunction, it is significantly less common. So we will stick to \( \forall \) and \( \exists \) in this class.

**Remark 1.32.** According to Def. 1.30, the existential statement (1.12b) is true if, and only if, \( P(x) \) is true for at least one \( x \in A \). So if there is precisely one such \( x \), then (1.12b) is true; and if there are several different \( x \in A \) such that \( P(x) \) is true, then (1.12b) is still true. Uniqueness statements are often of particular importance, and one sometimes writes

\[ \exists! x \in A \quad P(x) \quad (1.13) \]

for the statement “there exists a unique \( x \in A \) such that \( P(x) \) is true”. This notation can be defined as an abbreviation for

\[ \exists x \in A \left( P(x) \land \forall y \in A \left( P(y) \Rightarrow x = y \right) \right) \quad (1.14) \]
1 Foundations: Mathematical Logic and Set Theory

Example 1.33. Here are some examples of uniqueness statements:

\[ \exists! \ n \in \mathbb{N} \ n > 10. \quad (F) \quad (1.15a) \]
\[ \exists! \ 12 > n > 10. \quad (T) \quad (1.15b) \]
\[ \exists! \ 11 > n > 10. \quad (F) \quad (1.15c) \]
\[ \exists! \ x^2 = -1. \quad (F) \quad (1.15d) \]
\[ \exists! \ x^2 = 1. \quad (F) \quad (1.15e) \]
\[ \exists! \ x^2 = 0. \quad (T) \quad (1.15f) \]

Remark 1.34. As for propositional calculus, we also have some important rules for predicate calculus:

(a) Consider the negation of a universal statement, \( \neg \forall_{x \in A} P(x) \), which is true if, and only if, \( P(x) \) does not hold for each \( x \in A \), i.e. if, and only if, there exists at least one \( x \in A \) such that \( \neg P(x) \) is true. We have just proved the rule

\[ \neg \forall_{x \in A} P(x) \iff \exists_{x \in A} \neg P(x). \quad (1.16a) \]

Similarly, consider the negation of an existential statement. We claim the corresponding rule is

\[ \neg \exists_{x \in A} P(x) \iff \forall_{x \in A} \neg P(x). \quad (1.16b) \]

Indeed, we can prove (1.16b) from (1.16a):

\[ \neg \exists_{x \in A} P(x) \quad \text{Th. 1.11(k)} \quad \iff \quad \neg \exists_{x \in A} \neg P(x) \quad \text{(1.16a)} \quad \iff \quad \neg \forall_{x \in A} \neg P(x) \quad \text{Th. 1.11(k)} \quad \iff \quad \forall_{x \in A} \neg P(x). \quad (1.17) \]

One can interpret (1.16) as a generalization of the De Morgan’s laws Th. 1.11(i),(j).

One can actually generalize (1.16) even a bit more: If a statement starts with several quantifiers, then one negates the statement by replacing each \( \forall \) with \( \exists \) and vice versa plus negating the predicate after the quantifiers (see the example in (1.20e) below).

(b) If \( A, B \) are sets and \( P(x, y) \) denotes a predicate of both \( x \) and \( y \), then \( \forall_{x \in A} \forall_{y \in B} P(x, y) \) and \( \forall_{x \in A} \forall_{y \in B} P(x, y) \) both hold true if, and only if, \( P(x, y) \) holds true for each \( x \in A \) and each \( y \in B \), i.e. the order of two consecutive universal quantifiers does not matter:

\[ \forall_{x \in A} \forall_{y \in B} P(x, y) \iff \forall_{y \in B} \forall_{x \in A} P(x, y) \quad (1.18a) \]

In the same way, we obtain the following rule:

\[ \exists_{x \in A} \exists_{y \in B} P(x, y) \iff \exists_{y \in B} \exists_{x \in A} P(x, y). \quad (1.18b) \]
If $A = B$, one also uses abbreviations of the form
\[ \forall_{x,y \in A} P(x, y) \quad \text{for} \quad \forall x \in A \forall y \in A P(x, y), \quad (1.19a) \]
\[ \exists_{x,y \in A} P(x, y) \quad \text{for} \quad \exists x \in A \exists y \in A P(x, y). \quad (1.19b) \]

Generalizing rules (1.18), we can always commute identical quantifiers. Caveat: Quantifiers that are not identical must not be commuted (see Ex. 1.35(d) below).

**Example 1.35. (a)** Negation of universal and existential statements:

Negation of (1.11a) : \( \exists x \in \mathbb{N} \neg (x + 1 > 0) \). (F) \( (1.20a) \)

Negation of (1.11b) : \( \exists x \in \mathbb{R} \neg (x \leq 0) \). (T) \( (1.20b) \)

Negation of (1.11c) : \( \forall n \in \mathbb{N} n \leq 10 \). (F) \( (1.20c) \)

Negation of (1.11d) : \( \forall x \in \mathbb{R} \neg (x^2 = -1) \). (T) \( (1.20d) \)

Negation of (1.11e) : \( \exists n \in \mathbb{N} \forall m \in \mathbb{N} \neg (m > n) \). (F) \( (1.20e) \)

(b) As a more complicated example, consider the negation of the uniqueness statement (1.13), i.e. of (1.14):

\[ \neg \exists x \in A \forall \neg P(x) \quad \leftrightarrow \quad \neg \exists x \in A \left( P(x) \land \forall y \in A (P(y) \Rightarrow x = y) \right) \]

\( (1.16b), \text{Th. } 1.11(a) \)

\[ \forall x \in A \neg P(x) \land \forall y \in A (\neg P(y) \lor x = y) \]

\( \text{Th. } 1.11(j),(k) \)

\[ \neg P(x) \lor \forall y \in A (\neg P(y) \lor x = y) \]

\( \text{Th. } 1.11(j),(k) \)

\[ \forall x \in A \left( \neg P(x) \lor \exists y \in A (P(y) \land x \neq y) \right) \]

\( \text{Th. } 1.11(a) \)

\[ \forall x \in A \left( P(x) \Rightarrow \exists y \in A (P(y) \land x \neq y) \right). \quad (1.21) \]

So how to decode the expression, we have obtained at the end? It states that if $P(x)$ holds for some $x \in A$, then there must be at least a second, different, element $y \in A$ such that $P(y)$ is true. This is, indeed, precisely the negation of $\exists! x \in A P(x)$.

(c) Identical quantifiers commute:

\[ \forall x \in \mathbb{R} \forall n \in \mathbb{N} x^{2n} \geq 0 \quad \leftrightarrow \quad \forall n \in \mathbb{N} x \in \mathbb{R} x^{2n} \geq 0, \quad (1.22a) \]

\[ \forall x \in \mathbb{R} \exists y \in \mathbb{R} ny > x^2 \quad \leftrightarrow \quad \forall x \in \mathbb{R} \exists n \in \mathbb{N} ny > x^2. \quad (1.22b) \]
The following example shows that different quantifiers do, in general, not commute (i.e. do not yield equivalent statements when commuted): While the statement
\[ \forall x \in \mathbb{R} \exists y \in \mathbb{R} \ y > x \quad (1.23a) \]
is true (for each real number \( x \), there is a bigger real number \( y \), e.g. \( y := x + 1 \) will do the job), the statement
\[ \exists y \in \mathbb{R} \forall x \in \mathbb{R} \ y > x \quad (1.23b) \]
is false (for example, since \( y > y \) is false). In particular, (1.23a) and (1.23b) are not equivalent.

Even though (1.13) provides useful notation, it is better not to think of \( \exists! \) as a quantifier. It is really just an abbreviation for (1.14), and it behaves very differently from \( \exists \) and \( \forall \): The following examples show that, in general, \( \exists! \) commutes neither with \( \exists \), nor with itself:
\[ \exists \! \! \exists_{n \in \mathbb{N} \ m \in \mathbb{N}} m < n \not\Rightarrow \exists! \exists_{n \in \mathbb{N} \ m \in \mathbb{N}} m < n \]
(the statement on the left is true, as one can choose \( n = 2 \), but the statement on the right is false, as \( \exists_{n \in \mathbb{N}} m < n \) holds for every \( m \in \mathbb{N} \)). Similarly,
\[ \exists! \exists_{n \in \mathbb{N} \ m \in \mathbb{N}} m < n \not\Rightarrow \exists! \exists!_{n \in \mathbb{N} \ m \in \mathbb{N}} m < n \]
(the statement on the left is still true and the statement on the right is still false (there is no \( m \in \mathbb{N} \) such that \( \exists!_{n \in \mathbb{N}} m < n \)).

**Remark 1.36.** One can make the following observations regarding the strategy for proving universal and existential statements:

(a) To prove that \( \forall x \in A \ P(x) \) is true, one must check the truth of \( P(x) \) for every element \( x \in A \) – examples are not enough!

(b) To prove that \( \forall x \in A \ P(x) \) is false, it suffices to find one \( x \in A \) such that \( P(x) \) is false – such an \( x \) is then called a counterexample and one counterexample is always enough to prove \( \forall x \in A \ P(x) \) is false!

(c) To prove that \( \exists x \in A \ P(x) \) is true, it suffices to find one \( x \in A \) such that \( P(x) \) is true – such an \( x \) is then called an example and one example is always enough to prove \( \exists x \in A \ P(x) \) is true!
The subfield of mathematical logic dealing with quantified statements is called **predicate calculus**. In general, one does not restrict the quantified variables to range only over elements of sets (as we have done above). Again, we refer to [EFT07] for a deeper treatment of the subject.

As an application of quantified statements, let us generalize the notion of union and intersection:

**Definition 1.37.** Let $I \neq \emptyset$ be a nonempty set, usually called an index set in the present context. For each $i \in I$, let $A_i$ denote a set (some or all of the $A_i$ can be identical).

(a) The intersection

$$\bigcap_{i \in I} A_i := \left\{ x : \forall i \in I \ x \in A_i \right\} \tag{1.24a}$$

consists of all elements $x$ that belong to every $A_i$.

(b) The union

$$\bigcup_{i \in I} A_i := \left\{ x : \exists i \in I \ x \in A_i \right\} \tag{1.24b}$$

consists of all elements $x$ that belong to at least one $A_i$. The union is called disjoint if, and only if, for each $i, j \in I$, $i \neq j$ implies $A_i \cap A_j = \emptyset$.

**Proposition 1.38.** Let $I \neq \emptyset$ be an index set, let $M$ denote a set, and, for each $i \in I$, let $A_i$ denote a set. The following set-theoretic rules hold:

(a) $\left( \bigcap_{i \in I} A_i \right) \cap M = \bigcap_{i \in I} (A_i \cap M)$.

(b) $\left( \bigcup_{i \in I} A_i \right) \cup M = \bigcup_{i \in I} (A_i \cup M)$.

(c) $\left( \bigcap_{i \in I} A_i \right) \cup M = \bigcap_{i \in I} (A_i \cup M)$.

(d) $\left( \bigcup_{i \in I} A_i \right) \cap M = \bigcup_{i \in I} (A_i \cap M)$.

(e) $M \setminus \bigcap_{i \in I} A_i = \bigcup_{i \in I} (M \setminus A_i)$.

(f) $M \setminus \bigcup_{i \in I} A_i = \bigcap_{i \in I} (M \setminus A_i)$.

**Proof.** We prove (c) and (e) and leave the remaining proofs as an exercise.
(c):
\[ x \in \left( \bigcap_{i \in I} A_i \right) \cup M \iff x \in M \lor \forall_{i \in I} x \in A_i \iff \forall_{i \in I} (x \in A_i \lor x \in M) \]
\[ \iff x \in \bigcap_{i \in I} (A_i \cup M). \]

To justify the equivalence at \((\ast)\), we make use of Th. 1.11(b) and verify \(\Rightarrow\) and \(\Leftarrow\). For \(\Rightarrow\) note that the truth of \(x \in M\) implies \(x \in A_i \lor x \in M\) is true for each \(i \in I\). If \(x \in A_i\) is true for each \(i \in I\), then \(x \in A_i \lor x \in M\) is still true for each \(i \in I\). To verify \(\Leftarrow\), note that the existence of \(i \in I\) such that \(x \in M\) implies the truth of \(x \in M \lor \forall_{i \in I} x \in A_i\).

If \(x \in M\) is false for each \(i \in I\), then \(x \in A_i\) must be true for each \(i \in I\), showing \(x \in M \lor \forall_{i \in I} x \in A_i\) is true also in this case.

(e):
\[ x \in M \setminus \bigcap_{i \in I} A_i \iff x \in M \land \forall_{i \in I} x \in A_i \iff x \in M \land \exists_{i \in I} x \notin A_i \]
\[ \iff \exists_{i \in I} x \in M \setminus A_i \iff x \in \bigcup_{i \in I} (M \setminus A_i), \]

completing the proof.

**Example 1.39.** We have the following identities of sets:
\[ \bigcap_{x \in \mathbb{R}} \mathbb{N} = \mathbb{N}, \quad (1.25a) \]
\[ \bigcap_{n \in \mathbb{N}} \{1, 2, \ldots, n\} = \{1\}, \quad (1.25b) \]
\[ \bigcup_{x \in \mathbb{R}} \mathbb{N} = \mathbb{N}, \quad (1.25c) \]
\[ \bigcup_{n \in \mathbb{N}} \{1, 2, \ldots, n\} = \mathbb{N}, \quad (1.25d) \]
\[ \mathbb{N} \setminus \bigcup_{n \in \mathbb{N}} \{2n\} = \{1, 3, 5, \ldots\} = \bigcap_{n \in \mathbb{N}} (\mathbb{N} \setminus \{2n\}). \quad (1.25e) \]

Comparing with the notation of Def. 1.37, in (1.25a), for example, we have \(I = \mathbb{R}\) and \(A_i = \mathbb{N}\) for each \(i \in I\) (where, in (1.25a), we have written \(x\) instead of \(i\)). Similarly, in (1.25b), we have \(I = \mathbb{N}\) and \(A_n = \{1, 2, \ldots, n\}\) for each \(n \in I\).

2 Functions and Relations

2.1 Functions

**Definition 2.1.** Let \(A, B\) be sets. Given \(x \in A, y \in B\), the set
\[ (x, y) := \big\{ \{x\}, \{x, y\} \big\} \quad (2.1) \]
is called the \textit{ordered pair} (often shortened to just \textit{pair}) consisting of \(x\) and \(y\). The set of all such pairs is called the Cartesian product \(A \times B\), i.e.
\[
A \times B := \{(x, y) : x \in A \land y \in B\}. \tag{2.2}
\]

\textbf{Example 2.2.} Let \(A\) be a set.
\[
A \times \emptyset = \emptyset \times A = \emptyset, \tag{2.3a}
\]
\[
\{1, 2\} \times \{1, 2, 3\} = \{(1, 1), (1, 2), (1, 3), (2, 1), (2, 2), (2, 3)\} \tag{2.3b}
\]
\[
\neq \{1, 2, 3\} \times \{1, 2\} = \{(1, 1), (1, 2), (2, 1), (2, 2), (3, 1), (3, 2)\}. \tag{2.3c}
\]

Also note that, for \(x \neq y\),
\[
(x, y) = \{\{x\}, \{x, y\}\} \neq \{\{y\}, \{x, y\}\} = (y, x). \tag{2.4}
\]

\textbf{Definition 2.3.} Given sets \(A, B\), a \textit{function} or \textit{map} \(f\) is an assignment rule that assigns to each \(x \in A\) a unique \(y \in B\). One then also writes \(f(x)\) for the element \(y\). The set \(A\) is called the \textit{domain} of \(f\), denoted \(\mathcal{D}(f)\), and \(B\) is called the \textit{range} of \(f\), denoted \(\mathcal{R}(f)\). The information about a map \(f\) can be concisely summarized by the notation
\[
f : A \rightarrow B, \quad x \mapsto f(x), \tag{2.5}
\]
where \(x \mapsto f(x)\) is called the \textit{assignment rule} for \(f\), \(f(x)\) is called the \textit{image} of \(x\), and \(x\) is called a \textit{preimage} of \(f(x)\) (the image must be unique, but there might be several preimages). The set
\[
\text{graph}(f) := \{(x, y) \in A \times B : y = f(x)\} \tag{2.6}
\]
is called the \textit{graph} of \(f\) (not to be confused with pictures visualizing the function \(f\), which are also called graph of \(f\)). If one wants to be completely precise, then one identifies the function \(f\) with the ordered triple \((A, B, \text{graph}(f))\).

The set of all functions with domain \(A\) and range \(B\) is denoted by \(\mathcal{F}(A, B)\) or \(B^A\), i.e.
\[
\mathcal{F}(A, B) := B^A := \{ (f : A \rightarrow B) : A = \mathcal{D}(f) \land B = \mathcal{R}(f) \}. \tag{2.7}
\]

\textbf{Caveat:} Some authors reserve the word \textit{map} for continuous functions, but we use function and map synonymously.

\textbf{Definition 2.4.} Let \(A, B\) be sets and \(f : A \rightarrow B\) a function.

\textbf{(a)} If \(T\) is a subset of \(A\), then
\[
f(T) := \{ f(x) \in B : x \in T \}\tag{2.8}
\]
is called the \textit{image} of \(T\) under \(f\).

\textbf{(b)} If \(U\) is a subset of \(B\), then
\[
f^{-1}(U) := \{ x \in A : f(x) \in U \}\tag{2.9}
\]
is called the \textit{preimage} or \textit{inverse image} of \(U\) under \(f\).
(c) \( f \) is called \textit{injective} or \textit{one-to-one} if, and only if, every \( y \in B \) has at most one preimage, i.e. if, and only if, the preimage of \( \{y\} \) has at most one element:

\[
\text{\( f \) injective } \iff \forall_{y \in B} \left( f^{-1}\{y\} = \emptyset \lor \exists!_{x \in A} f(x) = y \right)
\]

\[
\iff \forall_{x_1,x_2 \in A} (x_1 \neq x_2 \Rightarrow f(x_1) \neq f(x_2)). \quad (2.10)
\]

(d) \( f \) is called \textit{surjective} or \textit{onto} if, and only if, every element of the range of \( f \) has a preimage:

\[
\text{\( f \) surjective } \iff \forall_{y \in B} \exists x \in A \ f(x) \iff \forall_{y \in B} f^{-1}\{y\} \neq \emptyset. \quad (2.11)
\]

(e) \( f \) is called \textit{bijective} if, and only if, \( f \) is injective and surjective.

\textbf{Example 2.5.} Examples of Functions:

\[
\begin{align*}
\text{\( f \)} &: \{1, 2, 3, 4, 5\} \longrightarrow \{1, 2, 3, 4, 5\}, & f(x) := -x + 6, \quad (2.12a) \\
\text{\( g \)} &: \mathbb{N} \longrightarrow \mathbb{N}, & g(n) := 2n, \quad (2.12b) \\
\text{\( h \)} &: \mathbb{N} \longrightarrow \{2, 4, 6, \ldots \}, & h(n) := 2n, \quad (2.12c) \\
\tilde{h} &: \mathbb{N} \longrightarrow \{2, 4, 6, \ldots \}, & \tilde{h}(n) := \begin{cases} 
\frac{n}{2} & \text{for } n \text{ even}, \\
\frac{n+1}{2} & \text{for } n \text{ odd},
\end{cases} \quad (2.12d) \\
\text{\( G \)} &: \mathbb{N} \longrightarrow \mathbb{R}, & G(n) := n/(n + 1), \quad (2.12e) \\
\text{\( F \)} &: \mathcal{P}(\mathbb{N}) \longrightarrow \mathcal{P}(\mathcal{P}(\mathbb{N})), & F(A) := \mathcal{P}(A). \quad (2.12f)
\end{align*}
\]

Instead of \( f(x) := -x + 6 \) in (2.12a), one can also write \( x \mapsto -x + 6 \) and analogously in the other cases. Also note that, in the strict sense, functions \( g \) and \( h \) are different, since their ranges are different (however, using the following Def. 2.4(a), they have the same \textit{image} in the sense that \( g(\mathbb{N}) = h(\mathbb{N}) \)). Furthermore,

\[
(f(\{1, 2\}) = \{5, 4\} = f^{-1}(\{1, 2\}), \quad \tilde{h}^{-1}(\{2, 4, 6\}) = \{1, 2, 3, 4, 5, 6\}, \quad (2.13)
\]

\( f \) is bijective; \( g \) is injective, but not surjective; \( h \) is bijective; \( \tilde{h} \) is surjective, but not injective. Can you figure out if \( G \) and \( F \) are injective and/or surjective?

\textbf{Example 2.6. (a)} For each nonempty set \( A \), the map \( \text{Id} : A \longrightarrow A, \text{Id}(x) := x \), is called the \textit{identity} on \( A \). If one needs to emphasize that \( \text{Id} \) operates on \( A \), then one also writes \( \text{Id}_A \) instead of \( \text{Id} \). The identity is clearly bijective.

\textbf{(b)} Let \( A, B \) be nonempty sets. A map \( f : A \longrightarrow B \) is called \textit{constant} if, and only if, there exists \( c \in B \) such that \( f(x) = c \) for each \( x \in A \). In that case, one also writes \( f \equiv c \), which can be read as “\( f \) is identically equal to \( c \)”. If \( f \equiv c, \emptyset \neq T \subseteq A \), and \( U \subseteq B \), then

\[
\begin{align*}
\text{\( f(T) = \{c\}, \quad f^{-1}(U) = \begin{cases} 
A & \text{for } c \in U, \\
\emptyset & \text{for } c \notin U.
\end{cases} \quad (2.14)
\end{align*}
\]

\( f \) is injective if, and only if, \( A = \{x\}; f \) is surjective if, and only if, \( B = \{c\} \).
Given $A \subseteq X$, the map
\[ \iota : A \rightarrow X, \quad \iota(x) := x, \]
(2.15)
is called inclusion (also embedding or imbedding). An inclusion is always injective; it is surjective if, and only if $A = X$, i.e. if, and only if, it is the identity on $A$.

Given $A \subseteq X$ and a map $f : X \rightarrow B$, the map $g : A \rightarrow B, g(x) = f(x)$, is called the restriction of $f$ to $A$; $f$ is called the extension of $g$ to $X$. In this situation, one also uses the notation $f|_A$ for $g$ (some authors prefer the notation $f|_A$ or $f|_A$).

There are several important rules regarding functions and set-theoretic operations. However, we will not make use of them in this class, and the interested student can find them in [Phi16, Th. 2.7].

**Definition 2.7.** The composition of maps $f$ and $g$ with $f : A \rightarrow B, g : C \rightarrow D$, and $f(A) \subseteq C$ is defined to be the map
\[ g \circ f : A \rightarrow D, \quad (g \circ f)(x) := g(f(x)). \]
(2.16)
The expression $g \circ f$ is read as “$g$ after $f$” or “$g$ composed with $f$”.

**Example 2.8.** Consider the maps
\[
\begin{align*}
f & : \mathbb{N} \rightarrow \mathbb{R}, \quad n \mapsto n^2, \quad (2.17a) \\
g & : \mathbb{N} \rightarrow \mathbb{R}, \quad n \mapsto 2n. \quad (2.17b)
\end{align*}
\]
We obtain $f(\mathbb{N}) = \{1, 4, 9, \ldots \} \subseteq D(g), g(\mathbb{N}) = \{2, 4, 6, \ldots \} \subseteq D(f)$, and the compositions
\[
\begin{align*}
(g \circ f) & : \mathbb{N} \rightarrow \mathbb{R}, \quad (g \circ f)(n) = g(n^2) = 2n^2, \quad (2.18a) \\
(f \circ g) & : \mathbb{N} \rightarrow \mathbb{R}, \quad (f \circ g)(n) = f(2n) = 4n^2. \quad (2.18b)
\end{align*}
\]
showing that composing functions is, in general, not commutative, even if the involved functions have the same domain and the same range.

**Proposition 2.9.** Consider maps $f : A \rightarrow B, g : C \rightarrow D, h : E \rightarrow F$, satisfying $f(A) \subseteq C$ and $g(C) \subseteq E$.

(a) **Associativity of Compositions:**
\[ h \circ (g \circ f) = (h \circ g) \circ f. \quad (2.19) \]

(b) **One has the following law for forming preimages:**
\[
\forall W \in P(D) \quad (g \circ f)^{-1}(W) = f^{-1}(g^{-1}(W)). \quad (2.20)
\]
Proof. (a): Both \( h \circ (g \circ f) \) and \((h \circ g) \circ f \) map \( A \) into \( F \). So it just remains to prove 
\[
(h \circ (g \circ f))(x) = ((h \circ g) \circ f)(x)
\]
for each \( x \in A \). One computes, for each \( x \in A \),
\[
(h \circ (g \circ f))(x) = h((g \circ f)(x)) = h((g(f(x)))) = (h \circ g)(f(x)) = ((h \circ g) \circ f)(x),
\]
(2.21)
establishing the case.

(b): Exercise. ■

Definition 2.10. A function \( g : B \to A \) is called a right inverse (resp. left inverse) of a function \( f : A \to B \) if, and only if, \( f \circ g = \text{Id}_B \) (resp. \( g \circ f = \text{Id}_A \)). Moreover, \( g \) is called an inverse of \( f \) if, and only if, it is both a right and a left inverse. If \( g \) is an inverse of \( f \), then one also writes \( f^{-1} \) instead of \( g \). The map \( f \) is called (right, left) invertible if, and only if, there exists a (right, left) inverse for \( f \).

Example 2.11. (a) Consider the map
\[
f : \mathbb{N} \to \mathbb{N}, \quad f(n) := 2n. \tag{2.22a}
\]
The maps
\[
g_1 : \mathbb{N} \to \mathbb{N}, \quad g_1(n) := \begin{cases} n/2 & \text{if } n \text{ even}, \\ 1 & \text{if } n \text{ odd}, \end{cases} \tag{2.22b}
\]
\[
g_2 : \mathbb{N} \to \mathbb{N}, \quad g_2(n) := \begin{cases} n/2 & \text{if } n \text{ even}, \\ 2 & \text{if } n \text{ odd}, \end{cases} \tag{2.22c}
\]
both constitute left inverses of \( f \). It follows from Th. 2.12(c) below that \( f \) does not have a right inverse.

(b) Consider the map
\[
f : \mathbb{N} \to \mathbb{N}, \quad f(n) := \begin{cases} n/2 & \text{for } n \text{ even}, \\ (n + 1)/2 & \text{for } n \text{ odd}. \end{cases} \tag{2.23a}
\]
The maps
\[
g_1 : \mathbb{N} \to \mathbb{N}, \quad g_1(n) := 2n, \tag{2.23b}
\]
\[
g_2 : \mathbb{N} \to \mathbb{N}, \quad g_2(n) := 2n - 1, \tag{2.23c}
\]
both constitute right inverses of \( f \). It follows from Th. 2.12(c) below that \( f \) does not have a left inverse.

(c) The map
\[
f : \mathbb{N} \to \mathbb{N}, \quad f(n) := \begin{cases} n - 1 & \text{for } n \text{ even}, \\ n + 1 & \text{for } n \text{ odd}. \end{cases} \tag{2.24a}
\]
is its own inverse, i.e. $f^{-1} = f$. For the map

$$g : \mathbb{N} \rightarrow \mathbb{N}, \quad g(n) := \begin{cases} 2 & \text{for } n = 1, \\ 3 & \text{for } n = 2, \\ 1 & \text{for } n = 3, \\ n & \text{for } n \notin \{1, 2, 3\}, \end{cases} \quad (2.24b)$$

the inverse is

$$g^{-1} : \mathbb{N} \rightarrow \mathbb{N}, \quad g^{-1}(n) := \begin{cases} 3 & \text{for } n = 1, \\ 1 & \text{for } n = 2, \\ 2 & \text{for } n = 3, \\ n & \text{for } n \notin \{1, 2, 3\}. \end{cases} \quad (2.24c)$$

While Examples 2.11(a),(b) show that left and right inverses are usually not unique, they are unique provided $f$ is bijective (see Th. 2.12(c)).

**Theorem 2.12.** Let $A, B$ be nonempty sets.

(a) $f : A \rightarrow B$ is right invertible if, and only if, $f$ is surjective.

(b) $f : A \rightarrow B$ is left invertible if, and only if, $f$ is injective.

(c) $f : A \rightarrow B$ is invertible if, and only if, $f$ is bijective. In this case, the right inverse and the left inverse are unique and both identical to the inverse.

**Proof.** (a): If $f$ is surjective, then, for each $y \in B$, there exists $x_y \in f^{-1}\{y\}$ such that $f(x_y) = y$. Define

$$g : B \rightarrow A, \quad g(y) := x_y \quad (2.25)$$

(note to the interested reader: the definition of $g$ is, in general, not as unproblematic as it might seem – $g$ is a so-called choice function, and its definition makes use of the axiom of choice, see [Phi16, Sec. A.4]). Then, for each $y \in B$, $f(g(y)) = y$, showing $g$ is a right inverse of $f$. Conversely, if $g : B \rightarrow A$ is a right inverse of $f$, then, for each $y \in B$, it is $y = f(g(y))$, showing that $g(y) \in A$ is a preimage of $y$, i.e. $f$ is surjective.

(b): Fix $a \in A$. If $f$ is injective, then, for each $y \in B$ with $f^{-1}\{y\} \neq \emptyset$, let $x_y$ denote the unique element in $A$ satisfying $f(x_y) = y$. Define

$$g : B \rightarrow A, \quad g(y) := \begin{cases} x_y & \text{for } f^{-1}\{y\} \neq \emptyset, \\ a & \text{otherwise}. \end{cases} \quad (2.26)$$

Then, for each $x \in A$, $g(f(x)) = x$, showing $g$ is a left inverse of $f$. Conversely, if $g : B \rightarrow A$ is a left inverse of $f$ and $x_1, x_2 \in A$ with $f(x_1) = f(x_2) = y$, then $x_1 = (g \circ f)(x_1) = g(f(x_1)) = g(f(x_2)) = (g \circ f)(x_2) = x_2$, showing $y$ has precisely one preimage and $f$ is injective.
(c): Assume \( g \) to be a left inverse of \( f \) and \( h \) to be a right inverse of \( f \). Then, for each \( y \in B \),
\[
g(y) = (g \circ (f \circ h))(y) = ((g \circ f) \circ h)(y) = h(y),
\]
showing \( g = h \). In particular, if \( f \) has an inverse \( f^{-1} \), then \( g = h = f^{-1} \). If \( f \) is invertible, then \( f \) is bijective by (a) and (b). If \( f \) is bijective, then, by (a) and (b), \( f \) has a left inverse \( g \) and a right inverse \( h \). By (2.27), \( g = h \), i.e. \( f \) is invertible. \( \blacksquare \)

**Theorem 2.13.** Consider maps \( f : A \rightarrow B \), \( g : B \rightarrow C \). If \( f \) and \( g \) are both injective (resp. both surjective, both bijective), then so is \( g \circ f \). Moreover, in the bijective case, one has
\[
(g \circ f)^{-1} = f^{-1} \circ g^{-1}.
\]

**Proof.** Exercise. \( \blacksquare \)

**Definition 2.14.**

(a) Given an index set \( I \) and a set \( A \), a map \( f : I \rightarrow A \) is sometimes called a family (of elements in \( A \)), and is denoted in the form \( f = (a_i)_{i \in I} \) with \( a_i := f(i) \). When using this representation, one often does not even specify \( f \) and \( A \), especially if the \( a_i \) are themselves sets.

(b) A sequence in a set \( A \) is a family of elements in \( A \), where the index set is the set of natural numbers \( \mathbb{N} \). In this case, one writes \( (a_n)_{n \in \mathbb{N}} \) or \( (a_1, a_2, \ldots) \). More generally, a family is called a sequence, given a bijective map between the index set \( I \) and a subset of \( \mathbb{N} \).

(c) Given a family of sets \((A_i)_{i \in I}\), we define the Cartesian product of the \( A_i \) to be the set of functions
\[
\prod_{i \in I} A_i := \left\{ f : I \rightarrow \bigcup_{j \in I} A_j : \forall i \in I \ f(i) \in A_i \right\}.
\]

If \( I \) has precisely \( n \) elements with \( n \in \mathbb{N} \), then the elements of the Cartesian product \( \prod_{i \in I} A_i \) are called (ordered) \( n \)-tuples, (ordered) triples for \( n = 3 \).

**Example 2.15.**

(a) Using the notion of family, we can now say that the intersection \( \bigcap_{i \in I} A_i \) and union \( \bigcup_{i \in I} A_i \) as defined in Def. 1.37 are the intersection and union of the family of sets \((A_i)_{i \in I}\), respectively. As a concrete example, let us revisit (1.25b), where we have
\[
(A_n)_{n \in \mathbb{N}}, \quad A_n := \{1, 2, \ldots, n\}, \quad \bigcap_{n \in \mathbb{N}} A_n = \{1\}.
\]
(b) Examples of Sequences:

Sequence in \{0, 1\} : \( (1, 0, 1, 0, 1, 0, \ldots) \), \hfill (2.31a)

Sequence in \( \mathbb{N} \) : \( (n^2)_{n \in \mathbb{N}} = (1, 4, 9, 16, 25, \ldots) \), \hfill (2.31b)

Sequence in \( \mathbb{R} \) : \( ((-1)^n \sqrt{n})_{n \in \mathbb{N}} = (-1, \sqrt{2}, -\sqrt{3}, \ldots) \), \hfill (2.31c)

Sequence in \( \mathbb{R} \) : \( (1/n)_{n \in \mathbb{N}} = \left( \frac{1}{2}, \frac{1}{3}, \ldots \right) \), \hfill (2.31d)

Finite Sequence in \( \mathcal{P}(\mathbb{N}) \) : \( \{3, 2, 1\}, \{2, 1\}, \{1\}, \emptyset \). \hfill (2.31e)

(c) The Cartesian product \( \prod_{i \in I} A_i \), where all sets \( A_i = A \), is the same as \( A^I \), the set of all functions from \( I \) into \( A \). So, for example, \( \prod_{n \in \mathbb{N}} \mathbb{R} = \mathbb{R}^\mathbb{N} \) is the set of all sequences in \( \mathbb{R} \). If \( I = \{1, 2, \ldots, n\} \) with \( n \in \mathbb{N} \), then

\[
\prod_{i \in I} A = A^{(1,2,\ldots,n)} = \prod_{i=1}^n A =: A^n
\]

is the set of all \( n \)-tuples with entries from \( A \).

In the following, we explain the common notation \( 2^A \) for the power set \( \mathcal{P}(A) \) of a set \( A \). It is related to a natural identification between subsets and their corresponding characteristic function.

**Definition 2.16.** Let \( A \) be a set and let \( B \subseteq A \) be a subset of \( A \). Then

\[
\chi_B : A \rightarrow \{0, 1\}, \quad \chi_B(x) := \begin{cases} 1 & \text{if } x \in B, \\ 0 & \text{if } x \notin B, \end{cases}
\]

is called the characteristic function of the set \( B \) (with respect to the universe \( A \)). One also finds the notations \( 1_B \) and \( \mathbb{1}_B \) instead of \( \chi_B \) (note that all the notations suppress the dependence of the characteristic function on the universe \( A \)).

**Proposition 2.17.** Let \( A \) be a set. Then the map

\[
\chi : \mathcal{P}(A) \rightarrow \{0, 1\}^A, \quad \chi(B) := \chi_B,
\]

is bijective (recall that \( \mathcal{P}(A) \) denotes the power set of \( A \) and \( \{0, 1\}^A \) denotes the set of all functions from \( A \) into \( \{0, 1\} \)).

**Proof.** \( \chi \) is injective: Let \( B, C \in \mathcal{P}(A) \) with \( B \neq C \). By possibly switching the names of \( B \) and \( C \), we may assume there exists \( x \in B \) such that \( x \notin C \). Then \( \chi_B(x) = 1 \), whereas \( \chi_C(x) = 0 \), showing \( \chi(B) \neq \chi(C) \), proving \( \chi \) is injective.

\( \chi \) is surjective: Let \( f : A \rightarrow \{0, 1\} \) be an arbitrary function and define \( B := \{x \in A : f(x) = 1\} \). Then \( \chi(B) = \chi_B = f \), proving \( \chi \) is surjective. \( \blacksquare \)

Proposition 2.17 allows one to identify the sets \( \mathcal{P}(A) \) and \( \{0, 1\}^A \) via the bijective map \( \chi \). This fact together with the common practise of set theory to identify the number 2 with the set \( \{0, 1\} \) (cf. Rem. 1.27 above) explains the notation \( 2^A \) for \( \mathcal{P}(A) \).
2.2 Relations

Definition 2.18. Given sets $A$ and $B$, a relation is a subset $R$ of $A \times B$ (if one wants to be completely precise, a relation is an ordered triple $(A,B,R)$, where $R \subseteq A \times B$). If $A = B$, then we call $R$ a relation on $A$. One says that $a \in A$ and $b \in B$ are related according to the relation $R$ if, and only if, $(a,b) \in R$. In this context, one usually writes $a R b$ instead of $(a,b) \in R$.

Example 2.19. (a) The relations we are probably most familiar with are $=$ and $\leq$.

The relation $R$ of equality, usually denoted $=$, makes sense on every nonempty set $A$:

$$R := \Delta(A) := \{(x,x) \in A \times A : x \in A\}. \quad (2.35)$$

The set $\Delta(A)$ is called the diagonal of the Cartesian product, i.e., as a subset of $A \times A$, the relation of equality is identical to the diagonal:

$$x = y \iff x R y \iff (x,y) \in R = \Delta(A). \quad (2.36)$$

Similarly, the relation $\leq$ on $\mathbb{R}$ is identical to the set

$$R_{\leq} := \{(x,y) \in \mathbb{R}^2 : x \leq y\}. \quad (2.37)$$

(b) Every function $f : A \rightarrow B$ is a relation, namely the relation

$$R_f = \{(x,y) \in A \times B : y = f(x)\} = \text{graph}(f). \quad (2.38)$$

Conversely, if $B \neq \emptyset$, then every relation $R \subseteq A \times B$ uniquely corresponds to the function

$$f_R : A \rightarrow \mathcal{P}(B), \quad f_R(x) = \{y \in B : x R y\}. \quad (2.39)$$

Definition 2.20. Let $R$ be a relation on the set $A$.

(a) $R$ is called reflexive if, and only if,

$$\forall_{x \in A} x R x, \quad (2.40)$$

i.e. if, and only if, every element is related to itself.

(b) $R$ is called symmetric if, and only if,

$$\forall_{x,y \in A} (x R y \Rightarrow y R x), \quad (2.41)$$

i.e. if, and only if, each $x$ is related to $y$ if, and only if, $y$ is related to $x$.

(c) $R$ is called antisymmetric if, and only if,

$$\forall_{x,y \in A} ((x R y \land y R x) \Rightarrow x = y), \quad (2.42)$$

i.e. if, and only if, the only possibility for $x$ to be related to $y$ at the same time that $y$ is related to $x$ is in the case $x = y$. 
(d) $R$ is called transitive if, and only if,

$$\forall x,y,z \in A \ ( (x \, R \, y \land y \, R \, z) \Rightarrow x \, R \, z ), \tag{2.43}$$

i.e. if, and only if, the relatedness of $x$ and $y$ together with the relatedness of $y$ and $z$ implies the relatedness of $x$ and $z$.

**Example 2.21.** The relations $=$ and $\leq$ on $\mathbb{R}$ (or $\mathbb{N}$) are reflexive, antisymmetric, and transitive; $=$ is also symmetric, whereas $\leq$ is not; $<$ is antisymmetric (since $x < y \land y < x$ is always false) and transitive, but neither reflexive nor symmetric. The relation

$$R := \{(x, y) \in \mathbb{N}^2 : (x, y \text{ are both even}) \lor (x, y \text{ are both odd})\} \tag{2.44}$$

on $\mathbb{N}$ is not antisymmetric, but reflexive, symmetric, and transitive. The relation

$$S := \{(x, y) \in \mathbb{N}^2 : y = x^2\} \tag{2.45}$$

is not transitive (for example, $2 S 4$ and $4 S 16$, but not $2 S 16$), not reflexive, not symmetric; it is only antisymmetric.

**Definition 2.22.** A relation $R$ on a set $A$ is called an **equivalence relation** if, and only if, $R$ is reflexive, symmetric, and transitive. If $R$ is an equivalence relations, then one often writes $x \sim y$ instead of $x \, R \, y$.

**Example 2.23.** (a) The equality relation $=$ is an equivalence relation on each $A \neq \emptyset$.

(b) The relation $R$ defined in (2.44) is an equivalence relation on $\mathbb{N}$.

(c) Given a disjoint union $A = \bigcup_{i \in I} A_i$ with every $A_i \neq \emptyset$ (which is sometimes called a **decomposition** of $A$), an equivalence relation on $A$ is defined by

$$x \sim y \iff \exists i \in I \ (x \in A_i \land y \in A_i). \tag{2.46}$$

Conversely, given an equivalence relation $\sim$ on a nonempty set $A$, we can construct a decomposition $A = \bigcup_{i \in I} A_i$ such that (2.46) holds: For each $x \in A$, define

$$[x] := \{y \in A : x \sim y\}, \tag{2.47}$$

called the **equivalence class** of $x$; each $y \in [x]$ is called a **representative** of $[x]$. One verifies that the properties of $\sim$ guarantee

$$([x] = [y] \Leftrightarrow x \sim y) \quad \land \quad ([x] \cap [y] = \emptyset \Leftrightarrow \neg (x \sim y)). \tag{2.48}$$

The set of all equivalence classes $I := A/\sim := \{[x] : x \in A\}$ is called the **quotient set** of $A$ by $\sim$, and $A = \bigcup_{i \in I} A_i$ with $A_i := i$ for each $i \in I$ is the desired decomposition of $A$.

**Definition 2.24.** A relation $R$ on a set $A$ is called a **partial order** if, and only if, $R$ is reflexive, antisymmetric, and transitive. If $R$ is a partial order, then one usually writes $x \leq y$ instead of $x \, R \, y$. A partial order $\leq$ is called a **total** or **linear order** if, and only if, for each $x, y \in A$, one has $x \leq y$ or $y \leq x$. 
Notation 2.25. Given a (partial or total) order \( \leq \) on \( A \neq \emptyset \), we write \( x < y \) if, and only if, \( x \leq y \) and \( x \neq y \), calling \( < \) the *strict* order corresponding to \( \leq \) (note that the strict order is never a partial order).

Definition 2.26. Let \( \leq \) be a partial order on \( A \neq \emptyset \), \( \emptyset \neq B \subseteq A \).

(a) \( x \in A \) is called *lower* (resp. *upper*) bound for \( B \) if, and only if, \( x \leq b \) (resp. \( b \leq x \)) for each \( b \in B \). Moreover, \( B \) is called *bounded from below* (resp. from above) if, and only if, there exists a lower (resp. upper) bound for \( B \); \( B \) is called *bounded* if, and only if, it is bounded from above and from below.

(b) \( x \in B \) is called *minimum* or just *min* (resp. *maximum* or *max*) of \( B \) if, and only if, \( x \) is a lower (resp. upper) bound for \( B \). One writes \( x = \min B \) if \( x \) is minimum and \( x = \max B \) if \( x \) is maximum.

(c) A maximum of the set of lower bounds of \( B \) (i.e. a largest lower bound) is called *infimum* of \( B \), denoted \( \inf B \); a minimum of the set of upper bounds of \( B \) (i.e. a smallest upper bound) is called *supremum* of \( B \), denoted \( \sup B \).

Example 2.27. (a) For each \( A \subseteq \mathbb{R} \), the usual relation \( \leq \) defines a total order on \( A \).

For \( A = \mathbb{R} \), we see that \( \mathbb{N} \) has 0 and 1 as lower bound with \( 0 = \min \mathbb{N} = \inf \mathbb{N} \). On the other hand, \( \mathbb{N} \) is unbounded from above. The set \( M := \{1, 2, 3\} \) is bounded with \( \min M = 1 \), \( \max M = 3 \). The positive real numbers \( \mathbb{R}^+ := \{x \in \mathbb{R} : x > 0\} \) have \( \inf \mathbb{R}^+ = 0 \), but they do not have a minimum (if \( x > 0 \), then \( 0 < x/2 < x \)).

(b) Consider \( A := \mathbb{N} \times \mathbb{N} \). Then

\[
(m_1, m_2) \leq (n_1, n_2) \iff m_1 \leq n_1 \land m_2 \leq n_2, \tag{2.49}
\]

defines a partial order on \( A \) that is not a total order (for example, neither \((1, 2) \leq (2, 1)\) nor \((2, 1) \leq (1, 2)\)). For the set

\[
B := \{(1, 1), (2, 1), (1, 2)\}, \tag{2.50}
\]

we have \( \inf B = \min B = (1, 1) \), \( B \) does not have a max, but \( \sup B = (2, 2) \) (if \( (m, n) \in A \) is an upper bound for \( B \), then \( (2, 1) \leq (m, n) \) implies \( 2 \leq m \) and \( (1, 2) \leq (m, n) \) implies \( 2 \leq n \), i.e. \( (2, 2) \leq (m, n) \); since \( (2, 2) \) is clearly an upper bound for \( B \), we have proved \( \sup B = (2, 2) \)).

A different order on \( A \) is the so-called *lexicographic order* defined by

\[
(m_1, m_2) \leq (n_1, n_2) \iff m_1 < n_1 \lor (m_1 = n_1 \land m_2 \leq n_2). \tag{2.51}
\]

In contrast to the order from (2.49), the lexicographic order does define a total order on \( A \).

Lemma 2.28. Let \( \leq \) be a partial order on \( A \neq \emptyset \), \( \emptyset \neq B \subseteq A \). Then the relation \( \geq \), defined by

\[
x \geq y \iff y \leq x, \tag{2.52}
\]
is also a partial order on \( A \). Moreover, using obvious notation, we have, for each \( x \in A \),

\[
\begin{align*}
    x \leq & \text{- lower bound for } B \iff x \geq & \text{- upper bound for } B, \quad (2.53a) \\
    x \leq & \text{- upper bound for } B \iff x \geq & \text{- lower bound for } B, \quad (2.53b) \\
x = \min_\leq & B \iff x = \max_\geq B, \quad (2.53c) \\
x = \max_\leq & B \iff x = \min_\geq B, \quad (2.53d) \\
x = \inf_\leq & B \iff x = \sup_\geq B, \quad (2.53e) \\
x = \sup_\leq & B \iff x = \inf_\geq B. \quad (2.53f)
\end{align*}
\]

**Proof.** Reflexivity, antisymmetry, and transitivity of \( \leq \) clearly imply the same properties for \( \geq \), respectively. Moreover

\[
x \leq \text{- lower bound for } B \iff \forall_{b \in B} x \leq b \iff \forall_{b \in B} b \geq x \iff x \geq \text{- upper bound for } B,
\]

proving (2.53a). Analogously, we obtain (2.53b). Next, (2.53c) and (2.53d) are implied by (2.53a) and (2.53b), respectively. Finally, (2.53e) is proved by

\[
x = \inf_\leq B \iff x = \max_\leq \{ y \in A : y \leq \text{- lower bound for } B \}
\]

\[
\iff x = \min_\geq \{ y \in A : y \geq \text{- upper bound for } B \} \iff x = \sup_\geq B,
\]

and (2.53f) follows analogously. \(\square\)

**Proposition 2.29.** Let \( \leq \) be a partial order on \( A \neq \emptyset, \emptyset \neq B \subseteq A \). The elements \( \max B, \min B, \sup B, \inf B \) are all unique, provided they exist.

**Proof.** Exercise. \(\square\)

**Definition 2.30.** Let \( A, B \) be nonempty sets with partial orders, both denoted by \( \leq \) (even though they might be different). A function \( f : A \rightarrow B \), is called (strictly) isotone, order-preserving, or increasing if, and only if,

\[
\forall_{x,y \in A} (x < y \Rightarrow f(x) \leq f(y)) \quad (\text{resp. } f(x) < f(y)); \quad (2.54a)
\]

\( f \) is called (strictly) antitone, order-reversing, or decreasing if, and only if,

\[
\forall_{x,y \in A} (x < y \Rightarrow f(x) \geq f(y)) \quad (\text{resp. } f(x) > f(y)); \quad (2.54b)
\]

Functions that are (strictly) isotone or antitone are called (strictly) monotone.

**Proposition 2.31.** Let \( A, B \) be nonempty sets with partial orders, both denoted by \( \leq \).

(a) A (strictly) isotone function \( f : A \rightarrow B \) becomes a (strictly) antitone function and vice versa if precisely one of the relations \( \leq \) is replaced by \( \geq \).

(b) If the order \( \leq \) on \( A \) is total and \( f : A \rightarrow B \) is strictly isotone or strictly antitone, then \( f \) is one-to-one.
(c) If the order $\leq$ on $A$ is total and $f : A \rightarrow B$ is invertible and strictly isotone (resp. antitone), then $f^{-1}$ is also strictly isotone (resp. antitone).

**Proof.** (a) is immediate from (2.54).

(b): Due to (a), it suffices to consider the case that $f$ is strictly isotone. If $f$ is strictly isotone and $x \neq y$, then $x < y$ or $y < x$ since the order on $A$ is total. Thus, $f(x) < f(y)$ or $f(y) < f(x)$, i.e. $f(x) \neq f(y)$ in every case, showing $f$ is one-to-one.

(c): Again, due to (a), it suffices to consider the isotone case. If $u, v \in B$ such that $u < v$, then $u = f(f^{-1}(u)), v = f(f^{-1}(v))$, and the isotonicity of $f$ imply $f^{-1}(u) < f^{-1}(v)$ (we are using that the order on $A$ is total – otherwise, $f^{-1}(u)$ and $f^{-1}(v)$ need not be comparable).

**Example 2.32.** (a) $f : \mathbb{N} \rightarrow \mathbb{N}, f(n) := 2n$, is strictly increasing, every constant map on $\mathbb{N}$ is both increasing and decreasing, but not strictly increasing or decreasing. All maps occurring in (2.24) are neither increasing nor decreasing.

(b) The map $f : \mathbb{R} \rightarrow \mathbb{R}, f(x) := -2x$, is invertible and strictly decreasing, and so is $f^{-1} : \mathbb{R} \rightarrow \mathbb{R}, f^{-1}(x) := -x/2$.

(c) The following counterexamples show that the assertions of Prop. 2.31(b),(c) are no longer correct if one does not assume the order on $A$ is total. Let $A$ be the set from (2.50) (where it had been called $B$) with the (nontotal) order from (2.49). The map

$$f : A \rightarrow \mathbb{N}, \begin{cases} f(1, 1) := 1, \\ f(1, 2) := 2, \\ f(2, 1) := 2, \end{cases} \tag{2.55}$$

is strictly isotone, but not one-to-one. The map

$$f : A \rightarrow \{1, 2, 3\}, \begin{cases} f(1, 1) := 1, \\ f(1, 2) := 2, \\ f(2, 1) := 3, \end{cases} \tag{2.56}$$

is strictly isotone and invertible, however $f^{-1}$ is not isotone (since $2 < 3$, but $f^{-1}(2) = (1, 2)$ and $f^{-1}(3) = (2, 1)$ are not comparable, i.e. $f^{-1}(2) \leq f^{-1}(3)$ is not true).

### 3 Natural Numbers, Induction, and the Size of Sets

#### 3.1 Induction and Recursion

One of the most useful proof techniques is the method of induction – it is used in situations, where one needs to verify the truth of statements $\phi(n)$ for each $n \in \mathbb{N}$, i.e. the truth of the statement

$$\forall_{n \in \mathbb{N}} \phi(n). \tag{3.1}$$
Induction is based on the fact that $\mathbb{N}$ satisfies the so-called *Peano axioms*:

**P1:** $\mathbb{N}$ contains a special element called *one*, denoted 1.

**P2:** There exists an injective map $S : \mathbb{N} \rightarrow \mathbb{N}\setminus\{1\}$, called the *successor function* (for each $n \in \mathbb{N}$, $S(n)$ is called the *successor* of $n$).

**P3:** If a subset $A$ of $\mathbb{N}$ has the property that $1 \in A$ and $S(n) \in A$ for each $n \in A$, then $A$ is equal to $\mathbb{N}$. Written as a formula, the third axiom is:

$$\forall A \in P(\mathbb{N}) \left( 1 \in A \land S(A) \subseteq A \Rightarrow A = \mathbb{N} \right).$$

**Remark 3.1.** In Def. 1.26, we had introduced the natural numbers $\mathbb{N} := \{1, 2, 3, \ldots \}$. The successor function is $S(n) = n + 1$. In axiomatic set theory, one starts with the Peano axioms and shows that the axioms of set theory allow the construction of a set $\mathbb{N}$ which satisfies the Peano axioms. One then *defines* $2 := S(1)$, $3 := S(2)$, $\ldots$, $n + 1 := S(n)$. The interested reader can find more details in [Phi16, Sec. D.1].

**Theorem 3.2** (Principle of Induction). *Suppose, for each $n \in \mathbb{N}$, $\phi(n)$ is a statement (i.e. a predicate of $n$ in the language of Def. 1.30). If (a) and (b) both hold, where

(a) $\phi(1)$ is true,

(b) $\forall n \in \mathbb{N} \left( \phi(n) \Rightarrow \phi(n + 1) \right)$,

then (3.1) is true, i.e. $\phi(n)$ is true for every $n \in \mathbb{N}$.

**Proof.** Let $A := \{n \in \mathbb{N} : \phi(n)\}$. We have to show $A = \mathbb{N}$. Since $1 \in A$ by (a), and

$$n \in A \Rightarrow \phi(n) \Rightarrow \phi(n + 1) \Rightarrow S(n) = n + 1 \in A,$$

i.e. $S(A) \subseteq A$, the Peano axiom P3 implies $A = \mathbb{N}$.  

**Remark 3.3.** To prove some $\phi(n)$ for each $n \in \mathbb{N}$ by induction according to Th. 3.2 consists of the following two steps:

(a) Prove $\phi(1)$, the so-called *base case*.

(b) Perform the *inductive step*, i.e. prove that $\phi(n)$ (the *induction hypothesis*) implies $\phi(n + 1)$.

**Example 3.4.** We use induction to prove the statement

$$\forall n \in \mathbb{N} \ (1 + 2 + \cdots + n = \frac{n(n+1)}{2}) :$$

(3.3)
Base Case \((n = 1)\): \(1 = \frac{1+2}{2}\), i.e. \(\phi(1)\) is true.

Induction Hypothesis: Assume \(\phi(n)\), i.e. \(1 + 2 + \cdots + n = \frac{n(n+1)}{2}\) holds.

Induction Step: One computes

\[
1 + 2 + \cdots + n + (n+1) = \frac{n(n+1)}{2} + n + 1 = \frac{n(n+1) + 2n + 2}{2} = \frac{n^2 + 3n + 2}{2} = \frac{(n+1)(n+2)}{2},
\]

i.e. \(\phi(n+1)\) holds and the induction is complete.

**Corollary 3.5.** *Theorem 3.2 remains true if (b) is replaced by

\[
\forall n \in \mathbb{N} \left( \left( \forall_{1 \leq m \leq n} \phi(m) \right) \Rightarrow \phi(n+1) \right).
\]

**Proof.** If, for each \(n \in \mathbb{N}\), we use \(\psi(n)\) to denote \(\forall_{1 \leq m \leq n} \phi(m)\), then (3.5) is equivalent to

\[
\forall_{n \in \mathbb{N}} \left( \psi(n) \Rightarrow \psi(n+1) \right),
\]

i.e. to Th. 3.2(b) with \(\phi\) replaced by \(\psi\). Thus, Th. 3.2 implies \(\psi(n)\) holds true for each \(n \in \mathbb{N}\), i.e. \(\phi(n)\) holds true for each \(n \in \mathbb{N}\). ■

**Corollary 3.6.** Let \(I\) be an index set. Suppose, for each \(i \in I\), \(\phi(i)\) is a statement. If there is a bijective map \(f : \mathbb{N} \rightarrow I\) and (a) and (b) both hold, where

(a) \(\phi(f(1))\) is true,

(b) \(\forall n \in \mathbb{N} \left( \phi(f(n)) \Rightarrow \phi(f(n+1)) \right)\),

then \(\phi(i)\) is true for every \(i \in I\).

**Finite Induction:** The above assertion remains true if \(f : \{1, \ldots, m\} \rightarrow I\) is bijective for some \(m \in \mathbb{N}\) and \(\mathbb{N}\) in (b) is replaced by \(\{1, \ldots, m-1\}\).

**Proof.** If, for each \(n \in \mathbb{N}\), we use \(\psi(n)\) to denote \(\phi(f(n))\), then Th. 3.2 shows \(\psi(n)\) is true for every \(n \in \mathbb{N}\). Given \(i \in I\), we have \(n := f^{-1}(i) \in \mathbb{N}\) with \(f(n) = i\), showing that \(\phi(i) = \phi(f(n)) = \psi(n)\) is true.

For the finite induction, let \(\psi(n)\) denote \((n \leq m \land \phi(f(n))) \lor n > m\). Then, for \(1 \leq n < m\), we have \(\psi(n) \Rightarrow \psi(n+1)\) due to (b). For \(n \geq m\), we also have \(\psi(n) \Rightarrow \psi(n+1)\) due to \(n \geq m \Rightarrow n+1 > m\). Thus, Th. 3.2 shows \(\psi(n)\) is true for every \(n \in \mathbb{N}\). Given \(i \in I\), it is \(n := f^{-1}(i) \in \{1, \ldots, m\}\) with \(f(n) = i\). Since \(n \leq m \land \psi(n) \Rightarrow \phi(f(n))\), we obtain that \(\phi(i)\) is true. ■

Apart from providing a widely employable proof technique, the most important application of Th. 3.2 is the possibility to define sequences inductively, using so-called recursion:
Theorem 3.7 (Recursion Theorem). Let $A$ be a nonempty set and $x \in A$. Given a sequence of functions $(f_n)_{n \in \mathbb{N}}$, where $f_n : A^n \rightarrow A$, there exists a unique sequence $(x_n)_{n \in \mathbb{N}}$ in $A$ satisfying the following two conditions:

(i) $x_1 = x$.

(ii) $\forall n \in \mathbb{N}, x_{n+1} = f_n(x_1, \ldots, x_n)$.

The same holds if $\mathbb{N}$ is replaced by an index set $I$ as in Cor. 3.6.

Proof. To prove uniqueness, let $(x_n)_{n \in \mathbb{N}}$ and $(y_n)_{n \in \mathbb{N}}$ be sequences in $A$, both satisfying (i) and (ii), i.e.

$$
\begin{align*}
    x_1 &= y_1 = x \quad \text{and} \\
    \forall n \in \mathbb{N}, (x_{n+1} &= f_n(x_1, \ldots, x_n) \land y_{n+1} = f_n(y_1, \ldots, y_n)).
\end{align*}
$$

We prove by induction (in the form of Cor. 3.5) that $(x_n)_{n \in \mathbb{N}} = (y_n)_{n \in \mathbb{N}}$, i.e.

$$
\forall n \in \mathbb{N}, x_n = y_n : (3.7)
$$

Base Case $(n = 1)$: $\phi(1)$ is true according to $(3.6a)$.

Induction Hypothesis: Assume $\phi(m)$ for each $m \in \{1, \ldots, n\}$, i.e. $x_m = y_m$ holds for each $m \in \{1, \ldots, n\}$.

Induction Step: One computes

$$
\begin{align*}
    x_{n+1} &\overset{(3.6b)}{=} f_n(x_1, \ldots, x_n) \overset{\phi(1), \ldots, \phi(n)}{=} f_n(y_1, \ldots, y_n) \overset{(3.6b)}{=} y_{n+1},
\end{align*}
$$

i.e. $\phi(n + 1)$ holds and the induction is complete.

To prove existence, we have to show that there is a function $F : \mathbb{N} \rightarrow A$ such that the following two conditions hold:

$$
\begin{align*}
    F(1) &= x, \quad (3.9a) \\
    \forall n \in \mathbb{N}, F(n + 1) &= f_n(F(1), \ldots, F(n)). \quad (3.9b)
\end{align*}
$$

To this end, let

$$
\mathcal{F} := \left\{ B \subseteq \mathbb{N} \times A : (1, x) \in B \land \forall n \in \mathbb{N}, \forall (a_1, \ldots, (n+1, f_n(a_1, \ldots, a_n)) \in B \right\}
$$

and

$$
G := \bigcap_{B \in \mathcal{F}} B.
$$
Note that $G$ is well-defined, as $\mathbb{N} \times A \subseteq \mathcal{F}$. Also, clearly, $G \in \mathcal{F}$. We would like to define $F$ such that $G = \text{graph}(F)$. For this to be possible, we will show, by induction,

$$\forall n \in \mathbb{N} \exists! (n, x_n) \in G.$$  \hspace{1cm} (3.12)

**Base Case ($n = 1$):** From the definition of $G$, we know $(1, x) \in G$. If $(1, a) \in G$ with $a \neq x$, then $H := G \setminus \{(1, a)\} \in \mathcal{F}$, implying $G \subseteq H$ in contradiction to $(1, a) \notin H$. This shows $a = x$ and proves $\phi(1)$.

**Induction Hypothesis:** Assume $\phi(m)$ for each $m \in \{1, \ldots, n\}$.

**Induction Step:** From the induction hypothesis, we know $\exists! (x_1, \ldots, x_n) \in A^n (1, x_1), \ldots, (n, x_n) \in G$.

Thus, if we let $x_{n+1} := f_n(x_1, \ldots, x_n)$, then $(n + 1, x_{n+1}) \in G$ by the definition of $G$. If $(n + 1, a) \in G$ with $a \neq x_{n+1}$, then $H := G \setminus \{(n + 1, a)\} \in \mathcal{F}$ (using the uniqueness of the $(1, x_1), \ldots, (n, x_n) \in G$), implying $G \subseteq H$ in contradiction to $(n + 1, a) \notin H$. This shows $a = x_{n+1}$, proves $\phi(n + 1)$, and completes the induction.

Due to (3.12), we can now define $F : \mathbb{N} \rightarrow A, F(n) := x_n$, and the definition of $G$ then guarantees the validity of (3.9). \hfill \blacksquare

**Example 3.8.** In many applications of Th. 3.7, one has functions $g_n : A \rightarrow A$ and uses

$$\forall n \in \mathbb{N} \left(f_n : A^n \rightarrow A, \quad f_n(a_1, \ldots, a_n) := g_n(a_n)\right).$$  \hspace{1cm} (3.13)

Here are some important concrete examples:

(a) The **factorial function** $F : \mathbb{N}_0 \rightarrow \mathbb{N}, n \mapsto n!$, is defined recursively by

$$0! := 1, \quad 1! := 1, \quad \forall n \in \mathbb{N} \quad (n + 1)! := (n + 1) \cdot n!.$$  \hspace{1cm} (3.14a)

i.e. we have $A = \mathbb{N}$ and $g_n(x) := (n + 1) \cdot x$. So we obtain

$$(n!)_{n \in \mathbb{N}_0} = (1, 1, 2, 6, 24, 120, \ldots).$$  \hspace{1cm} (3.14b)

(b) For each $a \in \mathbb{R}$ and each $d \in \mathbb{R}$, we define the following **arithmetic progression** (also called **arithmetic sequence**) recursively by

$$a_1 := a, \quad \forall n \in \mathbb{N} \quad a_{n+1} := a_n + d.$$  \hspace{1cm} (3.15a)

i.e. we have $A = \mathbb{R}$ and $g_n = g$ with $g(x) := x + d$. For example, for $a = 2$ and $d = -0.5$, we obtain

$$(a_n)_{n \in \mathbb{N}} = (2, 1.5, 1, 0.5, 0, -0.5, -1, -1.5, \ldots).$$  \hspace{1cm} (3.15b)
For each \( a \in \mathbb{R} \) and each \( q \in \mathbb{R} \setminus \{0\} \), we define the following geometric progression (also called geometric sequence) recursively by

\[
x_1 := a, \quad \forall n \in \mathbb{N}, \quad x_{n+1} := x_n \cdot q.
\]

(3.16a)

i.e. we have \( A = \mathbb{R} \) and \( g_n = g \) with \( g(x) := x \cdot q \). For example, for \( a = 3 \) and \( q = -2 \), we obtain

\[
(x_n)_{n \in \mathbb{N}} = (3, -6, 12, -24, 48, \ldots).
\]

(3.16b)

For the time being, we will continue to always specify \( A \) and the \( g_n \) or \( f_n \) in subsequent recursive definitions, but in the literature, most of the time, the \( g_n \) or \( f_n \) are not provided explicitly.

**Example 3.9.** (a) The Fibonacci sequence consists of the Fibonacci numbers, defined recursively by

\[
F_0 := 0, \quad F_1 := 1, \quad \forall n \in \mathbb{N}, \quad F_{n+1} := F_n + F_{n-1}.
\]

(3.17a)

i.e. we have \( A = \mathbb{N}_0 \) and

\[
f_n : A^n \to A, \quad f_n(a_1, \ldots, a_n) := \begin{cases} 1 & \text{for } n = 1, \\ a_n + a_{n-1} & \text{for } n \geq 2. \end{cases}
\]

(3.17b)

So we obtain

\[
(F_n)_{n \in \mathbb{N}_0} = (0, 1, 1, 2, 3, 5, 8, 13, 21, 34, 55, 89, \ldots).
\]

(3.17c)

(b) For \( A := \mathbb{N} \), \( x := 1 \), and

\[
f_n : A^n \to A, \quad f_n(a_1, \ldots, a_n) := a_1 + \cdots + a_n,
\]

(3.18a)

one obtains

\[
x_1 = 1, \quad x_2 = f_1(1) = 1, \quad x_3 = f_2(1, 1) = 2, \quad x_4 = f_3(1, 1, 2) = 4, \\
x_5 = f_4(1, 1, 2, 4) = 8, \quad x_6 = f_5(1, 1, 2, 4, 8) = 16, \ldots
\]

(3.18b)

**Definition 3.10.** (a) **Summation Symbol:** On \( A = \mathbb{R} \) (or, more generally, on every set where an addition \(+ : A \times A \to A\) is defined), define recursively, for each given (possibly finite) sequence \((a_1, a_2, \ldots)\) in \( A \):

\[
\sum_{i=1}^{1} a_i := a_1, \quad \sum_{i=1}^{n+1} a_i := a_{n+1} + \sum_{i=1}^{n} a_i \quad \text{for } n \geq 1,
\]

(3.19a)

i.e.

\[
f_n : A^n \to A, \quad f_n(x_1, \ldots, x_n) := x_n + a_{n+1}.
\]

(3.19b)

In (3.19a), one can also use other symbols for \( i \), except \( a \) and \( n \); for a finite sequence, \( n \) needs to be less than the maximal index of the finite sequence.
More generally, if $I$ is an index set and $\phi : \{1, \ldots, n\} \longrightarrow I$ a bijective map, then define

$$\sum_{i \in I} a_i := \sum_{i=1}^{n} a_{\phi(i)}. \quad (3.19c)$$

The commutativity of addition implies that the definition in (3.19c) is actually independent of the chosen bijective map $\phi$. Also define

$$\sum_{i \in \emptyset} a_i := 0 \quad (3.19d)$$

(for a general $A$, 0 is meant to be an element such that $a + 0 = 0 + a = a$ for each $a \in A$ and we can even define this if $0 \notin A$).

(b) Product Symbol: On $A = \mathbb{R}$ (or, more generally, on every set where a multiplication $\cdot : A \times A \longrightarrow A$ is defined), define recursively, for each given (possibly finite) sequence $(a_1, a_2, \ldots)$ in $A$:

$$\prod_{i=1}^{1} a_i := a_1, \quad \prod_{i=1}^{n+1} a_i := a_{n+1} \cdot \prod_{i=1}^{n} a_i \text{ for } n \geq 1, \quad (3.20a)$$

i.e.

$$f_n : A^n \longrightarrow A, \quad f_n(x_1, \ldots, x_n) := x_n \cdot a_{n+1}. \quad (3.20b)$$

In (3.20a), one can also use other symbols for $i$, except $a$ and $n$; for a finite sequence, $n$ needs to be less than the maximal index of the finite sequence.

More generally, if $I$ is an index set and $\phi : \{1, \ldots, n\} \longrightarrow I$ a bijective map, then define

$$\prod_{i \in I} a_i := \prod_{i=1}^{n} a_{\phi(i)}. \quad (3.20c)$$

The commutativity of multiplication implies that the definition in (3.20c) is actually independent of the chosen bijective map $\phi$. Also define

$$\prod_{i \in \emptyset} a_i := 1 \quad (3.20d)$$

(for a general $A$, 1 is meant to be an element such that $a \cdot 1 = 1 \cdot a = a$ for each $a \in A$ and we can even define this if $1 \notin A$).

Example 3.11. (a) Given $a, d \in \mathbb{R}$, let $(a_n)_{n \in \mathbb{N}}$ be the arithmetic sequence as defined in (3.15a). It is an exercise to prove by induction that

$$\forall n \in \mathbb{N} \quad a_n = a + (n - 1)d, \quad (3.21a)$$

$$\forall n \in \mathbb{N} \quad S_n := \sum_{i=1}^{n} a_i = \frac{n}{2} (a_1 + a_n) = \frac{n}{2} (2a + (n - 1)d), \quad (3.21b)$$

where the $S_n$ are called arithmetic sums.
(b) Given $a \in \mathbb{R}$ and $q \in \mathbb{R} \setminus \{0\}$, let $(x_n)_{n \in \mathbb{N}}$ be the geometric sequence as defined in (3.16a). We will prove by induction that

\[
\forall n \in \mathbb{N} \quad x_n = a q^{n-1},
\]

(3.22a)

For the induction proof of (3.22b), we assume $\phi$ where the $S_n$ are called geometric sums.

For the induction proof of (3.22a), $\phi(n)$ is $x_n = a q^{n-1}$. The base case, $\phi(1)$, is the statement $x_1 = a q^0 = a$, which is true. For the induction step, we assume $\phi(n)$ and compute

\[
x_{n+1} = x_n \cdot q \quad = \quad a q^{n-1} \cdot q = a q^n,
\]

(3.23)

showing $\phi(n) \Rightarrow \phi(n + 1)$ and completing the proof.

For $q = 1$, the sum $S_n$ is actually arithmetic with $d = 0$, i.e. $S_n = n a$ can be obtained from (3.21b). For the induction proof of (3.22b) with $q \neq 1$, $\phi(n)$ is $S_n = \frac{a(1-q^n)}{1-q}$. The base case, $\phi(1)$, is the statement $S_1 = \frac{a(1-q)}{1-q} = a$, which is true. For the induction step, we assume $\phi(n)$ and compute

\[
S_{n+1} = S_n + x_{n+1} = \frac{a(1-q^n)}{1-q} + a q^n = \frac{a (1-q^n) + a q^n (1-q)}{1-q} = \frac{a(1-q^{n+1})}{1-q},
\]

(3.24)

showing $\phi(n) \Rightarrow \phi(n + 1)$ and completing the proof.

### 3.2 Cardinality: The Size of Sets

Cardinality measures the size of sets. For a finite set $A$, it is precisely the number of elements in $A$. For an infinite set, it classifies the set’s degree or level of infinity (it turns out that not all infinite sets have the same size).

**Definition 3.12.** (a) The sets $A, B$ are defined to have the same cardinality or the same size if, and only if, there exists a bijective map $\varphi : A \rightarrow B$. One can show that this defines an equivalence relation on every set of sets (see [Phi16, Th. A.53]).

(b) The cardinality of a set $A$ is $n \in \mathbb{N}$ (denoted $\#A = n$) if, and only if, there exists a bijective map $\varphi : A \rightarrow \{1, \ldots, n\}$. The cardinality of $\emptyset$ is defined as 0, i.e. $\#\emptyset := 0$. A set $A$ is called finite if, and only if, there exists $n \in \mathbb{N}_0$ such that $\#A = n$; $A$ is called infinite if, and only if. $A$ is not finite, denoted $\#A = \infty$ (in the strict sense, this is an abuse of notation, since $\infty$ is not a cardinality – for example $\#\mathbb{N} = \infty$ and $\#\mathcal{P} (\mathbb{N}) = \infty$, but $\mathbb{N}$ and $\mathcal{P}(\mathbb{N})$ do not have the same cardinality, since the power set $\mathcal{P}(A)$ is always strictly bigger than $A$ (see Th. 3.20 below) – $\#A = \infty$ is merely an abbreviation for the statement “$A$ is infinite”). The interested student finds additional material regarding the uniqueness of finite cardinality in [Phi16, Th. A.61] and [Phi16, Cor. A.62], and regarding characterizations of infinite sets in [Phi16, Th. A.54].
3 NATURAL NUMBERS, INDUCTION, AND THE SIZE OF SETS

The set $A$ is called \textit{countable} if, and only if, $A$ is finite or $A$ has the same cardinality as $\mathbb{N}$. Otherwise, $A$ is called \textit{uncountable}.

\textbf{Theorem 3.13.} Let $A \neq \emptyset$ be a finite set.

\begin{itemize}
  \item[(a)] If $B \subseteq A$ with $A \neq B$, then $B$ is finite with $\#B < \#A$.
  \item[(b)] If $a \in A$, then $\#(A \setminus \{a\}) = \#A - 1$.
\end{itemize}

\textit{Proof.} For $\#A = n \in \mathbb{N}$, we use induction to prove (a) and (b) simultaneously, i.e. we show

$$
\forall n \in \mathbb{N} \left( \#A = n \implies \forall B \in \mathcal{P}(A) \setminus \{A\} \forall a \in A \#B \in \{0, \ldots, n-1\} \land \#(A \setminus \{a\}) = n - 1 \right).
$$

\textbf{Base Case} ($n = 1$): In this case, $A$ has precisely one element, i.e. $B = A \setminus \{a\} = \emptyset$, and $\#\emptyset = 0 = n - 1$ proves $\phi(1)$.

\textbf{Induction Step:} For the induction hypothesis, we assume $\phi(n)$ to be true, i.e. we assume (a) and (b) hold for each $A$ with $\#A = n$. We have to prove $\phi(n + 1)$, i.e., we consider $A$ with $\#A = n + 1$. From $\#A = n + 1$, we conclude the existence of a bijective map $\varphi : A \rightarrow \{1, \ldots, n+1\}$. We have to construct a bijective map $\psi : A \setminus \{a\} \rightarrow \{1, \ldots, n\}$.

To this end, set $k := \varphi(a)$ and define the auxiliary function

$$
f : \{1, \ldots, n+1\} \rightarrow \{1, \ldots, n+1\}, \quad f(x) := \begin{cases} n + 1 & \text{for } x = k, \\
k & \text{for } x = n + 1, \\
x & \text{for } x \notin \{k, n + 1\}. \end{cases}
$$

Then $f \circ \varphi : A \rightarrow \{1, \ldots, n+1\}$ is bijective by Th. 2.13, and

$$(f \circ \varphi)(a) = f(\varphi(a)) = f(k) = n + 1.$$  

Thus, the restriction $\psi := (f \circ \varphi) \upharpoonright A \setminus \{a\}$ is the desired bijective map $\psi : A \setminus \{a\} \rightarrow \{1, \ldots, n\}$, proving $\#(A \setminus \{a\}) = n$. It remains to consider the strict subset $B$ of $A$.

Since $B$ is a strict subset of $A$, there exists $a \in A \setminus B$. Thus, $B \subseteq A \setminus \{a\}$ and, as we have already shown $\#(A \setminus \{a\}) = n$, the induction hypothesis applies and yields $B$ is finite with $\#B \leq \#(A \setminus \{a\}) = n$, i.e. $\#B \in \{0, \ldots, n\}$, proving $\phi(n + 1)$, thereby completing the induction.

\textbf{Theorem 3.14.} For $\#A = \#B = n \in \mathbb{N}$ and $f : A \rightarrow B$, the following statements are equivalent:

\begin{itemize}
  \item[(i)] $f$ is injective.
  \item[(ii)] $f$ is surjective.
  \item[(iii)] $f$ is bijective.
\end{itemize}
Proof. It suffices to prove the equivalence of (i) and (ii).

If $f$ is injective, then $f : A \to f(A)$ is bijective. Since $\#A = n$, there exists a bijective map $\varphi : A \to \{1, \ldots, n\}$. Then $(\varphi \circ f^{-1}) : f(A) \to \{1, \ldots, n\}$ is also bijective, showing $\#f(A) = n$, i.e., according to Th. 3.13(a), $f(A)$ can not be a strict subset of $B$, i.e. $f(A) = B$, proving $f$ is surjective.

If $f$ is surjective, then $f$ has a right inverse $g : B \to A$ by Th. 2.12(a), i.e. $f \circ g = \text{Id}_B$. But this also means $f$ is a left inverse for $g$, such that $g$ must be injective by Th. 2.12(b). According to what we have already proved above, $g$ injective implies $g$ surjective, i.e. $g$ must be bijective. From Th. 2.12(c), we then know the left inverse of $g$ is unique, implying $f = g^{-1}$. In particular, $f$ is injective.

\[\square\]

Lemma 3.15. For each finite set $A$ (i.e. $\#A = n \in \mathbb{N}_0$) and each $B \subseteq A$, one has $\#(A \setminus B) = \#A - \#B$.

Proof. For $B = \emptyset$, the assertion is true since $\#(A \setminus B) = \#A = \#A - 0 = \#A - \#B$.

For $B \neq \emptyset$, the proof is conducted over the size of $B$, i.e. as a finite induction (cf. Cor. 3.6) over the set $\{1, \ldots, n\}$, showing

\[
\forall m \in \{1, \ldots, n\} \quad (\#B = m \Rightarrow \#(A \setminus B) = \#A - \#B).
\]

Base Case ($m = 1$): $\phi(1)$ is precisely the statement provided by Th. 3.13(b).

Induction Step: For the induction hypothesis, we assume $\phi(m)$ with $1 \leq m < n$. To prove $\phi(m + 1)$, consider $B \subseteq A$ with $\#B = m + 1$. Fix an element $b \in B$ and set $B_1 := B \setminus \{b\}$. Then $\#B_1 = m$ by Th. 3.13(b), $A \setminus B = (A \setminus B_1) \setminus \{b\}$, and we compute

\[
\#(A \setminus B) = \#((A \setminus B_1) \setminus \{b\}) \overset{\text{Th. 3.13(b)}}{=} \#(A \setminus B_1) - 1 \overset{\phi(m)}{=} \#A - \#B_1 - 1 = \#A - \#B,
\]

proving $\phi(m + 1)$ and completing the induction. \[\square\]

Theorem 3.16. If $A, B$ are finite sets, then $\#(A \cup B) = \#A + \#B - \#(A \cap B)$.

Proof. The assertion is clearly true if $A$ or $B$ is empty. If $A$ and $B$ are nonempty, then there exist $m, n \in \mathbb{N}$ such that $\#A = m$ and $\#B = n$, i.e. there are bijective maps $f : A \to \{1, \ldots, m\}$ and $g : B \to \{1, \ldots, n\}$.

We first consider the case $A \cap B = \emptyset$. We need to construct a bijective map $h : A \cup B \to \{1, \ldots, m + n\}$. To this end, we define

\[
h : A \cup B \to \{1, \ldots, m + n\}, \quad h(x) := \begin{cases} f(x) & \text{for } x \in A, \\ g(x) + m & \text{for } x \in B. \end{cases}
\]

The bijectivity of $f$ and $g$ clearly implies the bijectivity of $h$, proving $\#(A \cup B) = m + n = \#A + \#B$. 

\[\square\]
Finally, we consider the case of arbitrary $A, B$. Since $A \cup B = A \cup (B \setminus A)$ and $B \setminus A = B \setminus (A \cap B)$, we can compute
\[
\#(A \cup B) = \#(A \cup (B \setminus A)) = \#A + \#(B \setminus A)
\]
\[= \#A + \#(B \setminus (A \cap B)) \overset{\text{Lem. 3.15}}{=} \#A + \#B - \#(A \cap B),
\]
thereby establishing the case.

\[\text{Theorem 3.17.} \quad \text{If } (A_1, \ldots, A_n), \ n \in \mathbb{N}, \text{ is a finite sequence of finite sets, then}
\]
\[
\# \prod_{i=1}^{n} A_i = \#(A_1 \times \cdots \times A_n) = \prod_{i=1}^{n} \#A_i.
\]

\[\text{Proof.} \quad \text{If at least one } A_i \text{ is empty, then (3.25) is true, since both sides are 0.}
\]
The case where all $A_i$ are nonempty is proved by induction over $n$, i.e. we know $k_i := \#A_i \in \mathbb{N}$ for each $i \in \{1, \ldots, n\}$ and show by induction
\[
\forall n \in \mathbb{N} \quad \# \prod_{i=1}^{n} A_i = \prod_{i=1}^{n} k_i.
\]

Base Case ($n = 1$): $\prod_{i=1}^{1} A_i = \#A_1 = k_1 = \prod_{i=1}^{1} k_i$, i.e. $\phi(1)$ holds.

Induction Step: From the induction hypothesis $\phi(n)$, we obtain a bijective map $\varphi : A \to \{1, \ldots, N\}$, where $A := \prod_{i=1}^{n} A_i$ and $N := \prod_{i=1}^{n} k_i$. To prove $\phi(n + 1)$, we need to construct a bijective map $h : A \times A_{n+1} \to \{1, \ldots, N \cdot k_{n+1}\}$. Since $\#A_{n+1} = k_{n+1}$, there exists a bijective map $f : A_{n+1} \to \{1, \ldots, k_{n+1}\}$. We define
\[
h : A \times A_{n+1} \to \{1, \ldots, N \cdot k_{n+1}\},
\]
\[h(a_1, \ldots, a_n, a_{n+1}) := (f(a_{n+1}) - 1) \cdot N + \varphi(a_1, \ldots, a_n).
\]
Since $\varphi$ and $f$ are bijective, and since every $m \in \{1, \ldots, N \cdot k_{n+1}\}$ has a unique representation in the form $m = a \cdot N + r$ with $a \in \{0, \ldots, k_{n+1} - 1\}$ and $r \in \{1, \ldots, N\}$ (exercise), $h$ is also bijective. This proves $\phi(n+1)$ and completes the induction.

\[\text{Theorem 3.18.} \quad \text{For each finite set } A \text{ (i.e. } \#A = n \in \mathbb{N}_0\), one has } \#\mathcal{P}(A) = 2^n.
\]

\[\text{Proof.} \quad \text{The proof is conducted by induction by showing}
\]
\[
\forall n \in \mathbb{N}_0 \left( \#A = n \implies \#\mathcal{P}(A) = 2^n \right).
\]

Base Case ($n = 0$): For $n = 0$, we have $A = \emptyset$, i.e. $\mathcal{P}(A) = \{\emptyset\}$. Thus, $\#\mathcal{P}(A) = 1 = 2^0$, proving $\phi(0)$.

Induction Step: Assume $\phi(n)$ and consider $A$ with $\#A = n + 1$. Then $A$ contains at least one element $a$. For $B := A \setminus \{a\}$, we then know $\#B = n$ from Th. 3.13(b).
Moreover, setting $\mathcal{M} := \{ C \cup \{a\} : C \in \mathcal{P}(B) \}$, we have the disjoint decomposition $\mathcal{P}(A) = \mathcal{P}(B) \cup \mathcal{M}$. As the map $\varphi : \mathcal{P}(B) \to \mathcal{M}$, $\varphi(C) := C \cup \{a\}$, is clearly bijective, $\mathcal{P}(B)$ and $\mathcal{M}$ have the same cardinality. Thus,

$$\#\mathcal{P}(A) \stackrel{\text{Th. 3.16}}{=} \#\mathcal{P}(B) + \#\mathcal{M} = \#\mathcal{P}(B) + \#\mathcal{P}(B) \left( \phi(n) \right) = 2 \cdot 2^n = 2^{n+1},$$

thereby proving $\phi(n+1)$ and completing the induction.

**Remark 3.19.** In the proof of the following Th. 3.20, we will encounter a new proof technique that we did not use before, the so-called proof by contradiction, also called indirect proof. It is based on the observation, called the principle of contradiction, that $A \land \neg A$ is always false:

$$\begin{array}{c|cc}
A & \neg A & A \land \neg A \\
T & F & F \\
F & T & F
\end{array}
$$

Thus, one possibility of proving a statement $B$ to be true is to show $\neg B \Rightarrow A \land \neg A$ for some arbitrary statement $A$. Since the right-hand side of the implication is false, the left-hand side must also be false, proving $B$ is true.

**Theorem 3.20.** Let $A$ be a set. There can never exist a surjective map from $A$ onto $\mathcal{P}(A)$ (in this sense, the size of $\mathcal{P}(A)$ is always strictly bigger than the size of $A$; in particular, $A$ and $\mathcal{P}(A)$ can never have the same size).

**Proof.** If $A = \emptyset$, then there is nothing to prove. For nonempty $A$, as mentioned above, the idea is to conduct a proof by contradiction. To this end, assume there does exist a surjective map $f : A \to \mathcal{P}(A)$ and define

$$B := \{ x \in A : x \notin f(x) \}.$$  

(3.27)

Now $B$ is a subset of $A$, i.e. $B \in \mathcal{P}(A)$ and the assumption that $f$ is surjective implies the existence of $a \in A$ such that $f(a) = B$. If $a \in B$, then $a \notin f(a) = B$, i.e. $a \in B$ implies $a \in B \land \neg (a \in B)$, so that the principle of contradiction tells us $a \notin B$ must be true. However, $a \notin B$ implies $a \in f(a) = B$, i.e., this time, the principle of contradiction tells us $a \in B$ must be true. In conclusion, we have shown our original assumption that there exists a surjective map $f : A \to \mathcal{P}(A)$ implies $a \in B \land \neg (a \in B)$, i.e., according to the principle of contradiction, no surjective map from $A$ into $\mathcal{P}(A)$ can exist. ■

We conclude the section with a number of important results regarding the natural numbers and countability.

**Theorem 3.21.** (a) Every nonempty finite subset of a totally ordered set has a minimum and a maximum.

(b) Every nonempty subset of $\mathbb{N}$ has a minimum.
Proof. (a): Let $A$ be a set and let $\leq$ denote a total order on $A$. Moreover, let $\emptyset \neq B \subseteq A$. We show by induction

$$\forall n \in \mathbb{N} \ (\#B = n \Rightarrow B \text{ has a min}).$$

Base Case ($n = 1$): For $n = 1$, $B$ contains a unique element $b$, i.e. $b = \min B$, proving $\phi(1)$.

Induction Step: Suppose $\phi(n)$ holds and consider $B$ with $\#B = n + 1$. Let $b$ be one element from $B$. Then $C := B \setminus \{b\}$ has cardinality $n$ and, according to the induction hypothesis, there exists $c \in C$ satisfying $c = \min C$. If $c \leq b$, then $c \leq x$ for each $x \in B$, proving $c = \min B$. If $b \leq c$, then $b \leq x$ for each $x \in B$, proving $b = \min B$. In each case, $B$ has a min, proving $\phi(n + 1)$ and completing the induction.

(b): Let $\emptyset \neq A \subseteq \mathbb{N}$. We have to show $A$ has a min. If $A$ is finite, then $A$ has a min by (a). If $A$ is infinite, let $n$ be an element from $A$. Then the finite set $B := \{k \in A : k \leq n\}$ must have a min $m$ by (a). Since $m \leq x$ for each $x \in B$ and $m \leq n < x$ for each $x \in A \setminus B$, we have $m = \min A$.

**Proposition 3.22.** Every subset $A$ of $\mathbb{N}$ is countable.

Proof. Since $\emptyset$ is countable, we may assume $A \neq \emptyset$. From Th. 3.21(b), we know that every nonempty subset of $\mathbb{N}$ has a min. We recursively define a sequence in $A$ by

$$a_1 := \min A, \quad a_{n+1} := \begin{cases} \min A_n & \text{if } A_n := A \setminus \{a_i : 1 \leq i \leq n\} \neq \emptyset, \\ a_n & \text{if } A_n = \emptyset. \end{cases}$$

This sequence is the same as the function $f : \mathbb{N} \to A, f(n) = a_n$. An easy induction shows that, for each $n \in \mathbb{N}$, $a_n \neq a_{n+1}$ implies the restriction $f |_{\{1, \ldots, n+1\}}$ is injective. Thus, if there exists $n \in \mathbb{N}$ such that $a_n = a_{n+1}$, then $f |_{\{1, \ldots, k\}} : \{1, \ldots, k\} \to A$ is bijective, where $k := \min\{n \in \mathbb{N} : a_n = a_{n+1}\}$, showing $A$ is finite, i.e. countable. If there does not exist $n \in \mathbb{N}$ with $a_n = a_{n+1}$, then $f$ is injective. Another easy induction shows that, for each $n \in \mathbb{N}$, $f(\{1, \ldots, n\}) \supseteq \{k \in A : k \leq n\}$, showing $f$ is also surjective, proving $A$ is countable.

**Proposition 3.23.** For each set $A \neq \emptyset$, the following three statements are equivalent:

(i) $A$ is countable.

(ii) There exists an injective map $f : A \to \mathbb{N}$.

(iii) There exists a surjective map $g : \mathbb{N} \to A$.

Proof. Directly from the definition of countable in Def. 3.12(c), one obtains (i)$\Rightarrow$(ii) and (i)$\Rightarrow$(iii). To prove (ii)$\Rightarrow$(i), let $f : A \to \mathbb{N}$ be injective. Then $f : A \to f(A)$ is bijective, and, since $f(A) \subseteq \mathbb{N}$, $f(A)$ is countable by Prop. 3.22, proving $A$ is countable as well. To prove (iii)$\Rightarrow$(i), let $g : \mathbb{N} \to A$ be surjective. According to Th. 2.12(a), $g$ has a right inverse $f : A \to \mathbb{N}$, i.e. $g \circ f = \text{Id}_A$. But this means $g$ is a left inverse for $f$, showing $f$ is injective according to Th. 2.12(b). Then $A$ is countable by an application of (ii).
Theorem 3.24. If \((A_1, \ldots, A_n), n \in \mathbb{N},\) is a finite family of countable sets, then \(\prod_{i=1}^{n} A_i\) is countable.

Proof. We first consider the special case \(n = 2\) with \(A_1 = A_2 = \mathbb{N}\) and show the map

\[ \varphi : \mathbb{N} \times \mathbb{N} \rightarrow \mathbb{N}, \quad \varphi(m, n) := 2^m \cdot 3^n, \]

is injective: If \(\varphi(m, n) = \varphi(p, q),\) then \(2^m \cdot 3^n = 2^p \cdot 3^q.\) Moreover \(m \leq p\) or \(p \leq m.\) If \(m \leq p,\) then \(3^n = 2^{p-m} \cdot 3^q.\) Since \(3^n\) is odd, \(2^{p-m} \cdot 3^q\) must also be odd, implying \(p - m = 0, i.e. m = p.\) Moreover, we now have \(3^n = 3^q,\) implying \(n = q,\) showing \((m, n) = (p, q), i.e. \varphi\) is injective.

We now come back to the general case stated in the theorem. If at least one of the \(A_i\) is empty, then \(A\) is empty. So it remains to consider the case, where all \(A_i\) are nonempty. The proof is conducted by induction by showing

\[ \forall n \in \mathbb{N} \prod_{i=1}^{n} A_i \text{ is countable.} \]

Base Case \((n = 1): \phi(1)\) is merely the hypothesis that \(A_1\) is countable.

Induction Step: Assuming \(\phi(n),\) Prop. 3.23(ii) provides injective maps \(f_1 : \prod_{i=1}^{n} A_i \rightarrow \mathbb{N}\) and \(f_2 : A_{n+1} \rightarrow \mathbb{N}.\) To prove \(\phi(n+1),\) we provide an injective map \(h : \prod_{i=1}^{n+1} A_i \rightarrow \mathbb{N}: \)

\[ h : \prod_{i=1}^{n+1} A_i \rightarrow \mathbb{N}, \quad h(a_1, \ldots, a_n, a_{n+1}) := \varphi(f_1(a_1, \ldots, a_n), f_2(a_{n+1})). \]

The injectivity of \(f_1, f_2,\) and \(\varphi\) clearly implies the injectivity of \(h,\) thereby proving \(\phi(n + 1)\) and completing the induction. ■

Theorem 3.25. If \((A_i)_{i \in I}\) is a countable family of countable sets (i.e. \(\emptyset \neq I\) is countable and each \(A_i, i \in I,\) is countable), then the union \(A := \bigcup_{i \in I} A_i\) is also countable.

Proof. It suffices to consider the case that all \(A_i\) are nonempty. Moreover, according to Prop. 3.23(iii), it suffices to construct a surjective map \(\varphi : \mathbb{N} \rightarrow A.\) Also according to Prop. 3.23(iii), the countability of \(I\) and the \(A_i\) provides us with surjective maps \(f : \mathbb{N} \rightarrow I\) and \(g_i : \mathbb{N} \rightarrow A_i.\) Define

\[ F : \mathbb{N} \times \mathbb{N} \rightarrow A, \quad F(m, n) := g_{f(m)}(n). \]

Then \(F\) is surjective: Given \(x \in A,\) there exists \(i \in I\) such that \(x \in A_i.\) Since \(f\) is surjective, there is \(m \in \mathbb{N}\) satisfying \(f(m) = i.\) Moreover, since \(g_i\) is surjective, there exists \(n \in \mathbb{N}\) with \(g_i(n) = x.\) Then \(F(m, n) = g_i(n) = x,\) verifying that \(F\) is surjective. As \(\mathbb{N} \times \mathbb{N}\) is countable by Th. 3.24, there exists a surjective map \(h : \mathbb{N} \rightarrow \mathbb{N} \times \mathbb{N}.\) Thus, \(F \circ h\) is the desired surjective map from \(\mathbb{N}\) onto \(A.\) Note: The axiom of choice (AC, see [Phi16, Sec. A.4]) is used when choosing each \(g_i\) from the set of all surjective maps from \(\mathbb{N}\) onto \(A_i.\) It has actually been shown that it is impossible to prove the theorem without using AC (cf. [Phi16, Rem. 3.18]). ■
4 Real Numbers

4.1 The Real Numbers as a Complete Totally Ordered Field

The set of real numbers, denoted \( \mathbb{R} \), is a set with special properties, namely a so-called *complete totally ordered field*. We already know what totally ordered means, but we still need to explain what a field is, what an ordered field is, and what it means for a total order to be complete. We begin with the last part.

**Definition 4.1.** A total order \( \leq \) on a nonempty set \( A \) is called *complete* if, and only if, every nonempty subset \( B \) of \( A \) that is bounded from above has a supremum, i.e.

\[
\forall B \in \mathcal{P}(A) \setminus \{\emptyset\} \left( \left( \exists x \in A \forall b \in B \ b \leq x \right) \Rightarrow \exists s \in A s = \sup B \right). \tag{4.1}
\]

**Lemma 4.2.** A total order \( \leq \) on a nonempty set \( A \) is complete if, and only if, every nonempty subset \( B \) of \( A \) that is bounded from below has an infimum.

**Proof.** According to Lem. 2.28, it suffices to prove one implication. We show that (4.1) implies that every nonempty \( B \) bounded from below has an infimum: Define

\[
C := \{x \in A : x \text{ is lower bound for } B\}. \tag{4.2}
\]

Then every \( b \in B \) is an upper bound for \( C \) and (4.1) implies there exists \( s = \sup C \in A \). To verify \( s = \inf B \), it remains to show \( s \in C \), i.e. that \( s \) is a lower bound for \( B \). However, every \( b \in B \) is an upper bound for \( C \) and \( s = \sup C \) is the min of all upper bounds for \( C \), i.e. \( s \leq b \) for each \( b \in B \), showing \( s \in C \). \( \blacksquare \)

**Definition 4.3.** Let \( A \) be a nonempty set with a map \( \circ : A \times A \rightarrow A, \ (x, y) \mapsto x \circ y \) (called a *composition* on \( A \), the examples we have in mind are addition and multiplication on \( \mathbb{R} \)). Then \( A \) is called a *group* with respect to \( \circ \) if, and only if, the following three conditions are satisfied:

1. **Associativity:** \( x \circ (y \circ z) = (x \circ y) \circ z \) holds for all \( x, y, z \in A \).
2. **There exists a neutral element** \( e \in A \), i.e. an element \( e \in A \) such that
   \[
   \forall x \in A \ x \circ e = x.
   \]
3. **For each** \( x \in A \), **there exists an inverse element** \( \overline{x} \in A \), i.e. an element \( \overline{x} \in A \) such that
   \[
   x \circ \overline{x} = e.
   \]

\( A \) is called a *commutative* or *abelian* group if, and only if, it is a group and satisfies the additional condition:
(iv) Commutativity: $x \circ y = y \circ x$ holds for all $x, y \in A$.

**Definition 4.4.** Let $A$ be a nonempty set with two maps

$$+ : A \times A \rightarrow A, \quad (x, y) \mapsto x + y,$$

$$\cdot : A \times A \rightarrow A, \quad (x, y) \mapsto x \cdot y$$

(+ is called addition and $\cdot$ is called multiplication; often one writes $xy$ instead of $x \cdot y$).

Then $A$ is called a **field** if, and only if, the following three conditions are satisfied:

(i) $A$ is a commutative group with respect to $+$. The neutral element with respect to $+$ is denoted $0$.

(ii) $A \setminus \{0\}$ is a commutative group with respect to $\cdot$. The neutral element with respect to $\cdot$ is denoted $1$.

(iii) Distributivity:

$$\forall x, y, z \in A \quad x \cdot (y + z) = x \cdot y + x \cdot z. \quad (4.5)$$

If $A$ is a field and $\leq$ is a total order on $A$, then $A$ is called a **totally ordered field** if, and only if, the following condition is satisfied:

(iv) Compatibility with Addition and Multiplication:

$$\forall x, y, z \in A \quad x \leq y \Rightarrow x + z \leq y + z, \quad (4.6a)$$

$$\forall x, y \in A \quad 0 \leq x \wedge 0 \leq y \Rightarrow 0 \leq xy. \quad (4.6b)$$

Finally, $A$ is called a **complete totally ordered field** if, and only if, $A$ is a totally ordered field that is complete in the sense of Def. 4.1.

**Theorem 4.5.** There exists a complete totally ordered field $\mathbb{R}$ (it is called the set of real numbers). Moreover, $\mathbb{R}$ is unique up to isomorphism, i.e. if $A$ is a complete totally ordered field, then there exists an isomorphism $\phi : A \rightarrow \mathbb{R}$, i.e. a bijective map $\phi : A \rightarrow \mathbb{R}$, satisfying

$$\forall x, y \in A \quad \phi(x + y) = \phi(x) + \phi(y), \quad (4.7a)$$

$$\forall x, y \in A \quad \phi(xy) = \phi(x)\phi(y), \quad (4.7b)$$

$$\forall x, y \in A \quad (x < y \Rightarrow \phi(x) < \phi(y)). \quad (4.7c)$$

It also turns out that the isomorphism is unique.

**Proof.** To really prove the existence of the real numbers by providing a construction is tedious and not easy. One possible construction is provided in [Phi16, Sec. D] (the existence proof is completed in [Phi16, Th. D.41], the results regarding the isomorphism can be found in [Phi16, Th. D.45]).
Theorem 4.6. The following statements and rules are valid in the set of real numbers \( \mathbb{R} \) (and, more generally, in every field):

(a) Inverse elements are unique. For each \( x \in \mathbb{R} \), the unique inverse with respect to addition is denoted by \(-x\). Also define \( y - x := y + (-x) \). For each \( x \in \mathbb{R} \setminus \{0\} \), the unique inverse with respect to multiplication is denoted by \( x^{-1} \). For \( x \neq 0 \), define the fractions \( \frac{y}{x} := y/x := yx^{-1} \) with numerator \( y \) and denominator \( x \).

(b) \(-(-x) = x\) and \((x^{-1})^{-1} = x\) for \( x \neq 0\).

(c) \((-x) + (-y) = -(x + y)\) and \(x^{-1}y^{-1} = (xy)^{-1}\) for \( x, y \neq 0\).

(d) \(x + a = y + a \Rightarrow x = y\) and, for \( a \neq 0\), \( xa = ya \Rightarrow x = y\).

(e) \(x \cdot 0 = 0\).

(f) \(x(-y) = -(xy)\).

(g) \((-x)(-y) = xy\).

(h) \(x(y - z) = xy - xz\).

(i) \(xy = 0 \Rightarrow x = 0 \vee y = 0\).

(j) Rules for Fractions:

\[
\frac{a}{c} + \frac{b}{d} = \frac{ad + bc}{cd}, \quad \frac{a}{c} \cdot \frac{b}{d} = \frac{ab}{cd}, \quad \frac{a/c}{b/d} = \frac{ad}{bc},
\]

where all denominators are assumed \( \neq 0\).

Proof. (a): Let \( a, b \) be additive inverses to \( x \). Then \( a = a + 0 = a + x + b = 0 + b = b \). The multiplicative case is proved completely analogously.

(b): \(-x + x = 0\) already shows that \( x \) is the inverse to \(-x\), i.e. \(-(-x) = x\). The multiplicative case is proved completely analogously.

(c): \(x + y + (-x) + (-y) = x - x + y - y = 0\), showing \((-x) + (-y)\) is the inverse to \((x + y)\). The multiplicative case is proved completely analogously.

(d): If \(x + a = y + a\), then \(x = x + a - a = y + a - a = y\). Again, the multiplicative case is proved completely analogously.

(e): One computes

\[
x \cdot 0 + x \cdot 1 \overset{(4.5)}{=} x \cdot (0 + 1) = x \cdot 1 = 0 + x \cdot 1,
\]

i.e. \(x \cdot 0 = 0\) follows from (d).

(f): \(xy + x(-y) = x(y - y) = x \cdot 0 = 0\), where we used (4.5) and (e). This shows \(x(-y)\) is the additive inverse to \(xy\).
(g): $xy = -(-(xy)) = -(x(-y)) = -((-y)x) = (-y)(-x)$, where (f) was used twice.

(h): $x(y - z) = x(y + (-z)) = xy + x(-z) = xy - xz$.

(i): If $xy = 0$ and $x \neq 0$, then $y = 1 \cdot y = x^{-1}xy = x^{-1} \cdot 0 = 0$.

(j): One computes
\[ \frac{a}{c} + \frac{b}{d} = ac^{-1} + bd^{-1} = ad^{-1}c^{-1} + bcc^{-1}d^{-1} = (ad + bc)(cd)^{-1} = \frac{ad + bc}{cd} \]
and
\[ \frac{a}{c} \cdot \frac{b}{d} = ac^{-1}bd^{-1} = ab(cd)^{-1} = \frac{ab}{cd} \]
and
\[ \frac{a/c}{b/d} = ac^{-1}(bd^{-1})^{-1} = ac^{-1}b^{-1}d = ad(bc)^{-1} = \frac{ad}{bc} , \]
completing the proof.

\[ \text{Theorem 4.7.} \] The following statements and rules are valid in the set of real numbers $\mathbb{R}$ (and, more generally, in every totally ordered field):

(a) $x \leq y \Rightarrow -x \geq -y$.

(b) $x \leq y \land z \geq 0 \Rightarrow xz \leq yz$ holds as well as $x \leq y \land z \leq 0 \Rightarrow xz \geq yz$.

(c) $x \neq 0 \Rightarrow x^2 := x \cdot x > 0$. In particular $1 > 0$.

(d) $x > 0 \Rightarrow 1/x > 0$, whereas $x < 0 \Rightarrow 1/x < 0$.

(e) If $0 < x < y$, then $x/y < 1$, $y/x > 1$, and $1/x > 1/y$.

(f) $x < y \land u < v \Rightarrow x + u < y + v$.

(g) $0 < x < y \land 0 < u < v \Rightarrow xu < yv$.

(h) $x < y \land 0 < \lambda < 1 \Rightarrow x < \lambda x + (1 - \lambda)y < y$. In particular $x < \frac{x + y}{2} < y$.

Proof. (a): Using (4.6a): $x \leq y \Rightarrow 0 \leq y - x \Rightarrow -y \leq -x$.

(b): One argues, for $z \geq 0$,
\[ x \leq y \Rightarrow 0 \leq y - x \overset{(4.6b)}{\Rightarrow} 0 \leq (y - x)z = yz - xz \Rightarrow xz \leq yz, \]
and, for $z \leq 0$,
\[ x \leq y \Rightarrow 0 \leq y - x \overset{(4.6b)}{\Rightarrow} 0 \leq (y - x)(-z) = xz - yz \Rightarrow xz \geq yz. \]

(c): From (4.6b), one obtains $x^2 \geq 0$. From Th. 4.6(i), one then gets $x^2 > 0$.

(d): If $x > 0$, then $x^{-1} < 0$ implies the false statement $1 = xx^{-1} < 0$, i.e. $x^{-1} > 0$. The case $x < 0$ is treated analogously.
(e): Using (d), we obtain from \(0 < x < y\) that \(\frac{x}{y} = xy^{-1} < yy^{-1} = 1\) and \(1 = xx^{-1} < yx^{-1} = y/x\).

(f): \(x < y \Rightarrow x + u < y + u\) and \(u < v \Rightarrow y + u < y + v\); both combined yield \(x + u < y + v\).

(g): \(0 < x < y \land 0 < u < v \Rightarrow xu < yu \land yu < yv \Rightarrow xu < yv\).

(h): Since \(0 < \lambda\) and \(1 - \lambda > 0\), \(x < y\) implies
\[
\lambda x < \lambda y \quad \land \quad (1-\lambda)x < (1-\lambda)y.
\]
Using (4.6a), we obtain
\[
x = \lambda x + (1-\lambda)x < \lambda x + (1-\lambda)y < \lambda y + (1-\lambda)y = y,
\]
completing the proof of the theorem.

\[\Box\]

**Theorem 4.8.** Let \(\emptyset \neq A, B \subseteq \mathbb{R}, \lambda \in \mathbb{R}\), and define
\[
A + B := \{a + b : a \in A \land b \in B\},
\]
\[
\lambda A := \{\lambda a : a \in A\}.
\]
If \(A\) and \(B\) are bounded, then
\[
\sup(A + B) = \sup A + \sup B,
\]
\[
\inf(A + B) = \inf A + \inf B,
\]
\[
\sup(\lambda A) = \begin{cases} 
\lambda \cdot \sup A & \text{for } \lambda \geq 0, \\
\lambda \cdot \inf A & \text{for } \lambda < 0,
\end{cases}
\]
\[
\inf(\lambda A) = \begin{cases} 
\lambda \cdot \inf A & \text{for } \lambda \geq 0, \\
\lambda \cdot \sup A & \text{for } \lambda < 0.
\end{cases}
\]

\[\Box\]

**Proof.** Exercise.

### 4.2 Important Subsets

**Remark 4.9.** We would like to recover the natural numbers \(\mathbb{N}\) as a subset of \(\mathbb{R}\). Indeed, if we start with 1 as the neutral element of multiplication and define \(2 := 1+1, 3 := 2+1, \ldots\), then \(\mathbb{N} := \{1, 2, \ldots\}\) is a subset of \(\mathbb{R}\), satisfying the Peano axioms P1, P2, P3 of Sec. 3.1. However, if one does actually construct \(\mathbb{R}\) according to the axioms of axiomatic set theory, then one starts by constructing \(\mathbb{N}\) first (basically as we did in Rem. 1.27 and Def. 1.26), constructing \(\mathbb{R}\) from \(\mathbb{N}\) in several steps (cf. [Phi16, Sec. D]). Depending on the construction used, the original set of natural numbers will typically not be the same set as the natural numbers as a subset of \(\mathbb{R}\). However, both sets will satisfy the Peano axioms and you will have a canonical bijection between the two sets. Which one you consider the “genuine” set of natural numbers depends on your personal taste.
and philosophy and is completely irrelevant. Any two models of \( \mathbb{N} \) will always produce equivalent results, since they must both satisfy the three Peano axioms.

We now introduce a zoo of important subsets of \( \mathbb{R} \) together with corresponding notation:

\[
\begin{align*}
\mathbb{N} &:= \{1, 2, 3, \ldots \} \quad \text{(natural numbers)}, \quad (4.10a) \\
\mathbb{N}_0 &:= \mathbb{N} \cup \{0\}, \quad (4.10b) \\
\mathbb{Z}^- &:= \{-n : n \in \mathbb{N}\} \quad \text{(negative integers)}, \quad (4.10c) \\
\mathbb{Z} &:= \mathbb{Z}^- \cup \mathbb{N}_0 \quad \text{(integers)}, \quad (4.10d) \\
\mathbb{Q}^+ &:= \{m/n : m, n \in \mathbb{N}\} \quad \text{(positive rational numbers)}, \quad (4.10e) \\
\mathbb{Q}^+_0 &:= \mathbb{Q}^+ \cup \{0\} \quad \text{(nonnegative rational numbers)}, \quad (4.10f) \\
\mathbb{Q}^- &:= \{-q : q \in \mathbb{Q}^+\} \quad \text{(negative rational numbers)}, \quad (4.10g) \\
\mathbb{Q}^-_0 &:= \mathbb{Q}^- \cup \{0\} \quad \text{(nonpositive rational numbers)}, \quad (4.10h) \\
\mathbb{Q} &:= \mathbb{Q}^+_0 \cup \mathbb{Q}^- \quad \text{(rational numbers)}, \quad (4.10i) \\
\mathbb{R}^+ &:= \{x \in \mathbb{R} : x > 0\} \quad \text{(positive real numbers)}, \quad (4.10j) \\
\mathbb{R}^+_0 &:= \{x \in \mathbb{R} : x \geq 0\} \quad \text{(nonnegative real numbers)}, \quad (4.10k) \\
\mathbb{R}^- &:= \{x \in \mathbb{R} : x < 0\} \quad \text{(negative real numbers)}, \quad (4.10l) \\
\mathbb{R}^-_0 &:= \{x \in \mathbb{R} : x \leq 0\} \quad \text{(nonpositive real numbers)}. \quad (4.10m)
\end{align*}
\]

For \( a, b \in \mathbb{R} \) with \( a \leq b \), one also defines the following intervals:

\[
\begin{align*}
[a, b] &:= \{x \in \mathbb{R} : a \leq x \leq b\} \quad \text{(bounded closed interval)}, \quad (4.11a) \\
]a, b[ &:= \{x \in \mathbb{R} : a < x < b\} \quad \text{(bounded open interval)}, \quad (4.11b) \\
]a, b] &:= \{x \in \mathbb{R} : a \leq x < b\} \quad \text{(bounded half-open interval)}, \quad (4.11c) \\
[a, b[ &:= \{x \in \mathbb{R} : a \leq x < b\} \quad \text{(bounded half-open interval)}, \quad (4.11d) \\
]a, \infty[ &:= \{x \in \mathbb{R} : x \leq b\} \quad \text{(unbounded closed interval)}, \quad (4.11e) \\
]a, \infty[ &:= \{x \in \mathbb{R} : x < b\} \quad \text{(unbounded open interval)}, \quad (4.11f) \\
[a, \infty[ &:= \{x \in \mathbb{R} : a \leq x\} \quad \text{(unbounded closed interval)}, \quad (4.11g) \\
]a, \infty[ &:= \{x \in \mathbb{R} : a < x\} \quad \text{(unbounded open interval)}. \quad (4.11h)
\end{align*}
\]

For \( a = b \), one says that the intervals defined by (4.11a) – (4.11d) are degenerate or trivial, where \([a, a] = \{a\}\), \([a, a[=]a, a[=]a, a]=\emptyset\) – it is sometimes convenient to have included the degenerate cases in the definition. It is sometimes also useful to abandon the restriction \( a \leq b \), to let \( c := \min\{a, b\} \), \( d := \max\{a, b\} \), and to define

\[
[a, b] := [c, d], \quad ]a, b[:=[c, d[, \quad ]a, b]:= [c, d] \setminus \{a\}, \quad [a, b[:=[c, d] \setminus \{b\}. \quad (4.11i)
\]

**Theorem 4.10** (Archimedean Property). Let \( \epsilon, x \) be real numbers. If \( \epsilon > 0 \) and \( x > 0 \), then there exists \( n \in \mathbb{N} \) such that \( n \epsilon > x \).
Proof. We conduct the proof by contradiction: Suppose \( x \) is an upper bound for the set \( A := \{ n \in \mathbb{N} : n \in \mathbb{N} \} \). Since the order \( \leq \) on \( \mathbb{R} \) is complete, according to (4.1), there exists \( s \in \mathbb{R} \) such that \( s = \sup A \). In particular, \( s - \epsilon \) is not an upper bound for \( A \), i.e. there exists \( n \in \mathbb{N} \) satisfying \( n \epsilon > s - \epsilon \). But then \( (n+1) \epsilon > s \) in contradiction to \( s = \sup A \).

This shows \( x \) is not an upper bound for \( A \), thereby establishing the case. ■

5 Complex Numbers

5.1 Definition and Basic Arithmetic

According to Th. 4.7(c), \( x^2 \geq 0 \) holds for every real number \( x \in \mathbb{R} \), i.e. the equation \( x^2 + 1 = 0 \) has no solution in \( \mathbb{R} \). This deficiency of the real numbers motivates the effort to try to extend the field of real numbers to a larger field \( \mathbb{C} \), the so-called complex numbers. The two requirements that \( \mathbb{C} \) is to be a field containing \( \mathbb{R} \) and that there is to be some complex number \( i \in \mathbb{C} \) satisfying \( i^2 = -1 \) already dictates the following laws of addition and multiplication for complex numbers \( z = x + iy \) and \( w = u + iv \) with \( x, y, u, v \in \mathbb{R} \):

\[
\begin{align*}
z + w &= x + iy + u + iv = x + u + i(y + v), \\
z w &= (x + iy)(u + iv) = xu - yv + i(xv + yu).
\end{align*}
\]

Moreover, if \( x + iy = u + iv \), then \( (x - u)^2 = -(v - y)^2 \), i.e. \( x - u = 0 = v - y \), implying \( x = u \) and \( y = v \). This suggests to try defining complex numbers as pairs of real numbers. Indeed, this works:

**Definition 5.1.** We define the set of complex numbers \( \mathbb{C} := \mathbb{R} \times \mathbb{R} \), where, keeping in mind (5.1), addition on \( \mathbb{C} \) is defined by

\[
+ : \mathbb{C} \times \mathbb{C} \to \mathbb{C}, \quad ((x, y), (u, v)) \mapsto (x, y) + (u, v) := (x + u, y + v),
\]

and multiplication on \( \mathbb{C} \) is defined by

\[
: \mathbb{C} \times \mathbb{C} \to \mathbb{C}, \quad ((x, y), (u, v)) \mapsto (x, y) \cdot (u, v) := (xu - yv, xv + yu).
\]

**Theorem 5.2.** (a) The set of complex numbers \( \mathbb{C} \) with addition and multiplication as defined in Def. 5.1 forms a field, where \((0, 0)\) and \((1, 0)\) are the neutral elements with respect to addition and multiplication, respectively,

\[
-z := (-x, -y)
\]

is the additive inverse to \( z = (x, y) \), whereas

\[
z^{-1} := \frac{1}{z} := \left( \frac{x}{x^2 + y^2}, \frac{-y}{x^2 + y^2} \right)
\]

is the multiplicative inverse to \( z = (x, y) \neq (0, 0) \).
(b) Defining subtraction and division in the usual way, for each \( z, w \in \mathbb{C} \), by \( w - z := w + (-z) \), and \( w/z := wz^{-1} \) for \( z \neq (0,0) \), respectively, all the rules stated in Th. 4.6 are valid in \( \mathbb{C} \).

(c) The map 
\[
\iota : \mathbb{R} \to \mathbb{C}, \quad \iota(x) := (x,0),
\]
(5.5)
is a monomorphism, i.e. it is injective and satisfies
\[
\forall x,y \in \mathbb{R} \quad \iota(x + y) = \iota(x) + \iota(y),
\]
(5.6a)
\[
\forall x,y \in \mathbb{R} \quad \iota(xy) = \iota(x) \cdot \iota(y).
\]
(5.6b)

It is customary to identify \( \mathbb{R} \) with \( \iota(\mathbb{R}) \), as it usually does not cause any confusion. One then just writes \( x \) instead of \( (x,0) \).

Proof. All computations required for (a) and (c) are straightforward and are left as an exercise; (b) is a consequence of (a), since Th. 4.6 and its proof are valid in every field.

Notation 5.3. The number \( i := (0,1) \) is called the imaginary unit (note that, indeed, \( i^2 = i \cdot i = (0,1) \cdot (0,1) = (0 \cdot 0 - 1 \cdot 1, 0 \cdot 1 + 1 \cdot 0) = (-1,0) = -1 \)). Using \( i \), one obtains the commonly used representation of a complex number \( z = (x,y) \in \mathbb{C} \):
\[
z = (x,y) = x \cdot (1,0) + y \cdot (0,1) = x + iy,
\]
(5.7)
where one calls Re \( z := x \) the real part of \( z \) and Im \( z := y \) the imaginary part of \( z \). Moreover, \( z \) is called purely imaginary if, and only if, Re \( z = 0 \) (as a consequence of this convention, one has the (harmless) pathology that 0 is both real and purely imaginary).

Remark 5.4. There does not exist a total order \( \leq \) on \( \mathbb{C} \) that makes \( \mathbb{C} \) into a totally ordered field (i.e. no total order on \( \mathbb{C} \) can be compatible with addition and multiplication in the sense of (4.6)): Indeed, if there were such a total order \( \leq \) on \( \mathbb{C} \), then all the rules of Th. 4.7 had to be valid with respect to that total order \( \leq \). In particular, \( 0 < 1^2 = 1 \) and \( 0 < i^2 = -1 \) had to be valid by Th. 4.7(c), and, then, \( 0 < 1 + (-1) = 0 \) had to be valid by Th. 4.7(f). However, \( 0 < 0 \) is false, showing that there is no total order on \( \mathbb{C} \) that satisfies (4.6). Caveat: Of course, there do exist total orders on \( \mathbb{C} \), just none compatible with addition and multiplication – for example, the lexicographic order on \( \mathbb{R} \times \mathbb{R} \) (defined as it was in (2.51) for \( \mathbb{N} \times \mathbb{N} \)) constitutes a total order on \( \mathbb{C} \).

Definition and Remark 5.5. Conjugation: For each complex number \( z = x + iy \), we define its complex conjugate or just conjugate to be the complex number \( \bar{z} := x - iy \). We then have the following rules that hold for each \( z = x + iy, w = u + iv \in \mathbb{C} \):

(a) \( \bar{z} + \bar{w} = x+u-iy-iv = \bar{z} + \bar{w} \) and \( \bar{z}w = xu-yv-(xv+yu)i = (x-iy)(u-iv) = \bar{z} \bar{w} \).

(b) \( z + \bar{z} = 2x = 2 \text{Re} z \) and \( z - \bar{z} = 2yi = 2i \text{Im} z \).

(c) \( z = \bar{z} \iff x + iy = x - iy \iff y = 0 \iff z \in \mathbb{R} \).
(d) \( z\bar{z} = (x + iy)(x - iy) = x^2 + y^2 \in \mathbb{R}_0^+ \).

**Notation 5.6.** Exponentiation with Integer Exponents: Define recursively for each \( z \in \mathbb{C} \) and each \( n \in \mathbb{N}_0^* \):

\[
z^0 := 1, \quad \forall \ n \in \mathbb{N}_0 \quad z^{n+1} := z \cdot z^n, \quad \text{and for } z \neq 0: \quad z^{-n} := (z^{-1})^n. \quad (5.8)
\]

**Theorem 5.7.** Exponentiation Rules: Let \( z, w \in \mathbb{C} \). For \( z, w \neq 0 \), the following rules hold for every \( m, n \in \mathbb{Z} \); otherwise they hold for each \( m, n \in \mathbb{N}_0^* \):

(a) \( z^{m+n} = z^m \cdot z^n \).

(b) \( z^n w^n = (zw)^n \).

(c) \( (z^m)^n = z^{mn} \).

**Proof.** (a): First, we prove the statement for each \( m \in \mathbb{N}_0 \) by induction: The base case \((m = 0)\) is \( z^n = z^n \), which is true. For the induction step, we compute

\[
z^{m+1+n} = z \cdot z^{m+n} \quad \overset{\text{ind. hyp.}}{=} \quad z \cdot z^m \cdot z^n \overset{(5.8)}{=} z^{m+1} \cdot z^n,
\]

completing the induction step. The above prove allows \( n < 0 \) for \( z \neq 0 \). Interchanging \( m \) and \( n \) covers the case \( m < 0 \) and \( n \geq 0 \). If \( m < 0 \) and \( n < 0 \), then

\[
z^{m+n} = z^{-(m+n)} = (z^{-1})^{-m-n} = (z^{-1})^{-m} \cdot (z^{-1})^{-n} = z^m \cdot z^n.
\]

(b): For \( n \in \mathbb{N}_0 \), the statement is proved by induction: The base case \((n = 0)\) is \( z^0 w^0 = 1 = (zw)^0 \), which is true. For the induction step, we compute

\[
z^{n+1} w^{n+1} = z \cdot z^n \cdot w \cdot w^n \quad \overset{\text{ind. hyp.}}{=} \quad zw \cdot (zw)^n \overset{(5.8)}{=} (zw)^{n+1},
\]

completing the induction step. For \( n < 0 \) and \( z \neq 0 \):

\[
z^n w^n = (5.8) \quad (z^{-1})^{-n} = (z^{-1} w^{-1})^{-n} \quad \overset{\text{Th. 4.6(c)}}{=} \quad ((zw)^{-1})^{-n} \overset{(5.8)}{=} (zw)^n.
\]

(c): First, we prove the statement for each \( n \in \mathbb{N}_0 \) by induction: The base case \((n = 0)\) is \( (z^m)^0 = 1 = z^0 \), which is true. For the induction step, we compute

\[
(z^m)^{n+1} = z^m \cdot (z^m)^n \quad \overset{\text{ind. hyp.}}{=} \quad z^m \cdot z^{mn} \overset{(a)}{=} z^{m+n} = z^{m(n+1)},
\]

completing the induction step. From (a), we also have \( (z^m)^{-1} = z^{-m} \) for \( z \neq 0 \). Thus, for \( n < 0 \) and \( z \neq 0 \):

\[
(z^m)^n = (5.8) \quad ((z^m)^{-1})^{-n} = (z^{-m})^{-n} = z^{-m(n)} = z^{mn},
\]

thereby completing the proof. ■
5.2 Sign and Absolute Value (Modulus)

We face a certain conundrum regarding the handling of square roots. The problem is that we will need the notion of a continuous function to prove the existence of a unique square root $\sqrt{x}$ for every nonnegative real number $x$ and, in consequence, we will have to wait until Section 7.2.5 below to carry out this proof. On the other hand, it is extremely desirable to present the theory of convergence simultaneously for real and for complex numbers, which requires the notion of the absolute value or modulus of a complex number, to be defined in Def. 5.9(b) below as the square root of a nonnegative real number.

Faced with this difficulty, we will introduce the notion of square root now, assuming the existence, until we can add the proof in Section 7.2.5. Some students might be worried that this might lead to a circular argument, where our later proof of the existence of square roots would somehow make use of our previous assumption of that existence. Of course, we will be careful not to make such a circular (and, thereby, logically invalid) argument. The point is that for real numbers the notion of absolute value does in no way depend on the notion of a square root (see Lem. 5.10 below).

**Definition and Remark 5.8.** We define a nonnegative real number $y \in \mathbb{R}_0^+$ to be the square root of the nonnegative real number $x \in \mathbb{R}_0^+$ if, and only if, $y^2 = x$. If $y$ is the square root of $x$, then one uses the notation $\sqrt{x} := y$. We will see in Rem. and Def. 7.61 that every $x \in \mathbb{R}_0^+$ has a unique square root and that the function $f : \mathbb{R}_0^+ \to \mathbb{R}_0^+$, $f(x) := \sqrt{x}$, is strictly increasing (in particular, injective).

**Definition 5.9.** (a) The sign function is defined by

$$\text{sgn} : \mathbb{R} \to \mathbb{R}, \quad \text{sgn}(x) := \begin{cases} 1 & \text{for } x > 0, \\ 0 & \text{for } x = 0, \\ -1 & \text{for } x < 0. \end{cases} \quad (5.9)$$

It is emphasized that the sign function is only defined for real numbers (cf. Rem. 5.4)!

(b) The absolute value or modulus function is defined by

$$\text{abs} : \mathbb{C} \to \mathbb{R}_0^+, \quad z = x + iy \mapsto |z| := \sqrt{z \bar{z}} = \sqrt{x^2 + y^2}, \quad (5.10)$$

where the term absolute value is often preferred for real numbers $z \in \mathbb{R}$ and the term modulus is often preferred if one also considers complex numbers $z \notin \mathbb{R}$.

**Lemma 5.10.** For each $x \in \mathbb{R}$, one has

$$|x| = x \cdot \text{sgn}(x) = \begin{cases} x & \text{for } x \geq 0, \\ -x & \text{for } x < 0. \end{cases} \quad (5.11)$$

**Proof.** One has

$$|x| = \sqrt{x^2} = \begin{cases} x & \text{for } x \geq 0, \\ -x & \text{for } x < 0. \end{cases} \quad (5.12)$$
as claimed.

**Theorem 5.11.** The following rules hold for each \( z, w \in \mathbb{C} \):

(a) \( z \neq 0 \Rightarrow |z| > 0 \).
(b) \( ||z|| = |z| \).
(c) \( |z| = |\bar{z}| \).
(d) \( \max\{|\Re z|, |\Im z|\} \leq |z| \leq |\Re z| + |\Im z| \).
(e) \( |zw| = |z||w| \).
(f) For \( w \neq 0 \), one has \( |z| = |w| \).
(g) Triangle Inequality:
\[
|z + w| \leq |z| + |w|.
\] (5.13)
(h) Inverse Triangle Inequality:
\[
||z| - |w|| \leq |z - w|.
\] (5.14)

**Proof.** We carry out the proofs for \( z, w \in \mathbb{C} \). However, for \( z, w \in \mathbb{R} \), everything can easily be shown directly from (5.11), without making use of square roots.

Let \( z = x + iy \) with \( x, y \in \mathbb{R} \).

(a): If \( z \neq 0 \), then \( x \neq 0 \) or \( y \neq 0 \), i.e. \( x^2 > 0 \) or \( y^2 > 0 \) by Th. 4.7(c), implying \( x^2 + y^2 > 0 \) by Th. 4.7(f), i.e. \( |z| = \sqrt{x^2 + y^2} > 0 \).

(b): Since \( a := |z| \in \mathbb{R}^+_0 \), we have \( |a| = \sqrt{a^2} = a = |z| \).

(c): Since \( \bar{z} = x - iy \), we have \( |\bar{z}| = \sqrt{x^2 + (-y)^2} = \sqrt{x^2 + y^2} = |z| \).

(d): It is \( x = \Re z \), \( y = \Im z \). Let \( a := \max\{|x|, |y|\} \). As remarked in Def. and Rem. 5.8, the square root function is increasing and, thus, taking square roots in the chain of inequalities \( a^2 \leq x^2 + y^2 \leq (|x| + |y|)^2 \) implies \( a \leq |z| \leq |x| + |y| \) as claimed.

(e): As remarked in Def. and Rem. 5.8, the square root function is injective, and, thus, (e) follows from
\[
|zw|^2 = z\bar{w} w = |z|^2 |w|^2.
\]

(f): Let \( w = u + iv \) with \( u, v \in \mathbb{R} \). We first consider the special case \( z = 1 \). Applying the formula (5.4b) for the inverse to \( w \), one obtains
\[
|w^{-1}|^2 = \frac{u^2}{(u^2 + v^2)^2} + \frac{v^2}{(u^2 + v^2)^2} = \frac{1}{u^2 + v^2} = (|w|^{-1})^2,
\]
i.e. \( |w^{-1}| = |w|^{-1} \). Now (f) follows from (e): \( |z_w| = |zw| = |z||w| = |z||w|^{-1} = |z|^{-1} = |z|^{-1} = |z|^{-1} = |z|^{-1} \).
(g) follows from

\[ |z + w|^2 = (z + w)(\bar{z} + \bar{w}) = z\bar{z} + w\bar{z} + zw + w\bar{w} \]

Def. and Rem. 5.5(b)

\[ = |z|^2 + 2\text{Re}(z\bar{w}) + |w|^2 \]

\[ \leq |z|^2 + 2|z\bar{w}| + |w|^2 = (|z| + |w|)^2, \]

once again using that the square root function is increasing.

(h): Using (g), we obtain

\[ |z| = |z - w + w| \leq |z - w| + |w| \quad \Rightarrow \quad |z| - |w| \leq |z - w|, \]

\[ |w| = |w - z + z| \leq |z - w| + |z| \quad \Rightarrow \quad -(|z| - |w|) \leq |z - w|, \]

implying \(||z| - |w|| \leq |z - w|\) by (5.11) (notice \(|z| - |w| \in \mathbb{R}|\).

\[ \blacksquare \]

**Remark 5.12.** Each complex number \((x, y) = x + iy\) can be visualized as a point in the so-called complex plane, where the horizontal \(x\)-axis represents real numbers and the vertical \(y\)-axis represents purely imaginary numbers. Then the addition of complex numbers is precisely the vector addition of 2-dimensional vectors in the complex plane, and conjugation is represented by reflection through the \(x\)-axis. Moreover, the modulus \(|z|\) of a complex number is precisely its distance from the origin \((0, 0)\), and \(|z - w|\) is the distance between the points \(z = (x, y)\) and \(w = (u, v)\) in the plane. Complex multiplication can also be interpreted geometrically in the plane: If \(\phi\) denotes the angle that the vector representing \(z = (x, y)\) forms with the \(x\)-axis, and, likewise, \(\psi\) denotes the angle that the vector representing \(w = (u, v)\) forms with the \(x\)-axis, then \(zw\) is the vector of length \(|zw|\) that forms the angle \(\phi + \psi\) with the \(x\)-axis (we will better understand this geometrical interpretation of complex multiplication later (see Def. and Rem. 8.29), when writing complex numbers in the polar form \(z = x + iy = |z|\exp(i\phi)\), making use of the exponential function \(\exp\)).

### 5.3 Sums and Products

Here we compile some important rules involving sums and products of complex numbers (the exceptions are the estimates in Th. 5.13(d),(e) below, which actually require real numbers):

**Theorem 5.13.** (a) For each \(n \in \mathbb{N}\) and each \(\lambda, \mu, z_j, w_j \in \mathbb{C}, j \in \{1, \ldots, n\}\):

\[ \sum_{j=1}^{n} (\lambda z_j + \mu w_j) = \lambda \sum_{j=1}^{n} z_j + \mu \sum_{j=1}^{n} w_j. \]

(b) For each \(n \in \mathbb{N}_0\) and each \(z \in \mathbb{C}\):

\[ (1 - z)(1 + z + z^2 + \cdots + z^n) = (1 - z) \sum_{j=0}^{n} z^j = 1 - z^{n+1}. \]
(c) For each \( n \in \mathbb{N}_0 \) and each \( z, w \in \mathbb{C} \):

\[
w^{n+1} - z^{n+1} = (w - z) \sum_{j=0}^{n} z^j w^{n-j} = (w - z)(w^n + zw^{n-1} + \cdots + z^{n-1}w + z^n).
\]

(d) For each \( n \in \mathbb{N} \) and each \( x_j, y_j \in \mathbb{R}, j \in \{1, \ldots, n\} \):

\[
\left( \forall j \in \{1, \ldots, n\} \ x_j \leq y_j \right) \Rightarrow \sum_{j=1}^{n} x_j \leq \sum_{j=1}^{n} y_j,
\]

where equality can only hold if \( x_j = y_j \) for each \( j \in \{1, \ldots, n\} \).

(e) For each \( n \in \mathbb{N} \) and each \( x_j, y_j \in \mathbb{R}, j \in \{1, \ldots, n\} \):

\[
\left( \forall j \in \{1, \ldots, n\} \ 0 < x_j \leq y_j \right) \Rightarrow \prod_{j=1}^{n} x_j \leq \prod_{j=1}^{n} y_j,
\]

where equality can only hold if \( x_j = y_j \) for each \( j \in \{1, \ldots, n\} \).

(f) Triangle Inequality: For each \( n \in \mathbb{N} \) and each \( z_j \in \mathbb{C}, j \in \{1, \ldots, n\} \):

\[
\left| \sum_{j=1}^{n} z_j \right| \leq \sum_{j=1}^{n} |z_j|.
\]

Proof. In each case, the proof can be conducted by an easy induction. We carry out (c) and leave the other cases as exercises. For (c), the base case \( (n = 0) \) is provided by the true statement \( w^{0+1} - z^{0+1} = w - z = (w - z)z^0 w^0 - z^0 w^0 \). For the induction step, one computes

\[
(w - z) \sum_{j=0}^{n+1} z^j w^{n+1-j} = (w - z) \left( z^{n+1}w^0 + \sum_{j=0}^{n} z^j w^{n+1-j} \right)
\]

\[
= (w - z)z^{n+1} + (w - z) w \sum_{j=0}^{n} z^j w^{n-j}
\]

\[
\overset{\text{ind. hyp.}}{=} (w - z)z^{n+1} + w(w^{n+1} - z^{n+1}) = w^{n+2} - z^{n+2},
\]

completing the induction.

\[ \blacksquare \]

5.4 Binomial Coefficients and Binomial Theorem

The goal in this section is to expand \((z + w)^n\) into a sum. This sum involves the so-called binomial coefficients \( \binom{n}{k} \), which are also useful in other contexts. To obtain an idea for what to expect, let us compute the cases \( n = 0, 1, 2, 3 \): \((z + w)^0 = 1\), \((z + w)^1 = z + w\),
\[ (z + w)^2 = z^2 + 2zw + w^2, \ (z + w)^3 = z^3 + 3z^2w + 3zw^2 + w^3. \] One finds that the coefficients form what is known as Pascal’s triangle, which we write for \( n = 0, \ldots, 5 \):

\[
\begin{array}{cccccc}
\text{n = 0:} & 1 \\
\text{n = 1:} & 1 & 1 \\
\text{n = 2:} & 1 & 2 & 1 \\
\text{n = 3:} & 1 & 3 & 3 & 1 \\
\text{n = 4:} & 1 & 4 & 6 & 4 & 1 \\
\text{n = 5:} & 1 & 5 & 10 & 10 & 5 & 1
\end{array}
\] (5.15)

The entries of the \( n \)th row of Pascal’s triangle are denoted by \( \binom{n}{0}, \ldots, \binom{n}{n} \). One also observes that one obtains each entry of the \((n+1)\)st row, except the first and last entry, by adding the corresponding entries in row \( n \) to the left and to the right of the considered entry in row \( n + 1 \). The first and last entry of each row are always set to 1. This can be summarized as

\[
\forall n \in \mathbb{N}_0 \ \left( \binom{n}{0} = \binom{n}{n} = 1, \quad \binom{n+1}{k} = \binom{n}{k-1} + \binom{n}{k} \text{ for } k \in \{1, \ldots, n\} \right). \tag{5.16}
\]

The following Def. 5.14 provides a different and more general definition of binomial coefficients. We will then prove in Prop. 5.15 that the binomial coefficients as defined in Def. 5.14 do, indeed, satisfy (5.16).

**Definition 5.14.** For each \( \alpha \in \mathbb{C} \) and each \( k \in \mathbb{N}_0 \), we define the *binomial coefficient*

\[
\binom{\alpha}{0} := 1, \quad \binom{\alpha}{k} := \prod_{j=1}^{k} \frac{\alpha + 1 - j}{j} = \frac{\alpha(\alpha - 1) \cdots (\alpha - k + 1)}{1 \cdot 2 \cdots k} \text{ for } k \in \mathbb{N}. \tag{5.17}
\]

**Proposition 5.15.** (a) For each \( \alpha \in \mathbb{C} \) and each \( k \in \mathbb{N} \):

\[
\binom{\alpha}{0} = 1, \quad \binom{\alpha + 1}{k} = \binom{\alpha}{k-1} + \binom{\alpha}{k}. \tag{5.18}
\]

(b) For each \( n \in \mathbb{N}_0 \):

\[
\binom{n}{n} = 1. \tag{5.19}
\]

The above statements include (5.16) as a special case.

**Proof.** (a): The first identity is part of the definition in (5.17). For the second identity, we first observe, for each \( k \in \mathbb{N} \),

\[
\binom{\alpha}{k} = \prod_{j=1}^{k} \frac{\alpha + 1 - j}{j} = \frac{\alpha + 1 - k}{k} \prod_{j=1}^{k-1} \frac{\alpha + 1 - j}{j} = \binom{\alpha}{k-1} \frac{\alpha + 1 - k}{k}. \tag{5.20}
\]
which implies
\[
\binom{\alpha}{k-1} + \binom{\alpha}{k} = \binom{\alpha}{k-1} \left( 1 + \frac{\alpha + 1 - k}{k} \right) = \binom{\alpha}{k-1} \frac{\alpha + 1}{k} \\
= \frac{\alpha + 1}{k} \prod_{j=1}^{k-1} \frac{\alpha + 1 - j}{j} = \prod_{j=1}^{k} \frac{\alpha + 2 - j}{j} = \binom{\alpha + 1}{k}. \tag{5.21}
\]

(b): \( \binom{0}{0} = 1 \) according to (5.17). For \( n \in \mathbb{N} \), (5.19) is proved by induction. The base case \( (n = 1) \) is provided by the true statement \( \binom{1}{1} = 1 \). For the induction step, one computes
\[
\binom{n+1}{n+1} = \prod_{j=1}^{n+1} \frac{n + 1 + 1 - j}{j} = \frac{n+1}{n+1} \prod_{j=1}^{n+1} \frac{n+1 - j}{j} = \binom{n}{n} \text{ ind. hyp.} = 1, \tag{5.22}
\]
which completes the induction.

**Theorem 5.16** (Binomial Theorem). For each \( z, w \in \mathbb{C} \) and each \( n \in \mathbb{N}_0 \), the following formula holds:
\[
(z + w)^n = \sum_{k=0}^{n} \binom{n}{k} z^{n-k} w^k = z^n + \binom{n}{1} z^{n-1} w + \cdots + \binom{n}{n-1} z w^{n-1} + w^n. \tag{5.23}
\]

**Proof.** The proof is conducted via induction on \( n \). The base case \( (n = 0) \) is provided by the correct statement \( (z + w)^0 = \binom{0}{0} z^0 w^0 = 1 \). For the induction step, we first observe
\[
(z + w)^{n+1} = (z + w) (z + w)^n = z (z + w)^n + w (z + w)^n. \tag{5.24}
\]
Using the induction hypothesis, we now further manipulate the two terms on the right-hand side of (5.24):
\[
\begin{align*}
z (z + w)^n & \overset{\text{ind. hyp.}}{=} z \sum_{k=0}^{n} \binom{n}{k} z^{n-k} w^k = \sum_{k=0}^{n} \binom{n}{k} z^{n+1-k} w^k \\
\overset{(n+1)_{\text{ind. hyp.}}}{=} &= \sum_{k=0}^{n+1} \binom{n}{k} z^{n+1-k} w^k, \tag{5.25}
\end{align*}
\]
\[
\begin{align*}
w (z + w)^n & \overset{\text{ind. hyp.}}{=} w \sum_{k=0}^{n} \binom{n}{k} z^{n-k} w^k = \sum_{k=0}^{n} \binom{n}{k} z^{n-k} w^{k+1} \\
&= \sum_{k=1}^{n+1} \binom{n}{k-1} z^{n+1-k} w^k. \tag{5.26}
\end{align*}
\]
Plugging (5.25) and (5.26) into (5.24) yields

\[(z + w)^{n+1} = \binom{n}{0} z^{n+1} w^0 + \sum_{k=1}^{n+1} \left( \binom{n}{k} + \binom{n}{k-1} \right) z^{n+1-k} w^k\]

Prop. 5.15

\[(n + 1)^{n+1} = \binom{n+1}{0} z^{n+1} w^0 + \sum_{k=1}^{n+1} \binom{n+1}{k} z^{n+1-k} w^k\]

\[= \sum_{k=0}^{n+1} \binom{n+1}{k} z^{n+1-k} w^k,\]  \hspace{1cm} (5.27)

completing the induction. ■

The binomial theorem can now be used to infer a few more rules that hold for the binomial coefficients:

**Corollary 5.17.** One has the following identities:

\[\forall n \in \mathbb{N}_0 \quad \sum_{k=0}^{n} \binom{n}{k} = \binom{n}{0} + \binom{n}{1} + \cdots + \binom{n}{n} = 2^n,\]  \hspace{1cm} (5.28a)

\[\forall n \in \mathbb{N} \quad \sum_{k=0}^{n} \binom{n}{k} \left(-1\right)^k = \binom{n}{0} - \binom{n}{1} + \binom{n}{2} - \cdots + (-1)^n \binom{n}{n} = 0.\]  \hspace{1cm} (5.28b)

**Proof.** (5.28a) is just (5.23) with \(z = w = 1\); (5.28b) is just (5.23) with \(z = 1\) and \(w = -1\). ■

The formulas provided by the following proposition are also sometimes useful.

**Proposition 5.18.** (a) For each \(\alpha \in \mathbb{C}\) and each \(k \in \mathbb{N}_0\):

\[\sum_{j=0}^{k} \binom{\alpha + j}{j} = \binom{\alpha}{0} + \binom{\alpha + 1}{1} + \cdots + \binom{\alpha + k}{k} = \binom{\alpha + k + 1}{k}.\]  \hspace{1cm} (5.29)

(b) For each \(n, k \in \mathbb{N}_0\) with \(k \leq n\):

\[\binom{n}{k} = \frac{n!}{k!(n-k)!}.\]  \hspace{1cm} (5.30)

Moreover, for \(n \geq 1\), one has \(\binom{n}{k} = \#\mathcal{P}_k(\{1, \ldots, n\})\), where

\[\mathcal{P}_k(A) := \{B \in \mathcal{P}(A) : \#B = k\}.\]  \hspace{1cm} (5.31)

\[\text{denotes the set of all subsets of a set } A \text{ that have precisely } k \text{ elements.}\]

(c) For each \(n, k \in \mathbb{N}_0\):

\[\sum_{j=0}^{k} \binom{n + j}{n} = \binom{n}{n} + \binom{n + 1}{n} + \cdots + \binom{n + k}{n} = \binom{n + k + 1}{n + 1}.\]  \hspace{1cm} (5.32)
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Proof. The induction proofs of (a) and (b) are left as exercises. For (c), one computes

\[
\sum_{j=0}^{k} \binom{n+j}{n} \overset{(5.30)}{=} \sum_{j=0}^{k} \frac{(n+j)!}{n!(n+j-n)!} \overset{(5.30)}{=} \sum_{j=0}^{k} \binom{n+j}{j} \overset{(5.29)}{=} \binom{n+k+1}{k} \overset{(5.30)}{=} \frac{(n+k+1)!}{k!(n+1)!} = \binom{n+k+1}{n+1},
\]

thereby establishing the case. ■

6 Polynomials

6.1 Arithmetic of \( \mathbb{K} \)-Valued Functions

Notation 6.1. We will write \( \mathbb{K} \) in situations, where we allow \( \mathbb{K} \) to be \( \mathbb{R} \) or \( \mathbb{C} \).

Notation 6.2. If \( A \) is any nonempty set, then one can add and multiply arbitrary functions \( f, g : A \to \mathbb{K} \), and one can define several further operations to create new functions from \( f \) and \( g \):

\[
\begin{align*}
(f + g) : A &\to \mathbb{K}, \quad (f + g)(x) := f(x) + g(x), \quad (6.1a) \\
(\lambda f) : A &\to \mathbb{K}, \quad (\lambda f)(x) := \lambda f(x) \text{ for each } \lambda \in \mathbb{K}, \quad (6.1b) \\
(fg) : A &\to \mathbb{K}, \quad (fg)(x) := f(x)g(x), \quad (6.1c) \\
(f/g) : A &\to \mathbb{K}, \quad (f/g)(x) := f(x)/g(x) \text{ (assuming } g(x) \neq 0), \quad (6.1d) \\
\text{Re} f : A &\to \mathbb{R}, \quad (\text{Re} f)(x) := \text{Re}(f(x)), \quad (6.1e) \\
\text{Im} f : A &\to \mathbb{R}, \quad (\text{Im} f)(x) := \text{Im}(f(x)). \quad (6.1f)
\end{align*}
\]

For \( \mathbb{K} = \mathbb{R} \), we further define

\[
\begin{align*}
\max(f, g) : A &\to \mathbb{R}, \quad \max(f, g)(x) := \max\{f(x), g(x)\}, \quad (6.1g) \\
\min(f, g) : A &\to \mathbb{R}, \quad \min(f, g)(x) := \min\{f(x), g(x)\}, \quad (6.1h) \\
f^+ : A &\to \mathbb{R}, \quad f^+ := \max(f, 0), \quad (6.1i) \\
f^- : A &\to \mathbb{R}, \quad f^- := \max(-f, 0). \quad (6.1j)
\end{align*}
\]

Finally, once again also allowing \( \mathbb{K} = \mathbb{C} \),

\[
|f| : A \to \mathbb{R}, \quad |f|(x) := |f(x)|. \quad (6.1k)
\]

One calls \( f^+ \) and \( f^- \) the positive part and the negative part of \( f \), respectively. For \( \mathbb{R} \)-valued functions \( f \), we have

\[
|f| = f^+ + f^- . \quad (6.1l)
\]
6.2 Polynomials

Definition 6.3. Let \( n \in \mathbb{N}_0 \). Each function from \( K \) into \( K \), \( x \mapsto x^n \), is called a monomial. A function \( P \) from \( K \) into \( K \) is called a polynomial if, and only if, it is a linear combination of monomials, i.e. if, and only if \( P \) has the form

\[
P : K \rightarrow K, \quad P(x) = \sum_{j=0}^{n} a_j x^j = a_0 + a_1 x + \cdots + a_n x^n, \quad a_j \in K.
\]  

(6.2)

The \( a_j \) are called the coefficients of \( P \). The largest number \( d \leq n \) such that \( a_d \neq 0 \) is called the degree of \( P \), denoted \( \deg(P) \). If all coefficients are 0, then \( P \) is called the zero polynomial; the degree of the zero polynomial is defined as \(-1\) (in Th. 6.6(b) below, we will see that each polynomial of degree \( n \in \mathbb{N}_0 \) is uniquely determined by its coefficients \( a_0, \ldots, a_n \) and vice versa).

Polynomials of degree \( \leq 0 \) are constant. Polynomials of degree \( \leq 1 \) have the form \( P(x) = a + bx \) and are called affine functions (often they are also called linear functions, even though this is not really correct for \( a \neq 0 \), since every function \( P \) that is linear (in the sense of linear algebra) must satisfy \( P(0) = 0 \)). Polynomials of degree \( \leq 2 \) have the form \( P(x) = a + bx + cx^2 \) and are called quadratic functions.

Each \( \xi \in K \) such that \( P(\xi) = 0 \) is called a zero or a root of \( P \).

A rational function is a quotient \( P/Q \) of two polynomials \( P \) and \( Q \).

Remark 6.4. Let \( \lambda \in K \) and let \( P, Q \) be polynomials. Then \( \lambda P, P+Q, \) and \( PQ \) defined according to Not. 6.2 are polynomials as well. More precisely, if \( \lambda = 0 \) or \( P \equiv 0 \), then \( \lambda P = 0 \); if \( P \equiv 0 \), then \( P + Q = Q \); if \( Q \equiv 0 \), then \( P + Q = P \); if \( P \equiv 0 \) or \( Q \equiv 0 \), then \( PQ = 0 \). If \( \lambda \neq 0 \) and

\[
P(x) = \sum_{j=0}^{n} a_j x^j, \quad Q(x) = \sum_{j=0}^{m} b_j x^j,
\]  

with \( \deg(P) = n \geq 0, \quad \deg(Q) = m \geq 0, \quad n \geq m \geq 0, \) then, defining \( b_j := 0 \) for each \( j \in \{m+1, \ldots, n\} \) in case \( n > m \),

\[
(\lambda P)(x) = \sum_{j=0}^{n} (\lambda a_j) x^j, \quad \deg(\lambda P) = n, \tag{6.4a}
\]

\[
(P + Q)(x) = \sum_{j=0}^{n} (a_j + b_j) x^j, \quad \deg(P + Q) \leq n = \max\{m, n\}, \tag{6.4b}
\]

\[
(PQ)(x) = \sum_{j=0}^{m+n} c_j x^j, \quad \deg(PQ) = m + n, \tag{6.4c}
\]

where, setting \( a_k := 0 \) for each \( k \in \{n+1, \ldots, m+n\} \) and \( b_k := 0 \) for each \( k \in \{m+1, \ldots, m+n\} \),

\[
\forall \ j \in \{0, \ldots, n+m\} \quad c_j = \sum_{k=0}^{j} a_k b_{j-k}. \tag{6.4d}
\]
Formula (6.4c) can be proved by induction on \( m = \deg(Q) \in \mathbb{N}_0 \) as follows: For \( m = 0 \), we compute
\[
(PQ)(x) = b_0 \sum_{j=0}^{n} a_j x^j = \sum_{j=0}^{n+0} b_0 a_j x^j,
\]
i.e. \( c_j = b_0 a_j = \sum_{k=0}^{j} a_k b_{j-k} \), which establishes the base case, remembering \( b_{j-k} = 0 \) for \( j > k \). For the induction step, we compute, for \( \deg(Q) = m + 1 \),
\[
(PQ)(x) = \sum_{j=0}^{n} a_j x^j \sum_{\alpha=0}^{m+1} b_\alpha x^\alpha = \sum_{j=0}^{n} a_j x^j \left( b_{m+1} x^{m+1} + \sum_{\alpha=0}^{m} b_\alpha x^\alpha \right)
\]
\[
= \sum_{j=0}^{m+n+1} a_j b_{m+1} x^j + \sum_{j=0}^{m+n} \left( \sum_{k=0}^{j} a_k b_{j-k} \right) x^j
\]
\[
= \sum_{j=0}^{m+n+1} \left( \sum_{k=0}^{j} a_k b_{j-k} \right) x^j,
\]
which completes the induction step. There is a notational issue in the second and third line in of the above computation, since, in both lines, the \( b_{m+1} \) in the first sum is the actual \( b_{m+1} \) from \( Q \), but \( b_{m+1} = 0 \) in the second sum in both lines, which is due to the induction hypothesis being applied for \( m < m + 1 \). This is actually used when combining both sums in the last step, computing, for \( m + 1 \leq j \leq m + n \): \( a_{j-m-1} b_{m+1} x^j + a_{j-m-1} \cdot 0 \cdot x^j = a_{j-m-1} b_{m+1} x^j \). For \( j = m + n + 1 \), one has \( \sum_{k=0}^{m+n+1} a_k b_{m+n+1-k} = a_n b_{m+1} \), since \( b_{m+n+1-k} = 0 \) for \( n > k \) and \( a_k = 0 \) for \( k > n \).

Finally, \( \deg(PQ) = m + n \) follows from \( c_{m+n} = a_m b_n \neq 0 \).

**Theorem 6.5. (a)** For each polynomial \( P \) given in the form of (6.3) and each \( \xi \in \mathbb{K} \), we have the identity
\[
P(x) = \sum_{j=0}^{n} b_j (x - \xi)^j,
\]
where
\[
\forall j \in \{0, \ldots, n\} \quad b_j = \sum_{k=j}^{n} a_k \binom{k}{j} \xi^{k-j}, \quad \text{in particular} \quad b_0 = P(\xi), \quad b_n = a_n.
\]

**Theorem 6.5. (b)** If \( P \) is a polynomial with \( n := \deg(P) \geq 1 \), then, for each \( \xi \in \mathbb{K} \), there exists a polynomial \( Q \) with \( \deg(Q) = n - 1 \) such that
\[
P(x) = P(\xi) + (x - \xi) Q(x).
\]
In particular, if \( \xi \) is a zero of \( P \), then \( P(x) = (x - \xi) Q(x) \).
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Proof. (a): For $\xi = 0$, there is nothing to prove. For $\xi \neq 0$, defining the auxiliary variable $\eta := x - \xi$, we obtain $x = \xi + \eta$ and

$$P(x) = \sum_{k=0}^{n} a_k (\xi + \eta)^k = \sum_{j=0}^{n} \sum_{j=0}^{k} a_k \binom{k}{j} \xi^{k-j} \eta^j = \sum_{j=0}^{n} \sum_{k=j}^{n} a_k \binom{k}{j} \xi^{k-j} \eta^j,$$

with (6.5).

(b): According to (a), we have

$$P(x) = P(\xi) + (x - \xi) Q(x), \quad \text{with} \quad Q(x) = \sum_{j=1}^{n} b_j (x - \xi)^{j-1} = \sum_{j=0}^{n-1} b_{j+1} (x - \xi)^j,$$

proving (b).

Theorem 6.6. (a) If $P$ is a polynomial with $n := \deg(P) \geq 0$, then $P$ has at most $n$ zeros.

(b) Let $P, Q$ be polynomials as in (6.3) with $n = m$, $\deg(P) \leq n$, and $\deg(Q) \leq n$. If $P(x_j) = Q(x_j)$ at $n + 1$ distinct points $x_0, x_1, \ldots, x_n \in \mathbb{K}$, then $a_j = b_j$ for each $j \in \{0, \ldots, n\}$.

Consequence 1: If $P, Q$ with degree $\leq n$ agree at $n + 1$ distinct points, then $P = Q$.

Consequence 2: If we know $P = Q$, then they agree everywhere, in particular at $\max\{\deg(P), \deg(Q)\} + 1$ distinct points, which implies they have the same coefficients.

Proof. (a): For $n = 0$, $P$ is constant, but not the zero polynomial, i.e. $P \equiv a_0 \neq 0$ with no zeros as claimed. For $n \in \mathbb{N}$, the proof is conducted by induction. The base case ($n = 1$) is provided by the observation that $\deg(P) = 1$ implies $P$ is the affine function with $P(x) = a_0 + a_1 x$, $a_1 \neq 0$, i.e. $P$ has precisely one zero at $\xi = -a_0/a_1$. For the induction step, assume $\deg(P) = n + 1$. If $P$ has no zeros, then the assertion of (a) holds true. Otherwise, $P$ has at least one zero $\xi \in \mathbb{K}$, and, according to Th. 6.5(b), there exists a polynomial $Q$ such that $\deg(Q) = n$ and

$$P(x) = (x - \xi) Q(x).$$

From the induction hypothesis, we gather that $Q$ has at most $n$ zeros, i.e. (6.10) implies $P$ has at most $n + 1$ zeros, which completes the induction.

(b): If $P(x_j) = Q(x_j)$ at $n + 1$ distinct points $x_j$, then each of these points is a zero of $P - Q$. Thus $P - Q$ is a polynomial of degree $\leq n$ with at least $n + 1$ zeros. Then (a) implies $\deg(P - Q) = -1$, i.e. $P - Q$ is the zero polynomial, i.e. $a_j - b_j = 0$ for each $j \in \{0, \ldots, n\}$.
Remark 6.7. Let $P$ be a polynomial with $n := \deg(P) \geq 0$. According to Th. 6.6(a), $P$ has at most $n$ zeros. Using Th. 6.5(b) for an induction shows there exists $k \in \{0, \ldots, n\}$ and a polynomial $Q$ of degree $n - k$ such that

$$P(x) = Q(x) \prod_{j=1}^{k} (x - \xi_j) = (x - \xi_1)(x - \xi_2) \cdots (x - \xi_k)Q(x),$$

(6.11a)

where $Q$ does not have any zeros in $\mathbb{K}$ and $\{\xi_1, \ldots, \xi_k\} = \{\xi \in \mathbb{K} : P(\xi) = 0\}$ is the set of zeros of $P$. It can of course happen that $P$ does not have any zeros and $P = Q$ (no $\xi_j$ exist). It can also occur that some of the $\xi_j$ in (6.11a) are identical. Thus, we can rewrite (6.11a) as

$$P(x) = Q(x) \prod_{j=1}^{l} (x - \lambda_j)^{m_j} = (x - \lambda_1)^{m_1}(x - \lambda_2)^{m_2} \cdots (x - \lambda_l)^{m_l}Q(x),$$

(6.11b)

where $\lambda_1, \ldots, \lambda_l$, $l \in \{0, \ldots, k\}$, are the distinct zeros of $P$, and $m_j \in \mathbb{N}$ with $\sum_{j=1}^{l} m_j = k$. Then $m_j$ is called the multiplicity of the zero $\lambda_j$ of $P$.

7 Limits and Convergence in the Real and Complex Numbers

7.1 Sequences

Recall from Def. 2.14(b) that a sequence in $\mathbb{K}$ is a function $f : \mathbb{N} \to \mathbb{K}$, in this context usually denoted as $f = (z_n)_{n \in \mathbb{N}}$ or $(z_1, z_2, \ldots)$ with $z_n := f(n)$. Sometimes the sequence also has the form $(z_n)_{n \in I}$, where $I \neq \emptyset$ is a countable index set (e.g. $I = \mathbb{N}_0$) different from $\mathbb{N}$ (in the context of convergence (see the following Def. 7.1), $I$ must be $\mathbb{N}$ or it must have the same cardinality as $\mathbb{N}$, i.e. finite $I$ are not permissible).

Definition 7.1. The sequence $(z_n)_{n \in \mathbb{N}}$ in $\mathbb{K}$ is said to be convergent with limit $z \in \mathbb{K}$ if, and only if, for each $\epsilon > 0$, there exists an index $N \in \mathbb{N}$ such that $|z_n - z| < \epsilon$ for every index $n > N$. The notation for $(z_n)_{n \in \mathbb{N}}$ converging to $z$ is $\lim_{n \to \infty} z_n = z$ or $z_n \to z$ for $n \to \infty$. Thus, by definition,

$$\lim_{n \to \infty} z_n = z \iff \forall \epsilon \in \mathbb{R}^+ \exists N \in \mathbb{N} \forall n > N \ |z_n - z| < \epsilon.$$  

(7.1)

The sequence $(z_n)_{n \in \mathbb{N}}$ in $\mathbb{K}$ is called divergent if, and only if, it is not convergent.

Example 7.2. (a) For every constant sequence $(z_n)_{n \in \mathbb{N}} = (a)_{n \in \mathbb{N}}$ with $a \in \mathbb{K}$, one has $\lim_{n \to \infty} z_n = \lim_{n \to \infty} a = a$: Since, for each $n \in \mathbb{N}$, $|z_n - a| = |a - a| = 0$, one can choose $N = 1$ for each $\epsilon > 0$.

(b) $\lim_{n \to \infty} \frac{1}{n+a} = 0$ for each $a \in \mathbb{C}$: Here $z_n := 1/(n+a)$ (if $n = -a$, then set $z_n := w$ with $w \in \mathbb{C}$ arbitrary). Given $\epsilon > 0$, choose an arbitrary $N \in \mathbb{N}$ with $N \geq \epsilon^{-1} + |a|$. Then, for each $n \geq N$, we compute $|n+a| = |n-(-a)| \geq |n-a| = n - |a| > N - |a| \geq \epsilon^{-1}$, and, thus, $|z_n| = |n+a|^{-1} < \epsilon$ as desired.
(c) \((-1)^n\)\(_{n \in \mathbb{N}}\) is not convergent: We have \(z_n = 1\) for each even \(n\) and \(z_n = -1\) for each odd \(n\). Thus, for each \(z \neq 1\) and each even \(n\), \(|z_n - z| = |1 - z| > |1 - z|/2 =: \epsilon > 0\), i.e. \(z\) is not a limit of \((z_n)_{n \in \mathbb{N}}\). However, \(z = 1\) is also not a limit of the sequence, since, for each odd \(n\), \(|z_n - 1| = |-1 - 1| = 2 > 1 =: \epsilon > 0\), proving that the sequence has no limit.

**Theorem 7.3.** (a) Let \((z_n)_{n \in \mathbb{N}}\) be a sequence in \(\mathbb{C}\). Then \((z_n)_{n \in \mathbb{N}}\) is convergent in \(\mathbb{C}\) if, and only if, both \((\text{Re } z_n)_{n \in \mathbb{N}}\) and \((\text{Im } z_n)_{n \in \mathbb{N}}\) are convergent in \(\mathbb{R}\). Moreover, in that case,

\[
\lim_{n \to \infty} z_n = z \iff \lim_{n \to \infty} \text{Re } z_n = \text{Re } z \land \lim_{n \to \infty} \text{Im } z_n = \text{Im } z. \tag{7.2}
\]

(b) Let \((x_n)_{n \in \mathbb{N}}\) be a sequence in \(\mathbb{R}\) and \(z \in \mathbb{C}\). Then

\[
\lim_{n \to \infty} x_n = z \implies z \in \mathbb{R}. \tag{7.3}
\]

**Proof.** (a): Suppose \((z_n)_{n \in \mathbb{N}}\) converges to \(z \in \mathbb{C}\). Then, given \(\epsilon > 0\), there exists \(N \in \mathbb{N}\) such that, for each \(n > N\), \(|z_n - z| < \epsilon\). In consequence, for each \(n > N\),

\[
|\text{Re } z_n - \text{Re } z| = |\text{Re}(z_n - z)| \leq |z_n - z| < \epsilon, \tag{7.4}
\]

proving \(\lim_{n \to \infty} \text{Re } z_n = \text{Re } z\). The proof of \(\lim_{n \to \infty} \text{Im } z_n = \text{Im } z\) is completely analogous. Conversely, suppose there are \(x, y \in \mathbb{R}\) such that \(\lim_{n \to \infty} \text{Re } z_n = x\) and \(\lim_{n \to \infty} \text{Im } z_n = y\). Here we encounter, for the first time, what is sometimes called an \(\epsilon/2\) argument: Given \(\epsilon > 0\), there exists \(N \in \mathbb{N}\) such that, for each \(n > N\), \(|\text{Re } z_n - x| < \epsilon/2\) and \(|\text{Im } z_n - y| < \epsilon/2\), implying, for each \(n > N\),

\[
|z_n - (x + iy)| = |\text{Re } z_n + i \text{Im } z_n - (x + iy)| \\
\leq |\text{Re } z_n - x| + |i| |\text{Im } z_n - y| < \epsilon/2 + \epsilon/2 = \epsilon, \tag{7.5}
\]

proving \(\lim_{n \to \infty} z_n = x + iy\).

(b) is a direct consequence of (a).

**Example 7.4.** (a) According to Th. 7.3(a), we have

\[
\lim_{n \to \infty} \left(\sqrt{2} + \frac{i}{n - 17}\right) = \sqrt{2} + 0i = \sqrt{2}. \tag{7.2(a),(b)}
\]

(b) According to Th. 7.3(a) and Ex. 7.2(c), the sequence \(\left(\frac{1}{n} + (-1)^n i\right)_{n \in \mathbb{N}}\) is divergent.

Another important example relies on the following inequality:

**Proposition 7.5** (Bernoulli’s Inequality). For each \(n \in \mathbb{N}_0\) and each \(x \in [-1, \infty[\), we have

\[
(1 + x)^n \geq 1 + nx, \tag{7.6}
\]

with strict inequality whenever \(n > 1\) and \(x \neq 0\).
If $n = 0$, (7.6) reads $1 \geq 1$, for $n = 1$, (7.6) reads $1 + x \geq 1 + x$, for $n = 2$, (7.6) reads $(1 + x)^2 = 1 + 2x + x^2 \geq 1 + 2x$, all three statements being trivially true, in the case $n = 2$ with strict inequality for $x \neq 0$. We now proceed by induction for $n \geq 2$.

For the induction step, one estimates

$$(1 + x)^{n+1} = (1 + x)^n (1 + x)^{\text{ind. hyp.}} \geq (1 + nx) (1 + x) = 1 + (n + 1)x + nx^2$$

$$(7.7)$$

with strict inequality for $x \neq 0$.

**Example 7.6.** We have, for each $q \in \mathbb{C}$,

$$|q| < 1 \quad \Rightarrow \quad \lim_{n \to \infty} q^n = 0 : (7.8)$$

For $q = 0$, there is nothing to prove. For $0 < |q| < 1$, it is $|q|^{-1} > 1$, i.e. $h := |q|^{-1} - 1 > 0$. Thus, for each $\epsilon > 0$ and $N \geq 1/(\epsilon h)$, we obtain

$$n > N \quad \Rightarrow \quad |q|^{-n} = (1 + h)^n \geq 1 + nh > nh > 1/\epsilon \quad \Rightarrow \quad |q|^n = |q|^{-n} < \epsilon. \quad (7.9)$$

**Definition 7.7.** (a) Given $z \in \mathbb{K}$ and $\epsilon \in \mathbb{R}^+$, we call the set $B_\epsilon(z) := \{w \in \mathbb{K} : |w - z| < \epsilon\}$ the $\epsilon$-neighborhood of $z$ or, in anticipation of Calculus II, the (open) $\epsilon$-ball with center $z$ (in fact, for $\mathbb{K} = \mathbb{C}$, $B_\epsilon(z)$ represents an open disk in the complex plane with center $z$ and radius $\epsilon$, whereas, for $\mathbb{K} = \mathbb{R}$, $B_\epsilon(z) = [z - \epsilon, z + \epsilon]$ is the open interval with center $z$ and length $2\epsilon$). More generally, a set $U \subseteq \mathbb{K}$ is called a neighborhood of $z$ if, and only if, there exists $\epsilon > 0$ with $B_\epsilon(z) \subseteq U$ (so, for example, for $\epsilon > 0$, $B_\epsilon(z)$ is always a neighborhood of $z$, whereas $\mathbb{R}$ and $[z - \epsilon, \infty[$ are neighborhoods of $z$ for $\mathbb{K} = \mathbb{R}$, but not for $\mathbb{K} = \mathbb{C}$ ($[z - \epsilon, \infty[$ not even being defined for $z \notin \mathbb{R}$); the sets $\{z\}$, $\{w \in \mathbb{K} : \text{Re} w \geq \text{Re} z\}$, $\{w \in \mathbb{K} : \text{Re} w \geq \text{Re} z + \epsilon\}$ are never neighborhoods of $z$).

(b) If $\phi(n)$ is a statement for each $n \in \mathbb{N}$, then $\phi(n)$ is said to be true for almost all $n \in \mathbb{N}$ if, and only if, there exists a finite subset $A \subseteq \mathbb{N}$ such that $\phi(n)$ is true for each $n \in \mathbb{N} \setminus A$, i.e. if, and only if, $\phi(n)$ is always true, with the possible exception of finitely many cases.

**Remark 7.8.** In the language of Def. 7.7, the sequence $(z_n)_{n \in \mathbb{N}}$ converges to $z$ if, and only if, every neighborhood of $z$ contains almost all $z_n$.

**Definition 7.9.** The sequence $(z_n)_{n \in \mathbb{N}}$ in $\mathbb{K}$ is called bounded if, and only if, the set $\{|z_n| : n \in \mathbb{N}\}$ is bounded in the sense of Def. 2.26(a).

**Proposition 7.10.** Let $(z_n)_{n \in \mathbb{N}}$ be a sequence in $\mathbb{K}$.

(a) Limits are unique, that means if $z, w \in \mathbb{K}$ such that $\lim_{n \to \infty} z_n = z$ and $\lim_{n \to \infty} z_n = w$, then $z = w$.

(b) If $(z_n)_{n \in \mathbb{N}}$ is convergent, then it is bounded.
Proof. (a): Exercise.

(b): If \( \lim_{n \to \infty} z_n = z \), then \( A := \{|z_n| : |z_n - z| \geq 1\} \cup \{|z|\} \) is nonempty and finite. According to Th. 3.21(a), \( A \) has an upper bound \( M \). Then \( \max\{M, |z| + 1\} \) is an upper bound for \( \{|z_n| : n \in \mathbb{N}\} \), and 0 is always a lower bound, showing that the sequence is bounded.

**Proposition 7.11.** Let \( (z_n)_{n \in \mathbb{N}} \) be a sequence in \( \mathbb{C} \) with \( \lim_{n \to \infty} z_n = 0 \).

(a) If \( (b_n)_{n \in \mathbb{N}} \) is a sequence in \( \mathbb{C} \) such that there exists \( C \in \mathbb{R}^+ \) with \( |b_n| \leq C|z_n| \) for almost all \( n \), then \( \lim_{n \to \infty} b_n = 0 \).

(b) If \( (c_n)_{n \in \mathbb{N}} \) is a bounded sequence in \( \mathbb{C} \), then \( \lim_{n \to \infty} (c_n z_n) = 0 \).

**Proof.** (a): Given \( \epsilon > 0 \), there exists \( N \in \mathbb{N} \) such that \( |z_n| < \epsilon/C \) and \( |b_n| \leq C|z_n| \) for each \( n > N \). Then, for each \( n > N \), \( |b_n| \leq C|z_n| < \epsilon \), proving \( \lim_{n \to \infty} b_n = 0 \).

(b): If \( (c_n)_{n \in \mathbb{N}} \) is bounded, then there exists \( C \in \mathbb{R}^+ \) such that \( |c_n| \leq C \) for each \( n \in \mathbb{N} \). Thus, \( |c_n z_n| \leq C|z_n| \) for each \( n \in \mathbb{N} \), implying \( \lim_{n \to \infty} (c_n z_n) = 0 \) via (a).

**Example 7.12.** The sequences \( ((-1)^n)_{n \in \mathbb{N}} \) and \( (b)_{n \in \mathbb{N}} \) with \( b \in \mathbb{C} \) are bounded. Since, for each \( a \in \mathbb{C} \), \( \lim_{n \to \infty} \frac{1}{n+a} = 0 \) by Example 7.2(b), we obtain

\[
\lim_{n \to \infty} \frac{(-1)^n}{n+a} = \lim_{n \to \infty} \frac{b}{n+a} = 0
\]  

(7.10)

from Prop. 7.11(b).

**Theorem 7.13.** (a) Let \( (z_n)_{n \in \mathbb{N}} \) and \( (w_n)_{n \in \mathbb{N}} \) be sequences in \( \mathbb{C} \). Moreover, let \( z, w \in \mathbb{C} \) with \( \lim_{n \to \infty} z_n = z \) and \( \lim_{n \to \infty} w_n = w \). We have the following identities:

\[
\lim_{n \to \infty} (\lambda z_n) = \lambda z \quad \text{for each } \lambda \in \mathbb{C},
\]

(7.11a)

\[
\lim_{n \to \infty} (z_n + w_n) = z + w,
\]

(7.11b)

\[
\lim_{n \to \infty} (z_n w_n) = zw,
\]

(7.11c)

\[
\lim_{n \to \infty} \frac{z_n}{w_n} = \frac{z}{w} \quad \text{given all } w_n \neq 0 \text{ and } w \neq 0,
\]

(7.11d)

\[
\lim_{n \to \infty} |z_n| = |z|,
\]

(7.11e)

\[
\lim_{n \to \infty} \bar{z}_n = \bar{z},
\]

(7.11f)

\[
\lim_{n \to \infty} z_n^p = z^p \quad \text{for each } p \in \mathbb{N}.
\]

(7.11g)

(b) Let \( (x_n)_{n \in \mathbb{N}} \) and \( (y_n)_{n \in \mathbb{N}} \) be sequences in \( \mathbb{R} \). Moreover, let \( x, y \in \mathbb{R} \) with \( \lim_{n \to \infty} x_n = x \) and \( \lim_{n \to \infty} y_n = y \). Then

\[
\lim_{n \to \infty} \max\{x_n, y_n\} = \max\{x, y\},
\]

(7.12a)

\[
\lim_{n \to \infty} \min\{x_n, y_n\} = \min\{x, y\}.
\]

(7.12b)
(c) If, in the situation of (b) (i.e. for real sequences), $x_n \leq y_n$ holds for almost all $n \in \mathbb{N}$, then $x \leq y$. In particular, if almost all $x_n \geq 0$, then $x \geq 0$.

Proof. We start with the identities of (a).

(7.11a): For $\lambda = 0$, there is nothing to prove. For $\lambda \neq 0$ and $\epsilon > 0$, there exists $N \in \mathbb{N}$ such that, for each $n > N$, $|z_n - z| < \epsilon/|\lambda|$, implying

$$\forall_{n>N} |\lambda z_n - \lambda z| = |\lambda| |z_n - z| < \epsilon.$$  \hfill (7.13a)

(7.11b): Given $\epsilon > 0$, there exists $N \in \mathbb{N}$ such that, for each $n > N$, $|z_n - z| < \epsilon/2$ and $|w_n - w| < \epsilon/2$, implying

$$\forall_{n>N} |z_n + w_n - (z + w)| \leq |z_n - z| + |w_n - w| < \epsilon/2 + \epsilon/2 = \epsilon.$$  \hfill (7.13b)

(7.11c): Let $M_1 := \max\{|z|, 1\}$. According to Prop. 7.10(b), there exists $M_2 \in \mathbb{R}^+$ such that $M_2$ is an upper bound for $\{|w_n| : n \in \mathbb{N}\}$. Moreover, given $\epsilon > 0$, there exists $N \in \mathbb{N}$ such that, for each $n > N$, $|z_n - z| < \epsilon/(2M_2)$ and $|w_n - w| < \epsilon/(2M_1)$, implying

$$\forall_{n>N} \left( |z_n w_n - zw| = |(z_n - z)w_n + z(w_n - w)| \leq |w_n| \cdot |z_n - z| + |z| \cdot |w_n - w| < \frac{M_2 \epsilon}{2M_2} + \frac{M_1 \epsilon}{2M_1} = \epsilon \right).$$  \hfill (7.13c)

(7.11d): We first consider the case, where all $z_n = 1$. Given $\epsilon > 0$, there exists $N \in \mathbb{N}$ such that, for each $n > N$, $|w_n - w| < \epsilon |w|^2/2$ and $|w_n - w| < |w|/2$ (since $w \neq 0$ for this case), implying $|w| \leq |w - w_n| + |w_n| < |w|/2 + |w_n|$ and $|w_n| > |w|/2$. Thus,

$$\forall_{n>N} \left| \frac{1}{w_n} - \frac{1}{w} \right| = \left| \frac{w_n - w}{w_n w} \right| \leq \frac{2 |w_n - w|}{|w|^2} < \frac{2 \epsilon |w|^2}{2} = \epsilon.$$  \hfill (7.13d)

The general case now follows from (7.11c).

(7.11e): This is a consequence of the inverse triangle inequality (5.14): Given $\epsilon > 0$, there exists $N \in \mathbb{N}$ such that, for each $n > N$, $|z_n - z| < \epsilon$, implying

$$\forall_{n>N} ||z_n| - |z|| \leq |z_n - z| < \epsilon.$$  \hfill (7.13e)

(7.11f): Write $z_n = x_n + iy_n$ and $z = x + iy$ with $x_n, y_n, x, y \in \mathbb{R}$, $n \in \mathbb{N}$. Then we know $\lim_{n \to \infty} x_n = x$ and $\lim_{n \to \infty} y_n = y$ from (7.2), and

$$\lim_{n \to \infty} z_n = \lim_{n \to \infty} (x_n - iy_n) \overset{\text{(7.11a), (7.11b)}}{=} (x - iy) = \bar{z},$$  \hfill (7.13f)

which establishes the case.

(7.11g) follows by induction from (7.11c) (cf. (7.16b) below).

The proofs for the two identities of (b) are left as exercises.

(c): Proceeding by contraposition, assume $x > y$ and set $s := (x + y)/2$. Then $y < s < x$ and $y_n < s < x_n$ holds for almost all $n$, i.e. $x_n \leq y_n$ does not hold for almost all $n$. \qed
Example 7.14. (a) \( \lim_{n \to \infty} \frac{n+a}{n+b} = 1 \) for each \( a, b \in \mathbb{C} \): Here \( z_n := (n+a)/(n+b) \) (if \( n = -b \), then set \( z_n := w \) with \( w \in \mathbb{C} \) arbitrary). Using (7.11b) and (7.11d), one obtains
\[
\lim_{n \to \infty} \frac{n+a}{n+b} = \lim_{n \to \infty} \frac{1+a/n}{1+b/n} = \lim_{n \to \infty} 1 + \lim_{n \to \infty} \frac{a}{n} = \frac{1+0}{1+0} = 1. \tag{7.14}
\]

(b) Using (7.11b), (7.11d), and (7.11g), one obtains
\[
\lim_{n \to \infty} \frac{2n^5 - 3in^3 + 2i}{3n^5 + 17n} = \lim_{n \to \infty} \frac{2 - 3i/n^2 + 2i/n^5}{3 + 17/n^4} = \frac{2 + 0 + 0}{3 + 0} = \frac{2}{3}. \tag{7.15}
\]

Corollary 7.15. For \( k \in \mathbb{N} \), let \( (z_n^{(1)})_{n \in \mathbb{N}}, \ldots, (z_n^{(k)})_{n \in \mathbb{N}} \) be sequences in \( \mathbb{C} \). Moreover, let \( z^{(1)}, \ldots, z^{(k)} \in \mathbb{C} \) with \( \lim_{n \to \infty} z_n^{(j)} = z^{(j)} \) for each \( j \in \{1, \ldots, k\} \). Then
\[
\lim_{n \to \infty} \sum_{j=1}^{k} z_n^{(j)} = \sum_{j=1}^{k} z^{(j)}, \tag{7.16a}
\]
\[
\lim_{n \to \infty} \prod_{j=1}^{k} z_n^{(j)} = \prod_{j=1}^{k} z^{(j)}. \tag{7.16b}
\]

Proof. (7.16) follows by simple inductions from (7.11b) and (7.11c), respectively. \( \blacksquare \)

Theorem 7.16 (Sandwich Theorem). Let \( (x_n)_{n \in \mathbb{N}}, (y_n)_{n \in \mathbb{N}}, \) and \( (a_n)_{n \in \mathbb{N}} \) be sequences in \( \mathbb{R} \). If \( x_n \leq a_n \leq y_n \) holds for almost all \( n \in \mathbb{N} \), then
\[
\lim_{n \to \infty} x_n = \lim_{n \to \infty} y_n = x \in \mathbb{R} \quad \Rightarrow \quad \lim_{n \to \infty} a_n = x. \tag{7.17}
\]

Proof. Given \( \epsilon > 0 \), there exists \( N \in \mathbb{N} \) such that, for each \( n > N \), \( x_n \leq a_n \leq y_n \), \( |x_n - x| < \epsilon \), and \( |y_n - x| < \epsilon \), implying
\[
\forall \ n > N \quad x - \epsilon < x_n \leq a_n \leq y_n < x + \epsilon, \tag{7.18}
\]
which establishes the case. \( \blacksquare \)

Example 7.17. Since, 0 < \( \frac{1}{n!} \leq \frac{1}{n} \) holds for each \( n \in \mathbb{N} \), the Sandwich Th. 7.16 implies
\[
\lim_{n \to \infty} \frac{1}{n!} = 0. \tag{7.19}
\]

Definition 7.18. Let \( (x_n)_{n \in \mathbb{N}} \) be a sequence in \( \mathbb{R} \). The sequence is said to diverge to \( \infty \) (resp. to \( -\infty \)), denoted \( \lim_{n \to \infty} x_n = \infty \) (resp. \( \lim_{n \to \infty} x_n = -\infty \)) if, and only if, for each \( K \in \mathbb{R} \), almost all \( x_n \) are bigger (resp. smaller) than \( K \). Thus,
\[
\lim_{n \to \infty} x_n = \infty \iff \forall \ K \in \mathbb{R} \ \exists \ N \in \mathbb{N} \ \forall \ n > N \quad x_n > K, \tag{7.20a}
\]
\[
\lim_{n \to \infty} x_n = -\infty \iff \forall \ K \in \mathbb{R} \ \exists \ N \in \mathbb{N} \ \forall \ n > N \quad x_n < K. \tag{7.20b}
\]
Theorem 7.19. Suppose $S := (x_n)_{n \in \mathbb{N}}$ is a monotone sequence in $\mathbb{R}$ (increasing or decreasing). Defining $A := \{x_n : n \in \mathbb{N}\}$, the following holds:

\[
\lim_{n \to \infty} x_n = \begin{cases} 
\sup A & \text{if } S \text{ is increasing and bounded,} \\
\infty & \text{if } S \text{ is increasing and not bounded,} \\
\inf A & \text{if } S \text{ is decreasing and bounded,} \\
-\infty & \text{if } S \text{ is decreasing and not bounded.}
\end{cases}
\]  

(7.21)

Proof. We treat the increasing case; the decreasing case is proved completely analogously. If $A$ is bounded and $\epsilon > 0$, let $K := \sup A - \epsilon$; if $A$ is unbounded, then let $K \in \mathbb{R}$ be arbitrary. In both cases, since $K$ can not be an upper bound, there exists $N \in \mathbb{N}$ such that $x_N > K$. Since the sequence is increasing, for each $n > N$, $x_N \leq x_n$, showing $|\sup A - x_n| < \epsilon$ in the bounded case, and $x_n > K$ in the unbounded case.

Example 7.20. Theorem 7.19 implies

\[
\forall k \in \mathbb{N}, \quad \lim_{n \to \infty} n^k = \infty, \quad \lim_{n \to \infty} (-n^k) = -\infty.
\]  

(7.22)

It is sometimes necessary to consider so-called subsequences and reorderings of a given sequence. Here, we are interested in sequences in $\mathbb{R}$ or $\mathbb{C}$, but for subsequences and reorderings it is irrelevant in which set $A$ the sequence takes its values. As it presents virtually no extra difficulty to introduce the notions for general sequences, and since we will need to consider sequences with values in sets other than $\mathbb{R}$ or $\mathbb{C}$ in Calculus II, we admit general sequences in the following definition.

Definition 7.21. Let $A$ be an arbitrary nonempty set. Consider a sequence $\sigma : \mathbb{N} \to A$. Given a function $\phi : \mathbb{N} \to \mathbb{N}$ (that means $(\phi(n))_{n \in \mathbb{N}}$ constitutes a sequence of indices), the new sequence $(\sigma \circ \phi) : \mathbb{N} \to A$ is called a subsequence of $\sigma$ if, and only if, $\phi$ is strictly increasing (i.e. $1 \leq \phi(1) < \phi(2) < \ldots$). Moreover, $\sigma \circ \phi$ is called a reordering of $\sigma$ if, and only if, $\phi$ is bijective. One can write $\sigma$ in the form $(z_n)_{n \in \mathbb{N}}$ by setting $z_n := \sigma(n)$, and one can write $\sigma \circ \phi$ in the form $(w_n)_{n \in \mathbb{N}}$ by setting $w_n := (\sigma \circ \phi)(n) = z_{\phi(n)}$. Especially for a subsequence of $(z_n)_{n \in \mathbb{N}}$, it is also common to write $(z_{n_k})_{k \in \mathbb{N}}$. This notation corresponds to the one above if one lets $n_k := \phi(k)$. Analogous definitions work if the index set $\mathbb{N}$ of $\sigma$ is replaced by a general countable nonempty index set $I$.

Example 7.22. Consider the sequence $(1, 2, 3, \ldots)$. Then $(2, 4, 6, \ldots)$ constitutes a subsequence and $(2, 1, 4, 3, 6, 5, \ldots)$ constitutes a reordering. Using the notation of Def. 7.21, the original sequence is given by $\sigma : \mathbb{N} \to \mathbb{N}$, $\sigma(n) := n$; the subsequence is selected via $\phi_1 : \mathbb{N} \to \mathbb{N}$, $\phi_1(n) := 2n$; and the reordering is accomplished via $\phi_2 : \mathbb{N} \to \mathbb{N}$, $\phi_2(n) := \begin{cases} 
n + 1 & \text{if } n \text{ is odd,} \\
n - 1 & \text{if } n \text{ is even.}
\end{cases}$

Proposition 7.23. Let $(z_n)_{n \in \mathbb{N}}$ be a sequence in $\mathbb{C}$. If $\lim_{n \to \infty} z_n = z$, then every subsequence and every reordering of $(z_n)_{n \in \mathbb{N}}$ is also convergent with limit $z$. 
Conversely, if \( z \phi \) follows: For \( z m < n n>M \), one has contains infinitely many members of the sequence (i.e. \( \#\{n \in \mathbb{N}: z_n \in B_\varepsilon(z)\} = \infty \)).

Let \( (w_n)_{n \in \mathbb{N}} \) be a reordering of \( (z_n)_{n \in \mathbb{N}} \), i.e. there is a bijective function \( \phi: \mathbb{N} \rightarrow \mathbb{N} \) such that \( w_n = z_\phi(n) \). Let \( \epsilon \) and \( N \) be as before. Define

\[
M := \max\{\phi^{-1}(n) : n \leq N\}.
\]

(7.23)

As \( \phi \) is bijective, it is \( \phi(n) > N \) for each \( n > M \). Then, for each \( n > M \), one has \( w_n = z_\phi(n) \in B_\varepsilon(z) \), showing \( \lim_{n \to \infty} w_n = z \).

Definition 7.24. Let \( (z_n)_{n \in \mathbb{N}} \) be a sequence in \( \mathbb{K} \). A point \( z \in \mathbb{K} \) is called a cluster point or an accumulation point of the sequence if, and only if, for each \( \epsilon > 0 \), \( B_\varepsilon(z) \) contains infinitely many members of the sequence (i.e. \( \#\{n \in \mathbb{N}: z_n \in B_\varepsilon(z)\} = \infty \)).

Example 7.25. The sequence \((-1)^n)_{n \in \mathbb{N}}\) has cluster points 1 and \(-1\).

Proposition 7.26. A point \( z \in \mathbb{K} \) is a cluster point of the sequence \( (z_n)_{n \in \mathbb{N}} \) in \( \mathbb{K} \) if, and only if, the sequence has a subsequence converging to \( z \).

Proof. If \( (w_n)_{n \in \mathbb{N}} \) is a subsequence of \( (z_n)_{n \in \mathbb{N}} \), \( \lim_{n \to \infty} w_n = z \), then every \( B_\varepsilon(z) \), \( \epsilon > 0 \), contains infinitely many \( w_n \), i.e. infinitely many \( z_n \), i.e. \( z \) is a cluster point of \( (z_n)_{n \in \mathbb{N}} \). Conversely, if \( z \) is a cluster point of \( (z_n)_{n \in \mathbb{N}} \), then, inductively, define \( \phi: \mathbb{N} \rightarrow \mathbb{N} \) as follows: For \( \phi(1) \), choose the index \( k \) of any point \( z_k \) in \( B_\varepsilon(z) \) (such a point exists, since \( z \) is a cluster point of the sequence). Now assume that \( n > 1 \) and that \( \phi(m) \) have already been defined for each \( m < n \). Let \( M := \max\{\phi(m) : m < n\} \). Since \( B_\frac{1}{n}(z) \) contains infinitely many \( z_k \), there must be some \( z_k \in B_1(z) \) such that \( k > M \). Choose this \( k \) as \( \phi(n) \).

Thus, by construction, \( \phi \) is strictly increasing, i.e. \( (w_n)_{n \in \mathbb{N}} \) with \( w_n := z_\phi(n) \) is a subsequence of \( (z_n)_{n \in \mathbb{N}} \). Moreover, for each \( \epsilon > 0 \), there is \( N \in \mathbb{N} \) such that \( 1/N < \epsilon \). Then, for each \( n > N \), \( w_n \in B_\frac{1}{n}(z) \subseteq B_1(z) \subseteq B_\varepsilon(z) \), showing \( \lim_{n \to \infty} w_n = z \).

Theorem 7.27 (Bolzano-Weierstrass). Every bounded sequence \( S := (x_n)_{n \in \mathbb{N}} \) in \( \mathbb{K} \) has at least one cluster point in \( \mathbb{K} \). Moreover, for \( \mathbb{K} = \mathbb{R} \), the set \( A := \{x \in \mathbb{R} : x \text{ is cluster point of } S\} \) has a max \( x^* \in \mathbb{R} \) and a min \( x_* \in \mathbb{R} \), i.e. every bounded sequence in \( \mathbb{R} \) has a largest and a smallest cluster point. In addition, for each \( \epsilon > 0 \), the inequality \( x_* - \epsilon < x_n < x^* + \epsilon \) holds for almost all \( n \).

Proof. We first consider the case \( \mathbb{K} = \mathbb{R} \). Define

\[
A^* := \{x \in \mathbb{R} : x_n \leq x \text{ for almost all } n\},
\]

(7.24a)

\[
A_* := \{x \in \mathbb{R} : x_n \geq x \text{ for almost all } n\}.
\]

(7.24b)

We claim \( A^* \neq \emptyset \) is bounded from below and \( x^* = \max A = \inf A^* ; A_* \neq \emptyset \) is bounded from above and \( x_* = \min A = \sup A_* \). We prove the claim for \( A^* \) - the proof for \( A_* \) is...
conducted completely analogous. Let \( m, M \in \mathbb{R} \) be a lower and an upper bound for \( S \), respectively. Then \( M \in A^* \), showing \( A^* \neq \emptyset \); and \( m \) is a lower bound for \( A^* \). Since \( A^* \) is bounded from below, \( a := \inf A^* \in \mathbb{R} \) by the completeness of \( \mathbb{R} \). Moreover, for each \( \epsilon > 0 \), \( a - \epsilon \notin A^* \), as \( a \) is a lower bound for \( A^* \), i.e. \( x_n > a - \epsilon \) holds for infinitely many \( n \in \mathbb{N} \). On the other hand, \( a + \epsilon/2 \in A^* \) follows from \( a \) being the largest lower bound of \( A^* \), i.e. \( x_n > a + \epsilon/2 \) holds for only finitely many \( n \) (if any). In particular, we have shown \( x_n < a + \epsilon \) holds for almost all \( n \), and \( a - \epsilon < x_n < a + \epsilon \) must hold for infinitely many \( n \), showing \( a \) is a cluster point of \( S \). To see that \( a \) is the largest cluster point of \( S \) (i.e. \( a = \max A \)), we have to show that \( x > a \) implies \( x \) is not a cluster point of \( S \). However, letting \( \epsilon := x - a > 0 \), we had seen above that \( x_n > a + \epsilon/2 \) holds for only finitely many \( n \), i.e. \( B_{\epsilon/2}(x) \) contains only finitely many \( x_n \), showing \( x \) is not a cluster point of \( S \).

It now remains to consider the complex case, i.e. a bounded sequence \( (z_n)_{n \in \mathbb{N}} \) in \( \mathbb{C} \). For each \( n \in \mathbb{N} \), let \( z_n = x_n + iy_n \) with \( x_n, y_n \in \mathbb{R} \). Due to Th. 5.11(d), we have \(|x_n| \leq |z_n|\) and \(|y_n| \leq |z_n|\), i.e. the boundedness of \( S \) implies the boundedness of both \((x_n)_{n \in \mathbb{N}}\) and \((y_n)_{n \in \mathbb{N}}\). Then we know that \((x_n)_{n \in \mathbb{N}}\) has a cluster point \( x \) and, by Prop. 7.26, \( S \) has a subsequence \((z_{n_j})_{j \in \mathbb{N}}\) such that \( x = \lim_{j \to \infty} x_{n_j} \). As the subsequence \((y_{n_j})_{j \in \mathbb{N}}\) is still bounded, it must have a cluster point \( y \) and a subsequence \((y_{n_{j_k}})_{k \in \mathbb{N}}\) such that \( y = \lim_{k \to \infty} y_{n_{j_k}} \). Since \( x = \lim_{k \to \infty} x_{n_{j_k}} \) as well, we now have \( \lim_{k \to \infty} z_{n_{j_k}} = x + iy =: z \), i.e. \( S \) has a subsequence converging to \( z \). According to Prop. 7.26, \( z \) is a cluster point of \( S \).

**Definition 7.28.** A sequence \((z_n)_{n \in \mathbb{N}}\) in \( \mathbb{C} \) is defined to be a Cauchy sequence if, and only if, for each \( \epsilon \in \mathbb{R}^+ \), there exists \( N \in \mathbb{N} \) such that \( |z_n - z_m| < \epsilon \) for each \( n, m > N \), i.e.

\[
(z_n)_{n \in \mathbb{N}} \text{ Cauchy } \iff \forall \epsilon \in \mathbb{R}^+ \; \exists N \in \mathbb{N} \; n, m > N \; |z_n - z_m| < \epsilon. \tag{7.25}
\]

**Theorem 7.29.** The sequence \((z_n)_{n \in \mathbb{N}}\) in \( \mathbb{C} \) is convergent if, and only if, it is a Cauchy sequence.

**Proof.** Suppose the sequence is convergent with \( \lim_{n \to \infty} z_n = z \). Then, given \( \epsilon > 0 \), there is \( N \in \mathbb{N} \) such that \( |z_n - z| < \epsilon/2 \) for each \( n > N \). If \( n, m > N \), then \( |z_n - z_m| \leq |z_n - z| + |z - z_m| < \frac{\epsilon}{2} + \frac{\epsilon}{2} = \epsilon \), establishing that \((z_n)_{n \in \mathbb{N}}\) is a Cauchy sequence.

Conversely, suppose the sequence is a Cauchy sequence. Using similar reasoning as in the proof of Prop. 7.10(b), we first show the sequence is bounded. If the sequence is Cauchy, then there exists \( N \in \mathbb{N} \) such that \( |z_n - z_m| < 1 \) for all \( n, m > N \). Thus, the set \( A := \{|z_n| : |z_n - z_{N+1}| \geq 1\} \cup \{|z_1|\} \subseteq \mathbb{R}^+ \) is nonempty and finite. According to Th. 3.21(a), \( A \) has an upper bound \( M \). Then \( \max \{M, |z_{N+1}| + 1\} \) is an upper bound for \( \{|z_n| : n \in \mathbb{N}\} \), showing that the sequence is bounded. From Th. 7.27, we obtain that the sequence has a cluster point \( z \). It remains to show \( \lim_{n \to \infty} z_n = z \). Given \( \epsilon > 0 \), choose \( N \in \mathbb{N} \) such that \( |z_n - z_m| < \epsilon/2 \) for all \( n, m > N \). Since \( z \) is a cluster point, there exists \( k > N \) such that \( |z_k - z| < \epsilon/2 \). Thus,

\[
\forall n > N \; |z_n - z| \leq |z_n - z_k| + |z_k - z| < \frac{\epsilon}{2} + \frac{\epsilon}{2} = \epsilon, \tag{7.26}
\]

proving \( \lim_{n \to \infty} z_n = z \).
Example 7.30. Consider the sequence $S := (s_n)_{n \in \mathbb{N}}$ defined by
\[ s_n := \sum_{k=1}^{n} \frac{1}{k} = 1 + \frac{1}{2} + \cdots + \frac{1}{n}. \] (7.27)

We claim $S$ is not a Cauchy sequence and, thus, not convergent by Th. 7.29: For each $N \in \mathbb{N}$, we find $n, m > N$ such that $s_n - s_m > 1/2$, namely $m = N + 1$ and $n = 2(N+1)$:
\[ s_{2(N+1)} - s_{N+1} = \sum_{k=N+2}^{2(N+1)} \frac{1}{k} = \frac{1}{N + 2} + \frac{1}{N + 3} + \cdots + \frac{1}{2(N+1)} > (N + 1) \cdot \frac{1}{2(N+1)} = \frac{1}{2}. \] (7.28)

While we have just seen that $S$ is not convergent, it is clearly increasing, i.e. Th. 7.19 implies $S$ is unbounded and $\lim_{n \to \infty} s_n = \infty$. Sequences defined by longer and longer sums are known as series and will be studied further in Sec. 7.3 below. The series of the present example is known as the harmonic series. It has become famous as the simplest example of a series that does not converge even though its summands converge to 0. In terms of the notation introduced in Sec. 7.3 below, we have shown
\[ \sum_{k=1}^{\infty} \frac{1}{k} = 1 + \frac{1}{2} + \frac{1}{3} + \cdots = \infty. \] (7.29)

7.2 Continuity

7.2.1 Definitions and First Examples

Roughly, a function is continuous if a small change in its input results in a small change of its output. For functions defined on an interval, the notion of continuity makes precise the idea of a function having no jump – no discontinuity – at some point $x$ in its domain. For example, we would say the sign function of (5.9) has precisely one jump – one discontinuity – at $x = 0$, whereas quadratic functions (or, more generally, polynomials) do not have any jumps – they are continuous.

Definition 7.31. Let $M \subseteq \mathbb{C}$. If $\zeta \in M$, then a function $f : M \to \mathbb{K}$ is said to be continuous in $\zeta$ if, and only if, for each $\epsilon > 0$, there is $\delta > 0$ such that the distance between the values $f(z)$ and $f(\zeta)$ is less than $\epsilon$, provided the distance between $z$ and $\zeta$ is less than $\delta$, i.e. if, and only if,
\[ \forall \epsilon \in \mathbb{R}^+ \exists \delta \in \mathbb{R}^+ \forall z \in M \left( |z - \zeta| < \delta \implies |f(z) - f(\zeta)| < \epsilon \right). \] (7.30)

Moreover, $f$ is called continuous if, and only if, $f$ is continuous in every $\zeta \in M$. The set of all continuous functions from $f : M \to \mathbb{K}$ is denoted by $C(M, \mathbb{K})$, $C(M) := C(M, \mathbb{R})$.

Example 7.32. (a) Every constant map $f : M \to \mathbb{K}$, $\emptyset \neq M \subseteq \mathbb{C}$, is continuous: In this case, given $\epsilon$, we can choose any $\delta > 0$ we want, say $\delta := 42$: If $\zeta, z \in M$, then $|f(\zeta) - f(z)| = 0 < \epsilon$, which holds independently of $\delta$, in particular, if $|\zeta - z| < \delta$. 

(b) Every affine function $f : \mathbb{K} \rightarrow \mathbb{K}$, $f(z) := az + b$ is continuous: For $a = 0$, this follows from (a). For $a \neq 0$, given $\epsilon > 0$, choose $\delta := \epsilon / |a|$. Then,

$$\forall \zeta, z \in \mathbb{K} \left( |z - \zeta| < \delta = \frac{\epsilon}{|a|} \Rightarrow |f(z) - f(\zeta)| = |az + b - a\zeta - b| = |a||z - \zeta| < |a|\frac{\epsilon}{|a|} = \epsilon \right).$$

(7.31)

(c) The sign function of (5.9) is not continuous: It is continuous in each $\xi \in \mathbb{R} \setminus \{0\}$, but not continuous in 0: If $\xi \neq 0$, then, given $\epsilon > 0$, choose $\delta := |\xi|$. If $|x - \xi| < \delta$, then $\text{sgn}(x) = \text{sgn}(\xi)$, i.e. $|\text{sgn}(x) - \text{sgn}(\xi)| = 0 < \epsilon$, proving continuity in $\xi$. However, at 0, for $\epsilon := 1/2$, we have

$$\forall \delta > 0 \left| \text{sgn}(0) - \text{sgn}(\delta/2) \right| = |0 - 1| = 1 > \frac{1}{2} = \epsilon,$$

(7.32)

showing $\text{sgn}$ is not continuous in 0.

Some subtleties arise from the possibility that $f$ can be defined on subsets of $\mathbb{C}$ with very different properties. The notions introduced in Def. 7.33 help to deal with these subtleties.

**Definition 7.33.** Let $M \subseteq \mathbb{C}$.

(a) The point $z \in \mathbb{C}$ is called a cluster point or accumulation point of $M$ if, and only if, each $\epsilon$-neighborhood of $z$, $\epsilon \in \mathbb{R}^+$, contains infinitely many points of $M$, i.e. if, and only if,

$$\forall \right( \#(M \cap B_\epsilon(z)) = \infty).$$

(7.33)

Note: A cluster point of $M$ is not necessarily in $M$.

(b) The point $z$ is called an isolated point of $M$ if, and only if, there is $\epsilon \in \mathbb{R}^+$ such that $B_\epsilon(z) \cap M = \{z\}$. Note: An isolated point of $M$ is always in $M$.

**Proposition 7.34.** If $M \subseteq \mathbb{C}$, then each point of $M$ is either a cluster point or an isolated point of $M$, i.e.

$$M = \{z \in M : z \text{ cluster point of } M\} \cup \{z \in M : z \text{ isolated point of } M\}.$$

(7.34)

**Proof.** Consider $z \in M$ that is not a cluster point of $M$. We have to show that $z$ is an isolated point of $M$. Since $z$ is not a cluster point of $M$, there exists $\tilde{\epsilon} > 0$ such that $A := (M \cap B_{\tilde{\epsilon}}(z)) \setminus \{z\}$ is finite. Define

$$\epsilon := \begin{cases} \min \{|a - z| : a \in A\} & \text{if } A \neq \emptyset, \\ \tilde{\epsilon} & \text{if } A = \emptyset. \end{cases}$$

(7.35)

Then $B_{\epsilon}(z) \cap M = \{z\}$, showing $z$ is an isolated point of $M$. Finally, the union in (7.34) is clearly disjoint. $\blacksquare$
Lemma 7.35. Let $M \subseteq \mathbb{C}$, $f : M \rightarrow \mathbb{K}$. If $\zeta$ is an isolated point of $M$, then $f$ is always continuous in $\zeta$.

Proof. Independently of the concrete definition of $f$, we know there is $\delta > 0$ such that $B_\delta(\zeta) \cap M = \{\zeta\}$. In other words, if $z \in M$ with $|z - \zeta| < \delta$, then $z = \zeta$, implying $|f(z) - f(\zeta)| = 0 < \epsilon$ for each $\epsilon > 0$, showing $f$ to be continuous in $\zeta$.

Example 7.36. (a) The sign function restricted to the set $M := ]-\infty, -1[ \cup \{0\} \cup [1, \infty[$, i.e.

$$
\text{sgn}(x) = \begin{cases} 
1 & \text{for } x \in [1, \infty[, \\
0 & \text{for } x = 0, \\
-1 & \text{for } x \in ]-\infty, -1[
\end{cases}
$$

is continuous: As in Ex. 7.32(c), one sees that $\text{sgn}$ is continuous in each $\xi \in M \setminus \{0\}$. However, now it is also continuous in $0$, since $0$ is an isolated point of $M$.

(b) Every function $f : \mathbb{N} \rightarrow \mathbb{K}$ is continuous, since every $n \in \mathbb{N}$ is an isolated point of $\mathbb{N}$ (due to $\{n\} = \mathbb{N} \cap B_{\frac{1}{2}}(n)$).

7.2.2 Continuity, Sequences, and Function Arithmetic

To make available the power of the results on convergent sequences from Sec. 7.1 to investigations regarding the continuity of functions, we need to understand the relationship between both notions. The core of this relationship is the contents of the following Th. 7.37, which provides a criterion allowing one to test continuity in terms of convergent sequences:

Theorem 7.37. Let $M \subseteq \mathbb{C}$, $f : M \rightarrow \mathbb{K}$. If $\zeta \in M$, then $f$ is continuous in $\zeta$ if, and only if, for each sequence $(z_n)_{n \in \mathbb{N}}$ in $M$ with $\lim_{n \rightarrow \infty} z_n = \zeta$, the sequence $(f(z_n))_{n \in \mathbb{N}}$ converges to $f(\zeta)$, i.e.

$$
\lim_{n \rightarrow \infty} z_n = \zeta \quad \Rightarrow \quad \lim_{n \rightarrow \infty} f(z_n) = f(\zeta). 
$$

(7.36)

Proof. If $\zeta \in M$ is an isolated point of $M$, then there is $\delta > 0$ such that $M \cap B_\delta(\zeta) = \{\zeta\}$. Then every $f : M \rightarrow \mathbb{K}$ is continuous in $\zeta$ according to Lem. 7.35. On the other hand, every sequence in $M$ converging to $\zeta$ must be finally constant and equal to $\zeta$, i.e. (7.36) is trivially valid at $\zeta$. Thus, the assertion of the theorem holds if $\zeta \in M$ is an isolated point of $M$.

If $\zeta \in M$ is not an isolated point of $M$, then $\zeta$ is a cluster point of $M$ according to Prop. 7.34. So, for the remainder of the proof, let $\zeta \in M$ be a cluster point of $M$. Assume that $f$ is continuous in $\zeta$ and $(z_n)_{n \in \mathbb{N}}$ is a sequence in $M$ with $\lim_{n \rightarrow \infty} z_n = \zeta$. For each $\epsilon > 0$, there is $\delta > 0$ such that $z \in M$ and $|z - \zeta| < \delta$ implies $|f(z) - f(\zeta)| < \epsilon$. Since $\lim_{n \rightarrow \infty} z_n = \zeta$, there is also $N \in \mathbb{N}$ such that, for each $n > N$, $|z_n - \zeta| < \delta$. Thus, for each $n > N$, $|f(z_n) - f(\zeta)| < \epsilon$, proving $\lim_{n \rightarrow \infty} f(z_n) = f(\zeta)$. Conversely, assume that $f$ is not continuous in $\zeta$. We have to construct a sequence $(z_n)_{n \in \mathbb{N}}$ in $M$ with
\[ \lim_{n \to \infty} z_n = \zeta, \] but \((f(z_n))_{n \in \mathbb{N}}\) does not converge to \(f(\zeta)\). Since \(f\) is not continuous in \(\zeta\), there must be some \(\epsilon_0 > 0\) such that, for each \(1/n, n \in \mathbb{N}\), there is at least one \(z_n \in M\) satisfying \(|z_n - \zeta| < 1/n\) and \(|f(z_n) - f(\zeta)| \geq \epsilon_0\). Then \((z_n)_{n \in \mathbb{N}}\) is a sequence in \(M\) with \(\lim_{n \to \infty} z_n = \zeta\) and \((f(z_n))_{n \in \mathbb{N}}\) does not converge to \(f(\zeta)\). \[\blacksquare\]

We can now apply the rules of Th. 7.13 to see that all the arithmetic operations defined in Not. 6.2 preserve continuity:

**Theorem 7.38.** Let \(M \subseteq \mathbb{C}\), \(f, g : M \to \mathbb{K}\), \(\lambda \in \mathbb{K}\), \(\zeta \in M\). If \(f, g\) are both continuous in \(\zeta\), then \(\lambda f, f + g, fg, f/g\) for \(g \neq 0\), \(|f|, \text{Re}\, f, \text{Im}\, f\) are all continuous in \(\zeta\). If \(\mathbb{K} = \mathbb{R}\), then \(\max(f, g), \min(f, g), f^+\) and \(f^-\), are also all continuous in \(\zeta\).

**Proof.** Let \((z_n)_{n \in \mathbb{N}}\) be a sequence in \(M\) such that \(\lim_{n \to \infty} z_n = \zeta\). Then the continuity of \(f\) and \(g\) in \(\zeta\) yields \(\lim_{n \to \infty} f(z_n) = f(\zeta)\) and \(\lim_{n \to \infty} g(z_n) = g(\zeta)\). Then

\[
\begin{align*}
(7.11a) \Rightarrow & \quad \lim_{n \to \infty} (\lambda f)(z_n) = (\lambda f)(\zeta), \\
(7.11b) \Rightarrow & \quad \lim_{n \to \infty} (f + g)(z_n) = (f + g)(\zeta), \\
(7.11c) \Rightarrow & \quad \lim_{n \to \infty} (fg)(z_n) = (fg)(\zeta), \\
(7.11d) \Rightarrow & \quad \lim_{n \to \infty} (f/g)(z_n) = (f/g)(\zeta), \\
(7.11e) \Rightarrow & \quad \lim_{n \to \infty} |f|(z_n) = |f|(\zeta), \\
(7.2) \Rightarrow & \quad \lim_{n \to \infty} (\text{Re}\, f)(z_n) = (\text{Re}\, f)(\zeta), \\
(7.2) \Rightarrow & \quad \lim_{n \to \infty} (\text{Im}\, f)(z_n) = (\text{Im}\, f)(\zeta).
\end{align*}
\]

If \(f, g\) are both \(\mathbb{R}\)-valued, then we also have

\[
\begin{align*}
(7.12a) \Rightarrow & \quad \lim_{n \to \infty} \max(f, g)(z_n) = \max(f, g)(\zeta), \\
(7.12b) \Rightarrow & \quad \lim_{n \to \infty} \min(f, g)(z_n) = \min(f, g)(\zeta),
\end{align*}
\]

and, finally, the continuity of \(f^+\) and \(f^-\) follows from the continuity of \(\max(f, g)\). \[\blacksquare\]

**Corollary 7.39.** A function \(f : M \to \mathbb{C}\), \(M \subseteq \mathbb{C}\), is continuous in \(\zeta \in M\) if, and only if, both \(\text{Re}\, f\) and \(\text{Im}\, f\) are continuous in \(\zeta\).

**Proof.** If \(f\) is continuous in \(\zeta\), then \(\text{Re}\, f\) and \(\text{Im}\, f\) are both continuous in \(\zeta\) by Th. 7.38. If \(\text{Re}\, f\) and \(\text{Im}\, f\) are both continuous in \(\zeta\), then, as

\[
f = \text{Re}\, f + i \text{Im}\, f,
\]

\(f\) is continuous in \(\zeta\), once again, by Th. 7.38. \[\blacksquare\]

**Example 7.40.** (a) The continuity of the absolute value function \(z \mapsto |z|\) on \(\mathbb{K}\) can be concluded directly from (7.11e) and, alternatively, from combining the continuity of \(f : \mathbb{K} \to \mathbb{K}\), \(f(z) = z\), according to Ex. 7.32(b), with the continuity of \(|f|\) according to Th. 7.38.
(b) Every polynomial \( P : \mathbb{K} \rightarrow \mathbb{K}, \ P(x) = \sum_{j=0}^{n} a_j x^j, \ a_j \in \mathbb{K}, \) is continuous: First note that every monomial \( x \mapsto x^j \) is continuous on \( \mathbb{K} \) by (7.11g). Then Th. 7.38 implies the continuity of \( x \mapsto a_j x^j \) on \( \mathbb{K} \). Now the continuity of \( P \) follows from (7.16a) or, alternatively, by an induction from the \( f + g \) part of Th. 7.38.

(c) Let \( P, Q : \mathbb{K} \rightarrow \mathbb{K} \), be polynomials and let \( A := Q^{-1}\{0\} \) the set of all zeros of \( Q \) (if any). Then the rational function \((P/Q) : \mathbb{K} \setminus A \rightarrow \mathbb{K}\) is continuous as a consequence of (b) plus the \( f/g \) part of Th. 7.38.

**Theorem 7.41.** Let \( D_f, D_g \subseteq \mathbb{C}, \ f : D_f \rightarrow \mathbb{C}, \ g : D_g \rightarrow \mathbb{K}, \ f(D_f) \subseteq D_g \). If \( f \) is continuous in \( \zeta \in D_f \) and \( g \) is continuous in \( f(\zeta) \in D_g \), then \( g \circ f : D_f \rightarrow \mathbb{K} \) is continuous in \( \zeta \). In consequence, if \( f \) and \( g \) are both continuous, then the composition \( g \circ f \) is also continuous.

**Proof.** Let \( \zeta \in D_f \) and assume \( f \) is continuous in \( \zeta \) and \( g \) is continuous in \( f(\zeta) \). If \((z_n)_{n \in \mathbb{N}}\) is a sequence in \( D_f \) such that \( \lim_{n \rightarrow \infty} z_n = \zeta \), then the continuity of \( f \) in \( \zeta \) implies that \( \lim_{n \rightarrow \infty} f(z_n) = f(\zeta) \). Then the continuity of \( g \) in \( f(\zeta) \) implies \( \lim_{n \rightarrow \infty} g(f(z_n)) = g(f(\zeta)) \), thereby establishing the continuity of \( g \circ f \) in \( \zeta \). \( \square \)

### 7.2.3 Bounded, Closed, and Compact Sets

Subsets \( A \) of \( \mathbb{C} \) (and even subsets of \( \mathbb{R} \)) can be extremely complicated. If the set \( A \) has one or more of the benign properties defined in the following, then this can often be exploited in some useful way (we will see an important example in Th. 7.54 below).

**Definition 7.42.** Consider \( A \subseteq \mathbb{C} \).

(a) \( A \) is called **bounded** if, and only if, \( A = \emptyset \) or the set \( \{|z| : z \in A\} \) is bounded in \( \mathbb{R} \) in the sense of Def. 2.26(a), i.e. if, and only if,

\[
\exists M \in \mathbb{R}^+ \quad A \subseteq B_M(0).
\]

(b) \( A \) is called **closed** if, and only if, every sequence in \( A \) that converges in \( \mathbb{C} \) has its limit in \( A \) (note that \( \emptyset \) is, thus, closed).

(c) \( A \) is called **compact** if, and only if, \( A \) is both closed and bounded.

**Example 7.43.** (a) Clearly, \( \emptyset \) and sets containing single points \( \{z\}, \ z \in \mathbb{C} \) are compact. The sets \( \mathbb{C} \) and \( \mathbb{R} \) are simple examples of closed sets that are not bounded.

(b) Let \( a, b \in \mathbb{R}, \ a < b \). Each bounded interval \([a, b[, \ [a, b], \ ]a, b[, \ ]a, b] \) is, indeed, bounded (e.g. by \( M := 2 \max\{|a|, |b|\} \)). If \((x_n)_{n \in \mathbb{N}}\) is a sequence in \([a, b]\), converging to \( x \in \mathbb{R} \), then Th. 7.13(c) shows \( a \leq x \leq b \), i.e. \( x \in [a, b] \) and \([a, b] \) is, indeed, closed. Analogously, one sees that the unbounded intervals \([a, \infty[, \ ]-\infty, a] \) are also closed. On the other hand, open and half-open intervals are not closed: For sufficiently large \( n \), the convergent sequence \((b - \frac{1}{n})_{n \in \mathbb{N}}\) is in \([a, b]\), but \( \lim_{n \rightarrow \infty}(b - \frac{1}{n}) = b \notin [a, b[, \) and the other cases are treated analogously. In particular, only intervals of the form \([a, b]\) (and trivial intervals) are compact.
(c) For each \( \epsilon > 0 \) and each \( z \in \mathbb{C} \), the set \( B_\epsilon(z) \) is bounded (since \( B_\epsilon(z) \subseteq B_{\epsilon + |z|}(0) \) by the triangle inequality), but not closed (since, for sufficiently large \( n \in \mathbb{N} \), \( (z + \epsilon - \frac{1}{n})_{n \in \mathbb{N}} \) is a sequence in \( B_\epsilon(z) \), converging to \( z + \epsilon \notin B_\epsilon(z) \)). In particular, \( B_\epsilon(z) \) is not compact.

**Proposition 7.44. (a)** Finite unions of bounded (resp. closed, resp. compact) sets are bounded (resp. closed, resp. compact), i.e. if \( A_1, \ldots, A_n \subseteq \mathbb{C}, n \in \mathbb{N}, \) are bounded (resp. closed, resp. compact), then \( A := \bigcup_{j=1}^n A_j \) is also bounded (resp. closed, resp. compact).

**Proposition 7.44. (b)** Arbitrary (i.e. finite or infinite) intersections of bounded (resp. closed, resp. compact) sets are bounded (resp. closed, resp. compact), i.e. if \( I \neq \emptyset \) is an arbitrary index set and, for each \( j \in I \), \( A_j \subseteq \mathbb{C} \) is bounded (resp. closed, resp. compact), then \( A := \bigcap_{j \in I} A_j \) is also bounded (resp. closed, resp. compact).

**Proof. (a):** Exercise.

(b): Fix \( j_0 \in I \). If all \( A_j, j \in I, \) are bounded, then, in particular, there is \( M \in \mathbb{R}^+ \) such that \( A_{j_0} \subseteq B_M(0) \). Thus, \( A = \bigcap_{j \in I} A_j \subseteq A_{j_0} \subseteq B_M(0) \) shows \( A \) is also bounded. If all \( A_j, j \in I, \) are closed and \( (a_n)_{n \in \mathbb{N}} \) is a sequence in \( A \) that converges to some \( z \in \mathbb{C} \), then \( (a_n)_{n \in \mathbb{N}} \) is a sequence in each \( A_j, j \in I, \) and, since each \( A_j \) is closed, \( z \in A \) for each \( j \in I \), i.e. \( z \in A = \bigcap_{j \in I} A_j \). If all \( A_j, j \in I, \) are compact, then they are all closed and bounded and, thus, \( A \) is closed and bounded, i.e. \( A \) is compact. \( \blacksquare \)

**Example 7.45. (a)** According to Prop. 7.44(a), all finite subsets of \( \mathbb{C} \) are compact.

(b) \( \mathbb{N} = \bigcup_{n \in \mathbb{N}} \{n\} \) shows that infinite unions of compact sets can be unbounded, and \( ]0, 1[ = \bigcup_{n \in \mathbb{N}} [\frac{1}{n}, 1 - \frac{1}{n}] \) shows that infinite unions of compact sets are not always closed.

Many more examples of closed sets can be obtained as preimages of closed sets under continuous maps according to the following remark:

**Remark 7.46.** In Calculus II, it will be shown in the more general context of maps \( f \) between metric spaces that a map \( f \) is continuous if, and only if, all preimages \( f^{-1}(A) \) under \( f \) of closed sets \( A \) are closed. Here, we will only prove the following special case:

\[
f : \mathbb{C} \rightarrow \mathbb{K} \text{ continuous and } A \subseteq \mathbb{K} \text{ closed } \Rightarrow \ f^{-1}(A) \subseteq \mathbb{C} \text{ closed.} \tag{7.38}\]

Indeed, suppose \( f \) is continuous and \( A \subseteq \mathbb{K} \) is closed. If \( (z_n)_{n \in \mathbb{N}} \) is a sequence in \( f^{-1}(A) \) with \( \lim_{n \to \infty} z_n = z \in \mathbb{C} \), then \( (f(z_n))_{n \in \mathbb{N}} \) is a sequence in \( A \). The continuity of \( f \) then implies \( \lim_{n \to \infty} f(z_n) = f(z) \) and, then, \( f(z) \in A \), since \( A \) is closed. Thus, \( z \in f^{-1}(A) \), showing \( f^{-1}(A) \) is closed.

**Example 7.47. (a)** For each \( z \in \mathbb{C} \) and each \( r > 0 \), the closed disk \( \overline{B}_r(z) := \{w \in \mathbb{C} : |z - w| \leq r\} \) with radius \( r \) and center \( z \) is, indeed, closed by (7.38), since

\[
\overline{B}_r(z) = f^{-1}[0, r], \tag{7.39}\]

where \( f \) is the continuous map \( f : \mathbb{C} \rightarrow \mathbb{R}, f(w) := |z - w| \). Since \( \overline{B}_r(z) \) is clearly bounded, it is also compact.
(b) For each \( z \in \mathbb{C} \) and each \( r > 0 \), the circle (also called a 1-sphere) \( S_r(z) := \{ w \in \mathbb{C} : |z - w| = r \} \) with radius \( r \) and center \( z \) is closed by (7.38), since \( S_r(z) = f^{-1}\{r\} \), where \( f \) is the same map as in (7.39). Moreover, \( S_r(z) \) is also clearly bounded, and, thus, compact.

(c) According to (7.38), for each \( x \in \mathbb{R} \), the closed half-spaces \( \{ z \in \mathbb{C} : \text{Re} \ z \geq x \} = \text{Re}^{-1}[x, \infty[ \) and \( \{ z \in \mathbb{C} : \text{Im} \ z \geq x \} = \text{Im}^{-1}[x, \infty[ \) are, indeed, closed.

Theorem 7.48. A subset \( K \) of \( \mathbb{C} \) is compact if, and only if, every sequence in \( K \) has a subsequence that converges to some limit \( z \in K \).

Proof. If \( K \) is closed and bounded, and \( (z_n)_{n \in \mathbb{N}} \) is a sequence in \( K \), then the boundedness, the Bolzano-Weierstrass Th. 7.27, and Prop. 7.26 yield a subsequence that converges to some \( z \in \mathbb{C} \). However, since \( K \) is closed, \( z \in K \).

Conversely, assume every sequence in \( K \) has a subsequence that converges to some limit \( z \in K \). Let \( (z_n)_{n \in \mathbb{N}} \) be a sequence in \( K \) that converges to some \( \in \mathbb{C} \). Then this sequence must have a subsequence that converges to some \( z \in K \). However, according to Prop. 7.23, it must be \( w = z \in K \), showing \( K \) is closed. If \( K \) is not bounded, then there exists a sequence \( (z_n)_{n \in \mathbb{N}} \) in \( K \) such that \( \lim_{n \to \infty} |z_n| = \infty \). Every subsequence \( (z_{n_k})_{k \in \mathbb{N}} \) still has the property that \( \lim_{k \to \infty} |z_{n_k}| = \infty \), in particular, each subsequence is unbounded and can not converge to some \( z \in \mathbb{C} \) (let alone in \( K \)).

Caveat 7.49. In Calculus II, we will generalize the notion of compactness to subsets of so-called metric spaces, defining a set \( K \) to be compact if, and only if, every sequence in \( K \) has a subsequence that converges to some limit in \( K \). While it remains true that every compact set is closed and bounded, the converse does not(!) hold in general metric spaces (in general, even in closed sets, there exist bounded sequences that do not have convergent subsequences).

One reason that compact sets are useful is that real-valued continuous functions on compact sets assume a maximum and a minimum, which is the contents of Th. 7.54 below. In preparation, we now define maxima and minima for real-valued functions.

Definition 7.50. Let \( M \subseteq \mathbb{C} \), \( f : M \to \mathbb{R} \).

(a) Given \( z \in M \), \( f \) has a (strict) global min at \( z \) if, and only if, \( f(z) \leq f(w) \) \((f(z) < f(w))\) for each \( w \in M \setminus \{z\} \). Analogously, \( f \) has a (strict) global max at \( z \) if, and only if, \( f(z) \geq f(w) \) \((f(z) > f(w))\) for each \( w \in M \setminus \{z\} \). Moreover, \( f \) has a (strict) global extreme value at \( z \) if, and only if, \( f \) has a (strict) global min or a (strict) global max at \( z \).

(b) Given \( z \in M \), \( f \) has a (strict) local min at \( z \) if, and only if, there exists \( \epsilon > 0 \) such that \( f(z) \leq f(w) \) \((f(z) < f(w))\) for each \( w \in \{w \in M : |z - w| < \epsilon\} \setminus \{z\} \). Analogously, \( f \) has a (strict) local max at \( z \) if, and only if, there exists \( \epsilon > 0 \) such that \( f(z) \geq f(w) \) \((f(z) > f(w))\) for each \( w \in \{w \in M : |z - w| < \epsilon\} \setminus \{z\} \).
Moreover, \( f \) has a (strict) local extreme value at \( z \) if, and only if, \( f \) has a (strict) local min or a (strict) local max at \( z \).

**Remark 7.51.** In the context of Def. 7.50, it is immediate from the respective definitions that \( f \) has a (strict) global min at \( z \in M \) if, and only if, \(-f\) has a (strict) global max at \( z \). Moreover, the same holds if “global” is replaced by “local”. It is equally obvious that every (strict) global min/max is a (strict) local min/max.

**Theorem 7.52.** If \( K \subseteq \mathbb{C} \) is compact, and \( f : K \to \mathbb{C} \) is continuous, then \( f(K) \) is compact.

**Proof.** If \( (w_n)_{n \in \mathbb{N}} \) is a sequence in \( f(K) \), then, for each \( n \in \mathbb{N} \), there is some \( z_n \in K \) such that \( f(z_n) = w_n \). As \( K \) is compact, there is a subsequence \( (a_n)_{n \in \mathbb{N}} \) of \( (z_n)_{n \in \mathbb{N}} \) with \( \lim_{n \to \infty} a_n = a \) for some \( a \in K \). Then \( (f(a_n))_{n \in \mathbb{N}} \) is a subsequence of \( (w_n)_{n \in \mathbb{N}} \) and the continuity of \( f \) yields \( \lim_{n \to \infty} f(a_n) = f(a) \in f(K) \), showing that \( (w_n)_{n \in \mathbb{N}} \) has a convergent subsequence with limit in \( f(K) \). By Th. 7.48, we have therefore established that \( f(K) \) is compact. \( \blacksquare \)

**Lemma 7.53.** If \( K \) is a nonempty compact subset of \( \mathbb{R} \), then \( K \) contains a smallest and a largest element, i.e. there exist \( m,M \in K \) such that \( m \leq x \leq M \) for each \( x \in K \).

**Proof.** Since the compact set \( K \) is bounded, we know that

\[-\infty < m := \inf K \leq \sup K : = M < \infty.\]

According to the definition of the inf and sup as largest lower bound and smallest upper bound, respectively, for each \( n \in \mathbb{N} \), there must be elements \( x_n, y_n \in K \) such that \( m \leq x_n \leq m + \frac{1}{n} \) and \( M - \frac{1}{n} \leq y_n \leq M \). Since the compact set \( K \) is also closed, we get

\[m = \lim_{n \to \infty} x_n \in K \text{ and } M = \lim_{n \to \infty} y_n \in K.\]

**Theorem 7.54.** If \( K \subseteq \mathbb{C} \) is compact, and \( f : K \to \mathbb{R} \) is continuous, then \( f \) assumes its max and its min, i.e. there are \( z_m \in K \) and \( z_M \in K \) such that \( f \) has a global min at \( z_m \) and a global max at \( z_M \). In particular, the continuous function \( f \) assumes its max and min on each compact interval \( K = [a,b] \subseteq \mathbb{R}, \ a,b \in \mathbb{R} \).

**Proof.** Since \( K \) is compact and \( f \) is continuous, \( f(K) \subseteq \mathbb{R} \) is compact according to Th. 7.52. Then, by Lem. 7.53, \( f(K) \) contains a smallest element \( m \) and a largest element \( M \). This, in turn, implies that there are \( z_m, z_M \in K \) such that \( f(z_m) = m \) and \( f(z_M) = M \). \( \blacksquare \)

**Example 7.55.** On an unbounded set, a continuous function does not necessarily have a global max or a global min, as one can already see from \( x \mapsto x \). An example for a continuous function on a bounded, but not closed, interval, that does not have a global max is \( f : ]0,1] \to \mathbb{R}, \ f(x) := 1/x \), which is continuous by Th. 7.38.
7.2.4 Intermediate Value Theorem

**Theorem 7.56** (Bolzano’s Theorem). Let \( a, b \in \mathbb{R} \) with \( a < b \). If \( f : [a, b] \rightarrow \mathbb{R} \) is continuous with \( f(a) > 0 \) and \( f(b) < 0 \), then \( f \) has at least one zero in \([a, b] \). More precisely, the set \( A := f^{-1}\{0\} \) has a min \( \xi_1 \) and a max \( \xi_2 \), \( a < \xi_1 \leq \xi_2 < b \), where \( f > 0 \) on \([a, \xi_1]\) and \( f < 0 \) on \([\xi_2, b]\).

**Proof.** Let \( \xi_1 := \inf f^{-1}(\mathbb{R}_0^+) \).

(a): \( f(\xi_1) \leq 0 \): This is clear if \( \xi_1 = b \). If \( \xi_1 < b \), then, for each \( n \in \mathbb{N} \) sufficiently large, there exists \( x_n \in [\xi_1, \xi_1 + 1/n] \subseteq [a, b] \) such that \( f(x_n) \leq 0 \). Then \( \lim_{n \to \infty} x_n = \xi_1 \) and the continuity of \( f \) implies \( \lim_{n \to \infty} f(x_n) = f(\xi_1) \). Now \( f(\xi_1) \leq 0 \) is a consequence of Th. 7.13(c). In particular, (a) yields \( a < \xi_1 \) and \( f > 0 \) on \([a, \xi_1]\).

(b): \( f(\xi_1) \geq 0 \): The continuity of \( f \) implies \( \lim_{n \to \infty} f(\xi_1 - 1/n) = f(\xi_1) \) and, since we have already seen \( f(\xi_1 - 1/n) > 0 \) for each \( n \in \mathbb{N} \) sufficiently large, \( f(\xi_1) \geq 0 \) is again a consequence of Th. 7.13(c). In particular, we have \( \xi_1 < b \).

Combining (a) and (b), we have \( f(\xi_1) = 0 \) and \( a < \xi_1 < b \).

Defining \( \xi_2 := \sup f^{-1}(\mathbb{R}_0^+) \), \( f(\xi_2) = 0 \) and \( a < \xi_2 < b \) is shown completely analogous. Then \( f < 0 \) on \([\xi_2, b]\) is also clear as well as \( \xi_1 \leq \xi_2 \). ■

**Theorem 7.57** (Intermediate Value Theorem). Let \( a, b \in \mathbb{R} \) with \( a < b \). If \( f : [a, b] \rightarrow \mathbb{R} \) is continuous, then \( f \) assumes every value between \( f(a) \) and \( f(b) \), i.e.

\[
\left[ \min\{f(a), f(b)\}, \max\{f(a), f(b)\} \right] \subseteq f([a, b]). \tag{7.40}
\]

**Proof.** If \( f(a) = f(b) \), then there is nothing to prove. If \( f(a) < f(b) \) and \( \eta \in [f(a), f(b)] \), then consider the auxiliary function \( g : [a, b] \rightarrow \mathbb{R}, \ g(x) := \eta - f(x) \). Then \( g \) is continuous with \( g(a) = \eta - f(a) > 0 \) and \( g(b) = \eta - f(b) < 0 \). According to Bolzano’s Th. 7.56, there exists \( \xi \in [a, b] \) such that \( g(\xi) = \eta - f(\xi) = 0 \), i.e. \( f(\xi) = \eta \) as claimed. If \( f(b) < f(a) \) and \( \eta \in [f(b), f(a)] \), then consider the auxiliary function \( g : [a, b] \rightarrow \mathbb{R}, \ g(x) := f(x) - \eta \). Then \( g \) is continuous with \( g(a) = f(a) - \eta > 0 \) and \( g(b) = f(b) - \eta < 0 \). Once again, according to Bolzano’s Th. 7.56, there exists \( \xi \in [a, b] \) such that \( g(\xi) = f(\xi) - \eta = 0 \), i.e. \( f(\xi) = \eta \). ■

**Theorem 7.58.** If \( I \subseteq \mathbb{R} \) is an interval (of one of the 8 types listed in (4.11)) and \( f : I \rightarrow \mathbb{R} \) is continuous, then \( f(I) \) is also an interval (it can degenerate to a single point if \( f \) is constant). More precisely, if \( \emptyset \neq I = [a, b] \) is a compact interval, then \( \emptyset \neq f(I) = [\min f(I), \max f(I)] \); if \( I \) is not a compact interval, then one of the following 9 cases occurs:

\[
f(I) = \mathbb{R}, \tag{7.41a}
\]
\[
f(I) = ]-\infty, \sup f(I)[, \tag{7.41b}
\]
\[
f(I) = ]-\infty, \sup f(I)[, \tag{7.41c}
\]
\[
f(I) = [\inf f(I), \infty[ \tag{7.41d}
\]
\[
f(I) = [\inf f(I), \sup f(I)], \tag{7.41e}
\]

on the bounded half-open interval $I$:

- $f(I) = [\inf f(I), \sup f(I)]$, \hspace{100pt} (7.41f)
- $f(I) = \inf f(I), \infty]$, \hspace{100pt} (7.41g)
- $f(I) = \inf f(I), \sup f(I)]$, \hspace{100pt} (7.41h)
- $f(I) = \inf f(I), \sup f(I)]$. \hspace{100pt} (7.41i)

**Proof.** If $I$ is a compact interval, then we merely combine Th. 7.54 with Th. 7.57. Otherwise, let $\eta \in f(I)$. If $f(I)$ has an upper bound, then Th. 7.57 implies $[\eta, \sup f(I)] \subseteq f(I)$ and $f(I) \cap [\eta, \infty] \subseteq [\eta, \sup f(I)]$. If $f(I)$ does not have an upper bound, then Th. 7.57 implies $f(I) \cap [\eta, \infty[ = [\eta, \infty[$. Analogously, one obtains $f(I) \cap ]-\infty, \eta] = ]-\infty, \eta]$ or $f(I) \cap ]-\infty, \eta] = [\inf f(I), \eta]$ or $f(I) \cap ]-\infty, \eta] = \inf f(I), \eta]$, showing that there are precisely the 9 possibilities of (7.41) for $f(I) = (f(I) \cap ]-\infty, \eta]) \cup (f(I) \cap [\eta, \infty[)$. \hspace{100pt} ■

The above results will have striking consequences in the following Sec. 7.2.5.

**Example 7.59.** The piecewise affine function

$$f : [0, 1] \rightarrow \mathbb{R}, \quad f(x) := \begin{cases} (-1)^n \cdot n - \frac{2n+1}{n-1} \cdot \left( x - \frac{1}{n} \right) & \text{for } x \in \left[ \frac{1}{n}, \frac{1}{n-1} \right], \ n \text{ even,} \\ (-1)^n \cdot n + \frac{2n+1}{n-1} \cdot \left( x - \frac{1}{n} \right) & \text{for } x \in \left[ \frac{1}{n}, \frac{1}{n-1} \right], \ n \geq 3 \text{ odd,} \end{cases}$$

satisfies $f(1/n) = (-1)^n \cdot n$ for each $n \in \mathbb{N}$ and is an example of a continuous function on the bounded half-open interval $I := ]0, 1]$ with $f(I) = \mathbb{R}$.

### 7.2.5 Inverse Functions, Existence of Roots, Exponential Function, Logarithm

**Theorem 7.60.** Let $I \subseteq \mathbb{R}$ be an interval (of one of the 8 types listed in (4.11)). If $f : I \rightarrow \mathbb{R}$ is strictly increasing (resp. decreasing), then $f$ has an inverse function $f^{-1}$ defined on $J := f(I)$, i.e. $f^{-1} : J \rightarrow I$, and $f^{-1}$ is continuous and strictly increasing (resp. decreasing). If $f$ is also continuous, then $J$ must be an interval.

**Proof.** From Prop. 2.31(b), we know $f : I \rightarrow \mathbb{R}$ is one-to-one. Then $f : I \rightarrow f(I)$ is invertible and Prop. 2.31(c) shows $f^{-1}$ is strictly monotone in the same sense as $f$. We need to prove the continuity of $f^{-1}$. We assume $f$ to be strictly increasing (the case where $f$ is strictly decreasing then follows by considering $-f$). Let $\eta \in J$, $\varepsilon > 0$, and $\xi \in I$ with $f(\xi) = \eta$. If $J = \{\xi\}$, then $J = \{\eta\}$, and there is nothing to prove. It remains to consider three cases:

(a) $\xi = \min I$, i.e. $\xi$ is the left endpoint of $I$ (and $\xi \neq \max I$),

(b) $\xi = \max I$, i.e. $\xi$ is the right endpoint of $I$ (and $\xi \neq \min I$),

(c) $\xi$ is neither the min nor the max of $I$.

We carry out the proof for (c) and leave the (very similar) cases (a) and (b) as exercises. In case (c), there are points $\xi_1, \xi_2 \in B_\varepsilon(\xi) \cap I$ such that

$$\xi - \varepsilon < \xi_1 < \xi < \xi_2 < \xi + \varepsilon.$$  \hspace{100pt} (7.42)
As $f$ is strictly increasing, this implies
\[ f(\xi_1) < \eta < f(\xi_2). \]
Choose $\delta > 0$ such that
\[ f(\xi_1) < \eta - \delta < \eta < \eta + \delta < f(\xi_2). \]
Then
\[ \forall y \in J \cap B_\delta(\eta) \left( f(\xi_1) < y < f(\xi_2) \Rightarrow f^{-1}(y) < \xi_1 < f^{-1}(y) < \xi_2 \right). \]
proving the continuity of $f^{-1}$ in $\eta$. That $J$ must be an interval if $f$ is continuous was already shown in Th. 7.58.

**Remark and Definition 7.61 (Roots).** We are now in a position to fulfill the promise made in Def. and Rem. 5.8, i.e. to prove the existence of unique roots for nonnegative real numbers: For each $n \in \mathbb{N}$, the function $f : \mathbb{R}_0^+ \to \mathbb{R}$, $f(x) := x^n$, is continuous and strictly increasing with $J := f(\mathbb{R}_0^+) = \mathbb{R}_0^+$. Then Th. 7.60 implies the existence of a continuous and strictly increasing inverse function $f^{-1} : \mathbb{R}_0^+ \to \mathbb{R}_0^+$. For each $x \in \mathbb{R}_0^+$, we call $f^{-1}(x)$ the $n$th root of $x$ and write $\sqrt[n]{x} := x^{\frac{1}{n}} := f^{-1}(x)$. Then $(\sqrt[n]{x})^n = (x^{\frac{1}{n}})^n = x$ is immediate from the definition. Caveat: By definition, roots are always nonnegative and they are only defined for nonnegative numbers (when studying complex numbers and $\mathbb{C}$-valued functions more deeply in the field of Complex Analysis, one typically extends the notion of root, but we will not pursue this route in this class). As anticipated in Def. and Rem. 5.8, one also writes $\sqrt[n]{x}$ instead of $\sqrt[n]{x}$ and calls $\sqrt[n]{x}$ the square root of $x$.

**Remark and Definition 7.62.** It turns out that $\sqrt{2}$ (and many other roots) are not rational numbers, i.e. $\sqrt{2} \notin \mathbb{Q}$. This is easily proved by contradiction: If $\sqrt{2} \in \mathbb{Q}$, then there exist natural numbers $m, n \in \mathbb{N}$ such that $\sqrt{2} = \frac{m}{n}$. Moreover, by canceling possible factors of 2, we may assume at least one of the numbers $m, n$ is odd. Now $\sqrt{2} = \frac{m}{n}$ implies $m^2 = 2n^2$, i.e. $m^2$ and, thus, $m$ must be even. In consequence, there exists $p \in \mathbb{N}$ such that $m = 2p$, implying $2n^2 = m^2 = 4p^2$ and $n^2 = 2p^2$. Thus $n^2$ and $n$ must also be even, in contradiction to $m, n$ not both being even.

The elements of $\mathbb{R} \setminus \mathbb{Q}$ are called irrational numbers. It turns out that most real numbers are irrational numbers – one can show that $\mathbb{Q}$ is countable, whereas $\mathbb{R} \setminus \mathbb{Q}$ is not countable (actually, every interval contains countably many rational and uncountably many irrational numbers, see [Phi16, Sec. F], in particular, [Phi16, Th. F.1(c)] and [Phi16, Cor. F.4]).

**Theorem 7.63 (Inequality Between the Arithmetic Mean and the Geometric Mean).** If $n \in \mathbb{N}$ and $x_1, \ldots, x_n \in \mathbb{R}_0^+$, then
\[ \sqrt[n]{x_1 \cdots x_n} \leq \frac{x_1 + \cdots + x_n}{n}, \quad (7.43) \]
where the left-hand side is called the geometric mean and the right-hand side is called the arithmetic mean of the numbers $x_1, \ldots, x_n$. Equality occurs if, and only if, $x_1 = \cdots = x_n$. 
Proof. If at least one of the $x_j$ is 0, then (7.43) becomes the true statement $0 \leq \frac{x_1 + \cdots + x_n}{n}$ with strict equality if at least one $x_j > 0$. If $x_1 = \cdots = x_n = x$, then (7.43) also holds since both sides are equal to $x$. Thus, for the remainder of the proof, we assume all $x_j > 0$ and not all $x_j$ are equal. First, we consider the special case, where $\frac{x_1 + \cdots + x_n}{n} = 1$.

Since not all $x_j$ are equal, there exists $k$ with $x_k \neq 1$. We prove (7.43) by induction for $n \in \{2, 3, \ldots\}$ in the form

$$\left(\sum_{j=1}^{n} x_j = n \land \exists k \in \{1, \ldots, n\}, x_k \neq 1\right) \Rightarrow \prod_{j=1}^{n} x_j < 1.$$

Base Case ($n = 2$): Since $x_1 + x_2 = 2$, $0 < x_1, x_2$ and not both $x_1$ and $x_2$ are equal to 1, there is $\epsilon > 0$ such that $x_1 = 1 + \epsilon$ and $x_2 = 1 - \epsilon$, i.e. $x_1x_2 = 1 - \epsilon^2 < 1$, which establishes the base case. Induction Step: We now have $n \geq 2$ and $0 < x_1, \ldots, x_{n+1}$ with $\sum_{j=1}^{n+1} x_j = n + 1$ plus the existence of $k, l \in \{1, \ldots, n + 1\}$ such that $x_k = 1 + \alpha$, $x_l = 1 - \beta$ with $\alpha, \beta > 0$. Then define $y := x_k + x_l - 1 = 1 + \alpha - \beta$. One observes $y > 0$ (since $\beta < 1$) and

$$y + \sum_{j=1, j \neq k,l}^{n+1} x_j = -1 + \sum_{j=1}^{n+1} x_j \quad \text{ind. hyp.} \quad y \prod_{j=1, j \neq k,l}^{n+1} x_j \leq 1,$$

(we cannot exclude equality as $y$ and all the remaining $x_j$ might be equal to 1). Since $x_kx_l = (1 + \alpha)(1 - \beta) = 1 + \alpha - \beta - \alpha\beta = y - \alpha\beta < y$, we now infer $\prod_{j=1}^{n+1} x_j < 1$, concluding the induction proof. It remains to consider the case $\frac{x_1 + \cdots + x_n}{n} = \lambda > 0$, not all $x_j$ equal. One estimates

$$\sqrt[n]{x_1 \cdots x_n} = \lambda \sqrt[n]{\frac{x_1}{\lambda} \cdots \frac{x_n}{\lambda}} \quad \text{special case} \quad \lambda \frac{x_1 + \cdots + x_n}{\lambda n} = \frac{x_1 + \cdots + x_n}{n},$$

completing the proof of the theorem. \hfill \blacksquare

Corollary 7.64. For each $a \in \mathbb{R}_0^+ \setminus \{1\}$, $n \in \{2, 3, \ldots\}$, $p \in \{1, \ldots, n - 1\}$:

$$\sqrt[n]{a^p} < 1 + \frac{p}{n}(a - 1); \quad p = 1 \quad \text{yields} \quad \sqrt[n]{a} < 1 + \frac{a - 1}{n}. \quad (7.44)$$

Proof. The simple application

$$\sqrt[n]{a^p} = \sqrt[n]{a^p \prod_{j=1}^{n-p} 1} \quad \text{Th. 7.63} \quad \frac{pa + n - p}{n} = 1 + \frac{p}{n}(a - 1) \quad (7.45)$$

of Th. 7.63 establishes the case. \hfill \blacksquare

Example 7.65. We use (7.43) to show

$$\lim_{n \to \infty} \sqrt[n]{n} = 1: \quad (7.46)$$
First note $0 < x < 1 \Rightarrow 0 < x^n < 1$, i.e. $\sqrt[n]{n} > 1$ for each $(\sqrt[n]{n})^n = n > 1$. Now write $n$ as the product of $n$ factors $n = \sqrt[n]{\sqrt[n]{n} \cdot \Pi_{k=1}^{n-2} 1}$. Then, for $n > 1$,

$$\sqrt[n]{n} = \sqrt[n]{\sqrt[n]{n} \cdot \Pi_{k=1}^{n-2} 1} \overset{\text{Th. 7.63}}{\leq} \frac{2\sqrt{n} + n - 2}{n} < 1 + \frac{2}{\sqrt{n}}. \quad (7.47)$$

It is an exercise to show

$$\lim_{n \to \infty } \frac{1}{\sqrt{n}} = 0. \quad (7.48)$$

Now this together with $1 \leq \sqrt[n]{n} \leq 1 + \frac{2}{\sqrt{n}}$ and the Sandwich Th. 7.16 proves (7.46).

**Example 7.66** (Euler’s Number). We use Th. 7.63 to prove the limit

$$e := \lim_{n \to \infty } \left( 1 + \frac{1}{n} \right)^n \quad (7.49)$$

exists. It is known as Euler’s number. One can show it is an irrational number (see [Phi16, Sec. H.1]) and its first digits are $e = 2.71828\ldots$ It is of exceptional importance for analysis and mathematics in general, as it pops up in all kinds of different mathematical contexts. From Th. 7.63, we obtain

$$\forall n \in \mathbb{N} \quad \forall x \in [-n,1], \quad \left( 1 + \frac{x}{n} \right)^n = 1 \cdot \left( 1 + \frac{x}{n} \right)^n < \left( 1 + \frac{x}{n+1} \right)^{n+1}, \quad (7.50)$$

where we have used that, on both sides of the inequality in (7.50), there are $n+1$ factors having the same sum, namely $n + 1 + x$; and the inequality in (7.43) is strict, unless all factors are equal. We now apply (7.50) to the sequences $(a_n)_{n \in \mathbb{N}}, (b_n)_{n \in \mathbb{N}}, (c_n)_{n \in \mathbb{N}}$, where

$$\forall n \in \mathbb{N} \quad \left( a_n := \left( 1 + \frac{1}{n} \right)^n, \quad b_n := \left( 1 - \frac{1}{n} \right)^n, \quad c_n := \frac{1}{n+1} \right) \quad (7.51)$$

Applying (7.50) with $x = 1$ and $x = -1$, respectively, yields $(a_n)_{n \in \mathbb{N}}$ and $(b_n)_{n \in \mathbb{N}}$ are strictly increasing, and $(c_n)_{n \in \mathbb{N}}$ is strictly decreasing. On the other hand, $a_n < c_n$ holds for each $n \in \mathbb{N}$, showing $(a_n)_{n \in \mathbb{N}}$ is bounded from above by $c_1$, and $(c_n)_{n \in \mathbb{N}}$ is bounded from below by $a_1$. In particular, Th. 7.19 implies the convergence of both $(a_n)_{n \in \mathbb{N}}$ and $(c_n)_{n \in \mathbb{N}}$. Moreover, $\lim_{n \to \infty } c_n = \lim_{n \to \infty } \left( a_n (1+1/n) \right) = e \cdot 1 = e$, which, together with $a_n < e < c_n$ for each $n \in \mathbb{N}$, can be used to compute $e$ to an arbitrary precision.

**Definition 7.67.** Let $A \subseteq \mathbb{R}$ be a subset of the real numbers. Then $A$ is called dense in $\mathbb{R}$ if, and only if, every $\epsilon$-neighborhood of every real number contains a point from $A$, i.e. if, and only if,

$$\forall x \in \mathbb{R} \quad \forall \epsilon \in \mathbb{R}^+, \quad A \cap B_\epsilon(x) \neq \emptyset.$$

**Theorem 7.68.** (a) $\mathbb{Q}$ is dense in $\mathbb{R}$. 

(b) \( \mathbb{R} \setminus \mathbb{Q} \) is dense in \( \mathbb{R} \).

(c) For each \( x \in \mathbb{R} \), there exist sequences \((r_n)_{n \in \mathbb{N}}\) and \((s_n)_{n \in \mathbb{N}}\) in the rational numbers \( \mathbb{Q} \) such that \( x = \lim_{n \to \infty} r_n = \lim_{n \to \infty} s_n \), \((r_n)_{n \in \mathbb{N}}\) is strictly increasing and \((s_n)_{n \in \mathbb{N}}\) is strictly decreasing.

Proof. (a): Since each \( B_\epsilon(x) \) is an interval, it suffices to prove that every interval \( [a,b) \), \( a < b \), contains a rational number. If \( 0 \in [a,b) \), then there is nothing to prove. Suppose \( 0 < a < b \) and set \( \delta := b - a > 0 \). Choose \( n \in \mathbb{N} \) such that \( 1/n < \delta \) and let

\[
q := \max \left\{ \frac{k}{n} : k \in \mathbb{N} \land \frac{k}{n} < b \right\}.
\]

Then \( q \in \mathbb{Q} \) and \( a < q < b \). If \( a < b < 0 \), choose \( \delta \) and \( n \) as above, but let

\[
q := \min \left\{ -\frac{k}{n} : k \in \mathbb{N} \land -\frac{k}{n} > a \right\}.
\]

Then, once again, \( q \in \mathbb{Q} \) and \( a < q < b \).

(b): Analogous to (a), we show that every interval \( [a,b) \), \( a < b \), contains an irrational number: According to (a), we choose \( q \in \mathbb{Q} \cap [a,b) \), \( \delta := b - q > 0 \) and \( n \in \mathbb{N} \) such that \( \sqrt{2}/n < \delta \). Then \( a < \lambda := q + \sqrt{2}/n < b \) and also \( \lambda \in \mathbb{R} \setminus \mathbb{Q} \) (otherwise, \( \sqrt{2} = n(\lambda - q) \in \mathbb{Q} \)).

(c): Using (a), for each \( n \in \mathbb{N} \), we choose rational numbers \( r_n \) and \( s_n \) such that

\[
r_n \in \left[ x - \frac{1}{n}, x - \frac{1}{n+1} \right], \quad s_n \in \left[ x + \frac{1}{n+1}, x + \frac{1}{n} \right].
\]

Then, clearly, \((r_n)_{n \in \mathbb{N}}\) is strictly increasing, \((s_n)_{n \in \mathbb{N}}\) is strictly decreasing, and the Sandwich Th. 7.16 implies \( x = \lim_{n \to \infty} r_n = \lim_{n \to \infty} s_n \).

\[\square\]

**Definition and Remark 7.69** (Exponentiation). In Not. 5.6, we had defined \( a^x \) for \((a,x) \in \mathbb{C} \times \mathbb{N}_0\) and for \((a,x) \in (\mathbb{C} \setminus \{0\}) \times \mathbb{Z}\). We will now extend the definition to \((a,x) \in \mathbb{R}^+ \times \mathbb{R}\) (later, we will further extend the definition to \((a,z) \in \mathbb{R}^+ \times \mathbb{C}\)). The present extension to \((a,x) \in \mathbb{R}^+ \times \mathbb{R}\) is accomplished in two steps – first, in (a), for rational \( x \), then, in (b), for irrational \( x \).

(a) For rational \( x = k/n \) with \( k \in \mathbb{Z} \) and \( n \in \mathbb{N} \), define

\[
a^x := a^{k/n} := \sqrt[n]{a^k}.
\]

(7.52)

For this definition to make sense, we have to check it does not depend on the special representation of \( x \), i.e., we have to verify \( x = \frac{k}{n} = \frac{km}{nm} \) with \( k \in \mathbb{Z} \) and \( m,n \in \mathbb{N} \) implies \( a^{\frac{k}{n}} = a^{\frac{km}{nm}} \). To this end, observe, using Rem. and Def. 7.61,

\[
(a^{k/n})^{nm} = (\sqrt[n]{a^k})^{nm} = a^{km} \quad \text{and} \quad (a^{km/nm})^{nm} = (\sqrt[n]{a^{km}})^{nm} = a^{km}.
\]

(7.53)
proving $a^\frac{k}{a} = a^{\frac{km}{a}}$ (here, as in Rem. and Def. 7.61, we used that $\lambda \mapsto \lambda^N$ is one-to-one on $\mathbb{R}^+_0$ for each $N \in \mathbb{N}$). The exponentiation rules of Th. 5.7 now extend to rational exponents in a natural way, i.e., for each $a, b > 0$ and each $x, y \in \mathbb{Q}$:

\[
\begin{align*}
\alpha^x + y &= \alpha^x \alpha^y, \\
\alpha^b \alpha^x &= (\alpha \alpha)^x, \\
(\alpha^x)^y &= \alpha^{xy}.
\end{align*}
\]

(7.54a)

For the proof, by possibly multiplying numerator and denominator by some natural number, we can assume $x = k/n$ and $y = l/n$ with $k, l \in \mathbb{Z}$ and $n \in \mathbb{N}$. Then

\[
(a^x + y)^n = (a^{\frac{k+l}{n}})^n = a^{k+l}
\]

proving (7.54a);

\[
(a^b \alpha^x)^n = (a^{\frac{k}{n}})^n (b^{\frac{k}{n}})^n a^k b^k = (ab)^k
\]

proving (7.54b);

\[
(a^x)^{k+l} = (a^{\frac{k}{n}})^{k+l} (a^{\frac{l}{n}})^{k+l} = (ab)^{k+l}
\]

proving (7.54c).

Moreover, we obtain the following monotonicity rules for each $a, b \in \mathbb{R}^+$ and each $x, y \in \mathbb{Q}$:

\[
\begin{align*}
\forall \ x > 0 \quad (a < b) & \Rightarrow (\alpha^x < \alpha^y), \\
\forall \ x < 0 \quad (a < b) & \Rightarrow (\alpha^x > \alpha^y), \\
\forall \ a > 1 \quad (x < y) & \Rightarrow (\alpha^x < \alpha^y), \\
\forall \ 0 < a < 1 \quad (x < y) & \Rightarrow (\alpha^x > \alpha^y).
\end{align*}
\]

(7.55a)

(7.55b)

(7.55c)

(7.55d)

If $x = k/n$ with $k, n \in \mathbb{N}$ and $a < b$, then $a^{1/n} < b^{1/n}$ according to Rem. and Def. 7.61, which, in turn, implies $\alpha^x = (a^{1/n})^k < (b^{1/n})^k = b^x$, proving (7.55a); and $a^{-1} > b^{-1}$ implies $\alpha^{-x} = (a^{-1})^x > (b^{-1})^x = \alpha^{-x}$, proving (7.55b). If $x < y$, set $q := y - x > 0$. Then $1 < a$ and (7.55a) imply $1^q < a^q$, i.e. $\alpha^x < \alpha^y \alpha^q = \alpha^y$, proving (7.55c). Similarly, $0 < a < 1$ and (7.55a) imply $a^q < 1^q = 1$, i.e. $\alpha^y = \alpha^x a^q < \alpha^y$, proving (7.55d).

The following estimates will also come in handy: For $a \in \mathbb{R}^+$ and $x, y \in \mathbb{Q}$:

\[
\forall \ a > 1 \wedge x > 0 \quad \Rightarrow (\alpha^x - 1 < x \cdot \alpha^{x+1}),
\]

(7.56)

\[
\forall \ m \in \mathbb{N} \quad (x, y \in [-m, m] \quad \Rightarrow |\alpha^x - \alpha^y| \leq L |x - y|,
\]

(7.57)

where $L := \max\{a^{m+1}, (1/a)^{m+1}\}$. 


For $x \geq 1$, (7.56) is proved by $a^x < a^{x+1} < x \cdot a^{x+1} + 1$; for $x < 1$, write $x = p/n$ with $p, n \in \mathbb{N}$ and $p < n$, and apply (7.44) to obtain $a^x < 1 + x(a - 1) < 1 + xa < 1 + x \cdot a^{x+1}$. For the proof of (7.57), first consider $a > 1$. Moreover, by possibly renaming $x$ and $y$, we may assume $x < y$, i.e. $z := y - x > 0$. Thus, (7.56) holds with $x$ replaced by $z$. Multiplying the resulting inequality by $a^x$ yields

$$a^x a^z - a^x = a^y - a^x < z \cdot a^x a^{x+1} = (y - x) a^{x+1} \leq (y - x) a^{m+1},$$

proving (7.57) for $a > 1$. For $a = 1$, it is clearly true, and for $a < 1$, it is $a^{-1} > 1$, i.e.

$$|a^x - a^y| = |(a^{-1})^{-x} - (a^{-1})^{-y}| \leq |y - x| (a^{-1})^{m+1},$$

finishing the proof of (7.57).

(b) We now define $a^x$ for irrational $x$ by letting

$$a^x := \lim_{n \to \infty} a^{q_n}, \quad \text{where } (q_n)_{n \in \mathbb{N}} \text{ is a sequence in } \mathbb{Q} \text{ with } \lim_{n \to \infty} q_n = x. \quad (7.58)$$

For this definition to make sense, we have to know such sequences $(q_n)_{n \in \mathbb{N}}$ exist, which we do know from Th. 7.68(c). We also know from Th. 7.68(c) that there exists an increasing sequence $(q_n)_{n \in \mathbb{N}}$ in $\mathbb{Q}$ converging to $x$, in particular, bounded by $x$. Then, by (7.55c) and (7.55d), respectively, $(a^{q_n})_{n \in \mathbb{N}}$ is increasing for $a > 1$ and decreasing for $0 < a < 1$. Moreover, the sequence is bounded from above by $a^N$ with $N \in \mathbb{N}$, $N > x$, for $a > 1$; and bounded from below by 0 for $0 < a < 1$. In both cases, Th. 7.19 implies convergence of the sequence to some limit that we may call $a^x$. However, we still need to verify that, for each sequence $(r_n)_{n \in \mathbb{N}}$ in $\mathbb{Q}$ with $\lim_{n \to \infty} r_n = x$, the sequence $(a^{r_n})_{n \in \mathbb{N}}$ converges to the same limit $a^x$ in $\mathbb{R}$. If $\lim_{n \to \infty} r_n = x$, then $\lim_{n \to \infty} |q_n - r_n| = 0$. Since $(r_n)_{n \in \mathbb{N}}$ and $(q_n)_{n \in \mathbb{N}}$ are bounded, (7.57) implies

$$\exists L \in \mathbb{R}^+ \quad \forall n \in \mathbb{N} \quad |a^{q_n} - a^{r_n}| \leq L |q_n - r_n|, \quad (7.59)$$

such that Prop. 7.11(a) implies $\lim_{n \to \infty} |a^{q_n} - a^{r_n}| = 0$ and

$$\lim_{n \to \infty} a^{r_n} = \lim_{n \to \infty} (a^{r_n} - a^{q_n} + a^{q_n}) = 0 + a^x = a^x, \quad (7.60)$$

showing (7.58) does not depend on the chosen sequence.

**Proposition 7.70.** The exponentiation rules (7.54), the monotonicity rules (7.55), and the estimates (7.56) and (7.57) remain valid if $x, y \in \mathbb{Q}$ is replaced by $x, y \in \mathbb{R}$. Moreover, for each $a > 0$ and each sequence $(x_n)_{n \in \mathbb{N}}$ in $\mathbb{R}$:

$$\lim_{n \to \infty} x_n = x \in \mathbb{R} \quad \Rightarrow \quad \lim_{n \to \infty} a^{x_n} = a^x. \quad (7.61)$$

**Proof.** Given $x, y \in \mathbb{R}$, let $(p_n)_{n \in \mathbb{N}}$ and $(q_n)_{n \in \mathbb{N}}$ be sequences in $\mathbb{Q}$ such that $\lim_{n \to \infty} p_n = x$ and $\lim_{n \to \infty} q_n = y$.

We start by verifying (7.57). As we can assume $(p_n)_{n \in \mathbb{N}}$ and $(q_n)_{n \in \mathbb{N}}$ to be monotone, we may also assume $p_n, q_n \in [-m, m]$ for each $n \in \mathbb{N}$. Then the rational case of (7.57) implies

$$\forall n \in \mathbb{N} \quad |a^{p_n} - a^{q_n}| \leq L |p_n - q_n|,$$
and Th. 7.13(c) establishes the case. Then (7.61) also follows, since
\[ 0 \leq |x_n - a^x| \leq L |x_n - x| \to 0. \]

We deal with (7.54) next. For each \( a, b > 0 \):
\[
\begin{align*}
a^{x+y} & = \lim_{n \to \infty} a^{p_n + q_n} \quad \text{(7.54a)} \quad \lim_{n \to \infty} (a^{p_n} a^{q_n}) = a^x a^y, \\
a^x b^y & = \lim_{n \to \infty} a^{p_n} \lim_{n \to \infty} b^{p_n} = \lim_{n \to \infty} (a^{p_n} b^{p_n}) \quad \text{(7.54b)} \quad \lim_{n \to \infty} (ab)^{p_n} = (ab)^x, \\
\forall k \in \mathbb{N} \quad (a^x)^{q_k} & = \lim_{n \to \infty} (a^{p_n})^{q_k} \quad \text{(7.54c)} \quad \lim_{n \to \infty} a^{p_n q_k} = a^{x q_k}, \\
\Rightarrow \quad (a^x)^y & = \lim_{n \to \infty} (a^x)^{q_n} = \lim_{n \to \infty} a^{x q_n} \quad \text{(7.57)} \quad \Rightarrow a^{xy},
\end{align*}
\]
thereby proving (7.54).

Proceeding to (7.55c), let \( a > 1 \) and \( h > 0 \). If \((q_n)_{n \in \mathbb{N}}\) is a strictly increasing sequence in \( \mathbb{Q}^+ \) with \( \lim_{n \to \infty} q_n = h \), then \( a^h = \lim_{n \to \infty} a^{q_n} > a^{q_1} > 1 \). Thus, if \( x < y \), let \( h := y - x > 0 \) to obtain \( a^y = a^x a^h > a^x \), i.e. (7.55c). If \( 0 < a < 1 \) and \( x < y \), then \((1/a)^x < (1/a)^y\), yielding (7.55d). For (7.55a), consider \( x > 0 \) and \( 0 < a < b \). Then
\[
\frac{b}{a} > 1 \quad \Rightarrow \quad \frac{b^x}{a^x} = \left( \frac{b}{a} \right)^x > 1 \quad \Rightarrow \quad b^x > a^x,
\]
proving (7.55a). If \( x < 0 \) and \( 0 < a < b \), then \( a^x = (1/a)^{-x} > (1/b)^{-x} = b^x \), proving (7.55b).

Finally, it remains to verify (7.56). For \( x \geq 1 \), the proof for rational \( x \) still works for irrational \( x \). For \( 0 < x < 1 \), one uses the usual sequence \((q_n)_{n \in \mathbb{N}}\) in \( \mathbb{Q} \) with \( \lim_{n \to \infty} q_n = x \) and obtains (recalling \( a > 1 \))
\[
a^x = \lim_{n \to \infty} a^{q_n} \quad \text{(7.44)} \quad \leq \lim_{n \to \infty} (1 + q_n(a - 1)) = 1 + x(a - 1) < 1 + x \cdot a^{x+1},
\]
proving (7.56).

**Definition 7.71 (Exponential and Power Functions). (a)** Each function of the form
\[
f : \mathbb{R}^+ \to \mathbb{R}, \quad f(x) := x^\alpha, \quad \alpha \in \mathbb{R}, \quad (7.62)
\]
is called a **power function**. For \( \alpha > 0 \), the power function is extended to \( x = 0 \) by setting \( 0^\alpha := 0 \); for \( \alpha \in \mathbb{Z} \), it is defined on \( \mathbb{R} \setminus \{0\} \); for \( \alpha \in \mathbb{N}_0 \) even on \( \mathbb{R} \).

(b) Each function of the form
\[
f : \mathbb{R} \to \mathbb{R}^+, \quad f(x) := a^x, \quad a > 0, \quad (7.63)
is called a **(general) exponential function**. The case where \( a = e \) with \( e \) being Euler’s number from (7.49) is of particular interest and importance. Most of the time, when referring to an exponential function, one actually means \( x \mapsto e^x \). It is also common to write \( \exp(x) \) instead of \( e^x \).
Theorem 7.72. (a) Every power function as defined in Def. 7.71(a) is continuous on its respective domain. Moreover, for each $\alpha > 0$, it is strictly increasing on $[0, \infty[$; for each $\alpha < 0$, it is strictly decreasing on $]0, \infty[$.

(b) Every exponential function as defined in Def. 7.71(b) is continuous. Moreover, for each $a > 1$, it is strictly increasing; for each $0 < a < 1$, it is strictly decreasing.

Proof. (a): The monotonicity claims are provided by (7.55a) and (7.55b), respectively. For each $\alpha \in \mathbb{N}_0$, the power function is a polynomial, for each $\alpha \in \mathbb{Z}$, a rational function, i.e. continuity is provided by Ex. 7.40(b) and Ex. 7.40(c), respectively. For a general $\alpha \in \mathbb{R}$, the continuity proof on $\mathbb{R}^+$ will be postponed to Ex. 7.76(a) below, where it can be accomplished more easily. So it remains to show the continuity in $x = 0$ for $\alpha > 0$.

However, if $(x_n)_{n \in \mathbb{N}}$ is a sequence in $\mathbb{R}^+$ with $\lim_{n \to \infty} x_n = 0$ and $k \in \mathbb{N}$ with $1/k \leq \alpha$, then, at least for $n$ sufficiently large such that $x_n \leq 1$, $0 < x_n^\alpha \leq x_1^{1/k}$ by (7.55d). Then the continuity of $x \mapsto x^{1/k}$ implies $\lim_{n \to \infty} x_1^{1/k} = 0$ and the Sandwich Th. 7.16 implies $\lim_{n \to \infty} x_n^\alpha = 0$, proving continuity in $x = 0$.

(b): Everything has already been proved – continuity is provided by (7.61), monotonicity is provided by (7.55c) and (7.55d).

Remark and Definition 7.73 (Logarithm). According to Th. 7.72(b), for each $a \in \mathbb{R}^+ \setminus \{1\}$, the exponential function $f : \mathbb{R} \to \mathbb{R}^+$, $f(x) := a^x$, is continuous and strictly monotone with $f(\mathbb{R}) = \mathbb{R}^+$ (verify that the image is all of $\mathbb{R}^+$ as an exercise). Then Th. 7.60 implies the existence of a continuous and strictly monotone inverse function $f^{-1} : \mathbb{R}^+ \to \mathbb{R}$. For each $x \in \mathbb{R}^+$, we call $f^{-1}(x)$ the logarithm of $x$ to base $a$ and write $\log_a x := f^{-1}(x)$. The most important special case is where the base is Euler’s number, $a = e$. This is called the natural logarithm. Bases $a = 2$ and $a = 10$ also carry special names, binary and common logarithm, respectively. The notation is

$$\ln x := \log_e x, \quad \log_2 x := \log_2 x, \quad \log_{10} x := \log_{10} x,$$  

(7.64)

however, the notation in the literature varies – one finds log used instead of ln, lb, and lg; one also finds lg instead of lb. So you always need to verify what precisely is meant by either notation.

Corollary 7.74. For each $a \in \mathbb{R}^+ \setminus \{1\}$, the logarithm function $f : \mathbb{R}^+ \to \mathbb{R}$, $f(x) = \log_a x$ is continuous. For $a > 1$, it is strictly increasing; for $0 < a < 1$, it is strictly decreasing.

Theorem 7.75. One obtains the following logarithm rules:

\[
\begin{align*}
\forall a \in \mathbb{R}^+ \setminus \{1\} & \quad \log_a 1 = 0, & \text{(7.65a)} \\
\forall a \in \mathbb{R}^+ \setminus \{1\} & \quad \log_a a = 1, & \text{(7.65b)} \\
\forall a \in \mathbb{R}^+ \setminus \{1\} & \quad \forall x \in \mathbb{R}^+ & \quad a^{\log_a x} = x, & \text{(7.65c)} \\
\forall a \in \mathbb{R}^+ \setminus \{1\} & \quad \forall x \in \mathbb{R} & \quad \log_a a^x = x, & \text{(7.65d)}
\end{align*}
\]
Proof. All the rules are easy consequences of the logarithm being defined as the inverse function to \( f : \mathbb{R} \rightarrow \mathbb{R}^+, \ f(x) := a^x \).

(7.65a): It is \( \log_a 1 = f^{-1}(1) = 0 \), as \( f(0) = a^0 = 1 \).

(7.65b): It is \( \log_a a = f^{-1}(a) = 1 \), as \( f(1) = a^1 = a \).

(7.65c): It is \( a^{\log_a x} = f(f^{-1}(x)) = x \).

(7.65d): It is \( \log_a a^x = f^{-1}(f(x)) = x \).

(7.65e): It is \( \log_a (xy) = f^{-1}(xy) = f^{-1}(f(\log_a x + \log_a y)) = \log_a x + \log_a y \), since

\[
    f(\log_a x + \log_a y) = a^{\log_a x + \log_a y} = a^{\log_a x} a^{\log_a y} = xy. \tag{7.65c}
\]

(7.65f): It is \( \log_a (x^y) = f^{-1}(x^y) = f^{-1}(f(y \log_a x)) = y \log_a x \), since

\[
    f(y \log_a x) = a^{y \log_a x} = (a^{\log_a x})^y = x^y. \tag{7.65c}
\]

(7.65g) is just a combination of (7.65e) and (7.65f): \( \log_a (x/y) = \log_a (x^{y^{-1}}) = \log_a x - \log_a y \).

(7.65h) is just a special case of (7.65f): \( \log_a \sqrt[n]{x} = \log_a x^{1/n} = \frac{1}{n} \log_a x \).

(7.65i): One computes

\[
    (\log_b a) \log_a x = (\log_b a) \log_b a^{\log_a x} = \log_b x. \tag{7.65i}
\]

Thus, we have verified all the rules and concluded the proof.

**Example 7.76. (a) For each \( \alpha \in \mathbb{R} \), the power function

\[
    f : \mathbb{R}^+ \rightarrow \mathbb{R}, \quad f(x) := x^\alpha = e^{\alpha \ln x}, \quad (7.66)
\]

is continuous, which follows from Th. 7.41, since \( f = \exp \circ (\alpha \ln) \), \( \ln \) is continuous by Cor. 7.74, and \( \exp \) is continuous by Th. 7.72(b).
As a consequence of Th. 7.41, each of the following functions \(f_1, f_2, f_3\), where
\[f_1 : \mathbb{R} \rightarrow \mathbb{R}, \quad f_1(x) := \left(\exp(\lambda + x^2)\right)^\alpha,\]
\[f_2 : \mathbb{R} \rightarrow \mathbb{R}, \quad f_2(x) := \frac{1}{e^{\alpha x} + \lambda},\]
\[f_3 : \mathbb{R} \rightarrow \mathbb{R}, \quad f_3(x) := \frac{x^5}{(\lambda + |x|)^\alpha},\]
is continuous for each \(\alpha \in \mathbb{R}\) and each \(\lambda \in \mathbb{R}^+\).

### 7.3 Series

#### 7.3.1 Definition and Convergence

Series are a special type of sequences, namely sequences whose members arise from summing up the members of another sequence. We have, on occasion, already encountered series, for example the harmonic series \((s_n)_{n \in \mathbb{N}}\), whose members \(s_n\) were defined in (7.27). In the present section, we will study series more systematically.

**Definition 7.77.** Given a sequence \((a_n)_{n \in \mathbb{N}}\) in \(\mathbb{K}\) (or, more generally, in any set \(A\), where an addition is defined), the sequence \((s_n)_{n \in \mathbb{N}}\), where
\[
\forall \, n \in \mathbb{N} \quad s_n := \sum_{j=1}^{n} a_j,
\]
is called an (infinite) series and is denoted by
\[
\sum_{j=1}^{\infty} a_j := \sum_{j \in \mathbb{N}} a_j := (s_n)_{n \in \mathbb{N}}.
\]

The \(a_n\) are called the summands of the series, the \(s_n\) its partial sums. Moreover, each series \(\sum_{j=k}^{\infty} a_j\) with \(k \in \mathbb{N}\) is called a remainder (series) of the series \((s_n)_{n \in \mathbb{N}}\).

The example of the remainder series already shows that it is useful to allow countable index sets other than \(\mathbb{N}\). Thus, if \((a_j)_{j \in I}\), where \(I\) is a countable index set and \(\phi : \mathbb{N} \rightarrow I\) a bijective map, then define
\[
\sum_{j \in I} a_j := \sum_{j=1}^{\infty} a_{\phi(j)}
\]  
(compare the definition in (3.19c) for finite sums). Note that the definition depends on \(\phi\), which is suppressed in the notation \(\sum_{j \in I} a_j\).

For sequences in \(\mathbb{K}\), the notion of convergence is available, and, thus, it is also available for series arising from real or complex sequences (as such series are, again, sequences in \(\mathbb{K}\)).
Definition 7.78. If \((s_n)_{n \in \mathbb{N}}\) is a series with the \(s_n\) defined as in (7.67) and with summands \(a_j \in \mathbb{K}\), then the series is called convergent with limit \(s \in \mathbb{K}\) if, and only if, \(\lim_{n \to \infty} s_n = s\) in the sense of (7.1). In that case, one writes
\[
\sum_{j=1}^{\infty} a_j = s
\]
(7.70)
and calls \(s\) the sum of the series. The series is called divergent if, and only if, it is not convergent. We write \(\sum_{j=1}^{\infty} a_j = \infty\) (resp. \(\sum_{j=1}^{\infty} a_j = -\infty\)) if, and only if, \((s_n)_{n \in \mathbb{N}}\) diverges to \(\infty\) (resp. \(-\infty\)) in the sense of Def. 7.18.

Caveat 7.79. One has to use care as the symbol \(\sum_{j=1}^{\infty} a_j\) is used with two completely different meanings. If it is used according to (7.68), then it means a sequence; if it is used according to (7.70), then it means a real or complex number (or, possibly, \(\infty\) or \(-\infty\)). It should always be clear from the context, if it means a sequence or a number. For example, in the statement “the series \(\sum_{j=1}^{\infty} 2^{-j}\) is convergent”, it means a sequence; whereas in the statement “\(\sum_{j=1}^{\infty} 2^{-j} = 1\)”, it means a number.

Example 7.80. (a) For each \(q \in \mathbb{C}\) with \(|q| < 1\), \(\sum_{j=0}^{\infty} q^j\) is called a geometric series.

From (3.22b) (the reader is asked to go back and check that (3.22b) and its proof, indeed, remain valid for each \(q \in \mathbb{C}\)), we obtain the partial sums \(s_n = \sum_{j=0}^{n} q^j = \frac{1-q^{n+1}}{1-q}\). Since \(|q| < 1\), we know \(\lim_{n \to \infty} q^{n+1} = 0\) from Ex. 7.6. Thus, the series is convergent with
\[
\forall |q| < 1 \quad \sum_{j=0}^{\infty} q^j = \lim_{n \to \infty} s_n = \lim_{n \to \infty} \frac{1-q^{n+1}}{1-q} = \frac{1}{1-q}.
\]
(7.71)

(b) In Ex. 7.30, we obtained the divergence of the harmonic series:
\[
\sum_{k=1}^{\infty} \frac{1}{k} = \infty.
\]
(7.72)

Corollary 7.81. Let \(\sum_{j=1}^{\infty} a_j\) and \(\sum_{j=1}^{\infty} b_j\) be convergent series in \(\mathbb{C}\).

(a) Linearity:
\[
\forall \lambda, \mu \in \mathbb{C} \quad \sum_{j=1}^{\infty} (\lambda a_j + \mu b_j) = \lambda \sum_{j=1}^{\infty} a_j + \mu \sum_{j=1}^{\infty} b_j.
\]
(7.73)

(b) Complex Conjugation:
\[
\sum_{j=1}^{\infty} \overline{a_j} = \sum_{j=1}^{\infty} a_j.
\]
(7.74)

(c) Monotonicity:
\[
\left( \forall j \in \mathbb{N} \quad a_j, b_j \in \mathbb{R} \land a_j \leq b_j \right) \Rightarrow \sum_{j=1}^{\infty} a_j \leq \sum_{j=1}^{\infty} b_j.
\]
(7.75)
Each remainder series $\sum_{j=n+1}^{\infty} a_j$, $n \in \mathbb{N}$, converges, and, letting $S := \sum_{j=1}^{\infty} a_j$, $s_n := \sum_{j=1}^{n} a_j$, $r_n := \sum_{j=n+1}^{\infty} a_j$, one has
\[
\left( \forall n \in \mathbb{N} \ S = s_n + r_n \right), \quad \lim_{n \to \infty} a_n = \lim_{n \to \infty} r_n = 0. \tag{7.76}
\]

Proof. (a) follows from the first two identities of Th. 7.13(a), (b) is due to
\[
\sum_{j=1}^{\infty} a_j = \lim_{n \to \infty} \sum_{j=1}^{n} a_j \quad \text{Def. and Rem. 5.5(a)} = \lim_{n \to \infty} \sum_{j=1}^{n} a_j = \sum_{j=1}^{\infty} a_j,
\]
(c) follows from Th. 7.13(c), and, for (d), one computes
\[
\lim_{n \to \infty} a_n = \lim_{n \to \infty} (s_n - s_{n-1}) = S - S = 0,
\]
\[
\forall n \in \mathbb{N} \ r_n = \lim_{k \to \infty} \sum_{j=n+1}^{k} a_j = \lim_{k \to \infty} (s_k - s_n) = S - s_n,
\]
\[
\lim_{n \to \infty} r_n = \lim_{n \to \infty} (S - s_n) = S - S = 0,
\]
completing the proof. \(\blacksquare\)

### 7.3.2 Convergence Criteria

**Corollary 7.82.** Let $\sum_{j=1}^{\infty} a_j$ be series such that all $a_j \in \mathbb{R}^+_0$. If $s_n := \sum_{j=1}^{n} a_j$ are the partial sums of $\sum_{j=1}^{\infty} a_j$, then
\[
\lim_{n \to \infty} s_n = \begin{cases} 
\sup\{s_n : n \in \mathbb{N}\} & \text{if } (s_n)_{n \in \mathbb{N}} \text{ is bounded}, \\
\infty & \text{if } (s_n)_{n \in \mathbb{N}} \text{ is not bounded}. 
\end{cases} \tag{7.77}
\]

Proof. Since $(s_n)_{n \in \mathbb{N}}$ is increasing, (7.77) is a consequence of (7.21). \(\blacksquare\)

**Theorem 7.83.** Let $\sum_{j=1}^{\infty} a_j$ and $\sum_{j=1}^{\infty} b_j$ be series in $\mathbb{C}$ such that $|a_j| \leq |b_j|$ holds for each $j \geq k$ for some fixed $k \in \mathbb{N}$.

(a) If $\sum_{j=1}^{\infty} |b_j|$ is convergent, then $\sum_{j=1}^{\infty} a_j$ is convergent as well, and, moreover,
\[
\left| \sum_{j=k}^{\infty} a_j \right| \leq \sum_{j=k}^{\infty} |b_j|. \tag{7.78}
\]
(b) If $\sum_{j=1}^{\infty} a_j$ is divergent, then $\sum_{j=1}^{\infty} |b_j|$ is divergent as well.
Proof. Since (b) is merely the contraposition of (a), it suffices to prove (a). To this end, let \( s_n := \sum_{j=1}^{n} a_j \) and \( t_n := \sum_{j=1}^{n} |b_j| \) be the partial sums of \( \sum_{j=1}^{\infty} a_j \) and \( \sum_{j=1}^{\infty} |b_j| \), respectively. Since \( (t_n)_{n \in \mathbb{N}} \) converges, it must be a Cauchy sequence by Th. 7.29. Thus,

\[
\forall \epsilon \in \mathbb{R^+} \exists N \in \mathbb{N}, n > m > N \quad |t_n - t_m| = |b_{m+1}| + \cdots + |b_n| < \epsilon
\]

and the triangle inequality for finite sums implies

\[
\forall \epsilon \in \mathbb{R^+} \exists N \in \mathbb{N}, n > m > N \quad |s_n - s_m| = |a_{m+1} + \cdots + a_n| \leq |a_{m+1}| + \cdots + |a_n| \leq |b_{m+1}| + \cdots + |b_n| < \epsilon,
\]

showing \( (s_n)_{n \in \mathbb{N}} \) is a Cauchy sequence as well, i.e. convergent by Th. 7.29. Since the triangle inequality for finite sums also implies \( |\sum_{j=k}^{n} a_j| \leq \sum_{j=k}^{n} |b_j| \) for each \( n \geq k \), (7.78) is now a consequence of Th. 7.13(c). \( \blacksquare \)

Definition 7.84. A series \( \sum_{j=1}^{\infty} a_j \) in \( \mathbb{R} \) is called alternating if, and only if, its summands alternate between positive and negative signs, i.e. if \( \text{sgn}(a_{j+1}) = -\text{sgn}(a_j) \neq 0 \) for each \( j \in \mathbb{N} \).

Theorem 7.85 (Leibniz Criterion). Let \( \sum_{j=1}^{\infty} a_j \) be an alternating series. If the sequence \( (|a_n|)_{n \in \mathbb{N}} \) of absolute values is strictly decreasing and \( \lim_{n \to \infty} a_n = 0 \), then the series is convergent and

\[
\forall n \in \mathbb{N} \quad \exists r_n := \sum_{j=n+1}^{\infty} a_j = \theta_n a_{n+1}, \quad (7.79)
\]

that means the error made when approximating the limit by the partial sum \( s_n \) has the same sign as the first neglected summand \( a_{n+1} \), and its absolute value is less than \( |a_{n+1}| \).

Proof. We first consider the case where \( a_1 > 0 \), i.e. where there exists a strictly decreasing sequence of positive numbers \( (b_n)_{n \in \mathbb{N}} \) such that \( a_n = (-1)^{n+1} b_n \). As the \( b_n \) are strictly decreasing, we obtain \( b_n - b_{n+1} > 0 \) for each \( n \in \mathbb{N} \), such that the sequences \( (u_n)_{n \in \mathbb{N}} \) and \( (v_n)_{n \in \mathbb{N}} \), defined by

\[
\forall n \in \mathbb{N} \quad u_n := s_{2n} = \sum_{j=1}^{n} (b_{2j-1} - b_{2j}) = (b_1 - b_2) + (b_3 - b_4) + \cdots + (b_{2n-1} - b_{2n}),
\]

\[
\forall n \in \mathbb{N} \quad v_n := s_{2n+1} = b_1 - \sum_{j=1}^{n} (b_{2j} - b_{2j+1}) = b_1 - (b_2 - b_3) - (b_4 - b_5) - \cdots - (b_{2n} - b_{2n+1}),
\]

are strictly monotone, namely \( (u_n)_{n \in \mathbb{N}} \) strictly increasing and \( (v_n)_{n \in \mathbb{N}} \) strictly decreasing. Since, \( 0 < u_n < u_n + b_{2n+1} = v_n < b_1 \) for each \( n \in \mathbb{N} \), both sequences \( (u_n)_{n \in \mathbb{N}} \) and \( (v_n)_{n \in \mathbb{N}} \) are also bounded, and, thus, convergent by Th. 7.19, i.e. \( U := \lim_{n \to \infty} u_n \in \mathbb{R} \) and \( V := \lim_{n \to \infty} v_n \in \mathbb{R} \). Since

\[
V - U = \lim_{n \to \infty} (v_n - u_n) = \lim_{n \to \infty} (s_{2n+1} - s_{2n}) = \lim_{n \to \infty} a_{2n+1} = 0,
\]

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we obtain \( U = V \) and \( \lim_{n \to \infty} s_n = U \) and \( 0 < U < b_1 = a_1 \). In particular, there is \( \theta \in ]0, 1[ \) satisfying \( \sum_{j=1}^{\infty} a_j = \theta a_1 \).

In the case \( a_1 < 0 \), the above proof yields convergence of \( -\sum_{j=1}^{\infty} a_j = \sum_{j=1}^{\infty} (-a_j) \) with \( \sum_{j=1}^{\infty} (-a_j) = \theta (-a_1) \) for a suitable \( \theta \in ]0, 1[ \). However, this then yields, as before, \( \sum_{j=1}^{\infty} a_j = \theta a_1 \).

Applying the above result to each remainder series \( \sum_{j=n+1}^{\infty} a_j, n \in \mathbb{N} \), completes the proof of (7.79) and the theorem. \( \blacksquare \)

**Example 7.86.** (a) Each of the following alternating series clearly converges, as the Leibniz criterion of Th. 7.85 clearly applies in each case:

\[
\sum_{j=1}^{\infty} \frac{(-1)^{j+1}}{j} = 1 - \frac{1}{2} + \frac{1}{3} - \ldots, \quad (7.80a)
\]

\[
\sum_{j=1}^{\infty} \frac{(-1)^{j+1}}{2j-1} = 1 - \frac{1}{3} + \frac{1}{5} - \ldots, \quad (7.80b)
\]

\[
\sum_{j=1}^{\infty} \frac{(-1)^{j+1}}{\ln(j+1)} = \frac{1}{\ln 2} - \frac{1}{\ln 3} + \frac{1}{\ln 4} - \ldots \quad (7.80c)
\]

(b) To see that Th. 7.85 is false without its monotonicity requirement, take any divergent series with \( \sum_{j=1}^{\infty} a_j = \infty \), \( 0 < a_j \), \( \lim_{j \to \infty} a_j = 0 \) (for example the harmonic series), any convergent series with \( \sum_{j=1}^{\infty} c_j = s \in \mathbb{R}^+ \) and \( 0 < c_j \) (for example any geometric series with \( 0 < q < 1 \)), and define

\[ d_n := \begin{cases} a_{(n+1)/2} & \text{for } n \text{ odd,} \\ -c_{n/2} & \text{for } n \text{ even} \end{cases} \]

It is an exercise to show that \( \sum_{j=1}^{\infty} d_j \) is an alternating series with \( \lim_{n \to \infty} d_n = 0 \) and \( \sum_{j=1}^{\infty} d_j = \infty \).

**Definition 7.87.** The series \( \sum_{j=1}^{\infty} a_j \) in \( \mathbb{C} \) is said to be **absolutely convergent** if, and only if, \( \sum_{j=1}^{\infty} |a_j| \) is convergent.

**Corollary 7.88.** Every absolutely convergent series \( \sum_{j=1}^{\infty} a_j \) is also convergent and satisfies the triangle inequality for infinite series:

\[
\left| \sum_{j=1}^{\infty} a_j \right| \leq \sum_{j=1}^{\infty} |a_j|. \quad (7.81)
\]

**Proof.** The corollary is given by the special case \( a_j = b_j \) for each \( j \in \mathbb{N} \) of Th. 7.83(a). \( \blacksquare \)

**Theorem 7.89.** We consider the series \( \sum_{j=1}^{\infty} a_j \) in \( \mathbb{C} \).

(a) If \( \sum_{j=1}^{\infty} c_j \) is a convergent series such that \( c_j \in \mathbb{R}^+_0 \) and \( |a_j| \leq c_j \) for each \( j \in \mathbb{N} \), then \( \sum_{j=1}^{\infty} a_j \) is absolutely convergent.
(b) Root Test:

\[
\left( \exists \ 0 < q < 1 \left( \sqrt[\infty]{|a_n|} \leq q < 1 \text{ for almost all } n \in \mathbb{N} \right) \right) \\
\Rightarrow \sum_{j=1}^{\infty} a_j \text{ is absolutely convergent,} \quad (7.82a)
\]

\[
\# \left\{ n \in \mathbb{N} : \sqrt[|n|]{a_n} \geq 1 \right\} = \infty \Rightarrow \sum_{j=1}^{\infty} a_j \text{ is divergent.} \quad (7.82b)
\]

(c) Ratio Test: If all \( a_n \neq 0 \), then

\[
\left( \exists \ 0 < q < 1 \left( \left| \frac{a_{n+1}}{a_n} \right| \leq q < 1 \text{ for almost all } n \in \mathbb{N} \right) \right) \\
\Rightarrow \sum_{j=1}^{\infty} a_j \text{ is absolutely convergent,} \quad (7.83a)
\]

\[
\left| \frac{a_{n+1}}{a_n} \right| \geq 1 \text{ for almost all } n \in \mathbb{N} \Rightarrow \sum_{j=1}^{\infty} a_j \text{ is divergent.} \quad (7.83b)
\]

Proof. (a) is just another special case of Th. 7.83(a).

(b): If there is \( q \in ]0, 1[ \) and \( N \in \mathbb{N} \) such that \( \sqrt[\infty]{|a_n|} \leq q \) for each \( n > N \), i.e. \( |a_n| \leq q^n \) for each \( n > N \), then, by (7.71), \( \sum_{j=1}^{\infty} |a_j| \) is bounded by \( \frac{1}{1-q} + \sum_{j=1}^{N} |a_j| \) and, thus, convergent. If \( \sqrt[\infty]{|a_n|} \geq 1 \) for infinitely many \( n \in \mathbb{N} \), then \( |a_n| \geq 1 \) for infinitely many \( n \in \mathbb{N} \), showing that \( (a_n)_{n \in \mathbb{N}} \) does not converge to 0, proving the divergence of \( \sum_{j=1}^{\infty} a_j \).

(c): If there is \( q \in ]0, 1[ \) and \( N \in \mathbb{N} \) such that \( \left| \frac{a_{n+1}}{a_n} \right| \leq q \) for each \( n > N \), then, letting \( C := |a_{N+1}| \), an induction shows \( |a_{N+1+k}| \leq Cq^k \) for each \( k \in \mathbb{N} \), i.e., by (7.71), \( \sum_{j=1}^{\infty} |a_j| \) is bounded by \( \frac{C}{1-q} + \sum_{j=1}^{N+1} |a_j| \) and, thus, convergent. If there is \( N \in \mathbb{N} \) such that \( \left| \frac{a_{n+1}}{a_n} \right| \geq 1 \) for each \( n > N \), then \( |a_n| \geq |a_{N+1}| > 0 \) for each \( n > N \), showing \( (a_n)_{n \in \mathbb{N}} \) does not converge to 0 and proving the divergence of \( \sum_{j=1}^{\infty} a_j \).

**Caveat 7.90.** In (7.82a), it does not suffice to have \( \sqrt[\infty]{|a_n|} < 1 \) to conclude convergence, and, likewise, \( \left| \frac{a_{n+1}}{a_n} \right| < 1 \) does not suffice in (7.83a): As a counterexample, consider the harmonic series, which does not converge, but \( \sqrt[1/n]{n} < 1 \) for each \( n \geq 2 \) and \( \frac{1/(n+1)}{1/n} = \frac{n}{n+1} < 1 \) for each \( n \in \mathbb{N} \).

**Example 7.91. (a)** For each \( z \in \mathbb{C} \) with \( |z| < 1 \) and each \( p \in \mathbb{N}_0 \), the series \( \sum_{n=1}^{\infty} n^p z^n \) is absolutely convergent: We have \( \lim_{n \to \infty} \sqrt[n]{n^p} = 1 \) as a consequence of Ex. 7.65. This implies \( \lim_{n \to \infty} \sqrt[n]{|a_n|} = \lim_{n \to \infty} \sqrt[n]{n^p |z|^n} = |z| < 1 \). Thus, the root test of (7.82a) applies and proves convergence of the series.
Let \( z \in \mathbb{C} \). The series \( \sum_{n=1}^{\infty} \frac{z^n}{n^n} \) is absolutely convergent for \( |z| < e \) and divergent for \( |z| \geq e \), where \( e \) is Euler’s number from (7.49). We have, for each \( n \in \mathbb{N} \),

\[
\left| \frac{a_{n+1}}{a_n} \right| = \frac{|z|(n+1)n^n}{(n+1)n^{n+1}} = \frac{|z|n}{(1+\frac{1}{n})n} \to \frac{|z|}{e} \text{ for } n \to \infty.
\]

(7.84)

Thus, the ratio test of (7.83a) applies and proves absolute convergence of the series for \( |z| < e \). For \( |z| > e \), (7.83b) applies and proves divergence. Since, according to Ex. 7.66, \((1 + \frac{1}{n})^n < e \) for each \( n \in \mathbb{N} \), (7.83b) applies to prove divergence also for \( |z| = e \).

### 7.3.3 Absolute Convergence and Rearrangements

In general, one has to use care when dealing with infinite series, as convergence properties and even the limit in case of convergence can depend on the order of the summands (in obvious contrast to the situation of finite sums). For real series that are convergent, but not absolutely convergent, one has the striking Riemann rearrangement theorem [Phi16, Th. 7.93] that states one can choose an arbitrary number \( S \in \mathbb{R} \cup \{-\infty, \infty \} \) and reorder the summands such that the new series converges to \( S \) (actually, [Phi16, Th. 7.93] says even more, namely that one can prescribe an entire interval of cluster points for the rearranged series). However, the situation is better for absolutely convergent series.

In the present section, we will prove results that show the sum of absolutely convergent series does not depend on the order of the summands.

**Theorem 7.92.** Let \( \sum_{j=1}^{\infty} a_j \) and \( \sum_{j=1}^{\infty} b_j \) be series in \( \mathbb{C} \) such that \( (b_n)_{n \in \mathbb{N}} \) is a reordering of \( (a_n)_{n \in \mathbb{N}} \) in the sense of Def. 7.21. If \( \sum_{j=1}^{\infty} a_j \) is absolutely convergent, then so is \( \sum_{j=1}^{\infty} b_j \) and \( \sum_{j=1}^{\infty} a_j = \sum_{j=1}^{\infty} b_j \).

**Proof.** Let \( s_n := \sum_{j=1}^{n} a_j \), \( \tilde{s}_n := \sum_{j=1}^{n} |a_j| \), and \( t_n := \sum_{j=1}^{n} b_j \) denote the respective partial sums. We will show that \( \lim_{n \to \infty} (s_n - t_n) = 0 \). Given \( \epsilon > 0 \), since \( (\tilde{s}_n)_{n \in \mathbb{N}} \) is a Cauchy sequence by Th. 7.29, there exists \( N \in \mathbb{N} \), such that

\[
\forall n > m > N \quad |\tilde{s}_n - \tilde{s}_m| = |a_{m+1}| + \cdots + |a_n| < \epsilon.
\]

Since \((b_n)_{n \in \mathbb{N}}\) is a reordering of \((a_n)_{n \in \mathbb{N}}\), there exists a bijective map \( \phi : \mathbb{N} \to \mathbb{N} \) such that \( b_n = a_{\phi(n)} \) for each \( n \in \mathbb{N} \). Since \( \phi \) is bijective, there exists \( M \in \mathbb{N} \) such that \( \{1, 2, \ldots, N+1\} \subseteq \phi\{1, 2, \ldots, M\} \). Then \( n > M \) implies \( \phi(n) > N + 1 \), and

\[
\forall n > M \exists k \in \mathbb{N} \quad |s_n - t_n| \leq |a_{N+2}| + \cdots + |a_{N+k}| < \epsilon,
\]

since all \( a_j \) with \( j \leq N + 1 \) occur in both \( s_n \) and \( t_n \) and cancel in \( s_n - t_n \) (i.e. all \( a_j \) that do not cancel must have an index \( j > N + 1 \)). So we have shown that \( \lim_{n \to \infty} (s_n - t_n) = 0 \), which, in turn, implies

\[
\sum_{j=1}^{\infty} b_j = \lim_{n \to \infty} t_n = \lim_{n \to \infty} (t_n - s_n + s_n) = 0 + \sum_{j=1}^{\infty} a_j = \sum_{j=1}^{\infty} a_j.
\]
Applying this to $\tilde{s}_n := \sum_{j=1}^{n} |a_j|$ yields $\sum_{j=1}^{\infty} |b_j| = \sum_{j=1}^{\infty} |a_j|$, proving absolute convergence of $\sum_{j=1}^{\infty} b_j$.

**Theorem 7.93.** Let $I$ be an arbitrary infinite countable index set and let

$$I = \bigcup_{n \in \mathbb{N}} I_n$$

(7.85)

be a disjoint decomposition of $I$ into (empty, finite, or infinite) countable index sets $I_n$.

(a) If the series $\sum_{j \in I} a_j$ (cf. (7.69)) is absolutely convergent, then

$$\sum_{j \in I} a_j = \sum_{n=1}^{\infty} \sum_{\alpha \in I_n} a_{\alpha}. \quad (7.86)$$

(b) The following statements are equivalent:

(i) $\sum_{j \in I} a_j$ is absolutely convergent.

(ii) There exists a constant $C \in \mathbb{R}_0^+$ such that $\sum_{j \in J} |a_j| \leq C$ for each finite subset $J$ of $I$.

(iii) $\sum_{n=1}^{\infty} \sum_{\alpha \in I_n} |a_{\alpha}| < \infty$.

**Proof.** The proof needs some work; see, e.g., [Phi16, Th. 7.95].

**Example 7.94.** We apply Th. 7.93 to so-called double series, i.e. to series with index set $I := \mathbb{N} \times \mathbb{N}$. The following notation is common:

$$\sum_{m,n=1}^{\infty} a_{mn} := \sum_{(m,n) \in \mathbb{N} \times \mathbb{N}} a_{(m,n)}, \quad (7.87)$$

where one writes $a_{mn}$ (also $a_{m,n}$) instead of $a_{(m,n)}$. Recall from Th. 3.24 that $\mathbb{N} \times \mathbb{N}$ is countable. In general, the convergence properties of the double series and, if it exists, the value of the sum, will depend on the chosen bijection $\phi : \mathbb{N} \rightarrow \mathbb{N} \times \mathbb{N}$.

However, we will now assume our double series to be absolutely convergent. Then Th. 7.92 guarantees the sum does not depend on the chosen bijection and we can apply Th. 7.93. Applying Th. 7.93 to the decompositions

\begin{align}
\mathbb{N} \times \mathbb{N} &= \bigcup_{m \in \mathbb{N}} \{(m, n) : n \in \mathbb{N}\}, \quad (7.88a) \\
\mathbb{N} \times \mathbb{N} &= \bigcup_{n \in \mathbb{N}} \{(m, n) : m \in \mathbb{N}\}, \quad (7.88b) \\
\mathbb{N} \times \mathbb{N} &= \bigcup_{k \in \mathbb{N}} \{(m, n) \in \mathbb{N} \times \mathbb{N} : m + n = k\}, \quad (7.88c)
\end{align}
yields

\[
\sum_{(m,n) \in \mathbb{N} \times \mathbb{N}} a_{(m,n)} \quad \text{(7.88a)} \quad \sum_{m=1}^{\infty} \sum_{n=1}^{\infty} a_{mn} = \sum_{n=1}^{\infty} \sum_{m=1}^{\infty} a_{mn} \\
\sum_{k=2}^{\infty} \sum_{m+n=k} a_{mn} := \sum_{k=2}^{\infty} \sum_{m=1}^{k-1} a_{m,k-m}. \quad \text{(7.89)}
\]

**Theorem 7.95.** It is possible to compute the product of two absolutely convergent (real or complex) series \( \sum_{m=1}^{\infty} a_m \) and \( \sum_{m=1}^{\infty} b_m \) as a double series:

\[
\left( \sum_{m=1}^{\infty} a_m \right) \left( \sum_{m=1}^{\infty} b_m \right) = \sum_{m,n=1}^{\infty} a_m b_n = \sum_{k=2}^{\infty} \sum_{m=1}^{k-1} a_m b_{k-m} = \sum_{k=2}^{\infty} c_k,
\]

where \( c_k := \sum_{m=1}^{k-1} a_m b_{k-m} = a_1 b_{k-1} + a_2 b_{k-2} + \cdots + a_{k-1} b_1. \quad \text{(7.90)}\)

This form of computing the product is known as a Cauchy product.

**Proof.** We first show that \( \sum_{m,n=1}^{\infty} a_m b_n \) is absolutely convergent: By letting \( A := \sum_{m=1}^{\infty} |a_m| \) and \( B := \sum_{n=1}^{\infty} |b_n| \), we obtain

\[
\sum_{m=1}^{\infty} \sum_{n=1}^{\infty} |a_m b_n| = \sum_{m=1}^{\infty} (|a_m| B) = AB < \infty,
\]

i.e. \( \sum_{m,n=1}^{\infty} a_m b_n \) is absolutely convergent according to Th. 7.93(b)(iii). Now the second equality in (7.90) is just the third equality in (7.89), and the first equality in (7.90) also follows from (7.89):

\[
\sum_{m,n=1}^{\infty} a_m b_n = \sum_{m=1}^{\infty} \sum_{n=1}^{\infty} a_m b_n = \sum_{m=1}^{\infty} a_m \sum_{n=1}^{\infty} b_n = \left( \sum_{m=1}^{\infty} a_m \right) \left( \sum_{m=1}^{\infty} b_m \right),
\]

completing the proof. \( \blacksquare \)

Theorem 7.95 will be useful in Sec. 8.2 below.

### 7.3.4 \( b \)-Adic Representations of Real Numbers

We are mostly used to representing real numbers in the decimal system. For example, we write

\[
x = \frac{395}{3} = 131.\overline{5} = 1 \cdot 10^2 + 3 \cdot 10^1 + 1 \cdot 10^0 + \sum_{n=1}^{\infty} 6 \cdot 10^{-n}, \quad \text{(7.91a)}
\]

where

\[
\sum_{n=1}^{\infty} 6 \cdot 10^{-n} = 6 \cdot \left( \frac{1}{1 - \frac{1}{10}} - 1 \right) = 6 \cdot \frac{1}{9} = \frac{2}{3}.
\]
The decimal system represents real numbers as, in general, infinite series of decimal fractions. Digital computers represent numbers in the dual system, using base 2 instead of 10. For example, the number from (7.91a) has the dual representation

$$x = 10000011.10 = 2^7 + 2^1 + 2^0 + \sum_{n=0}^{\infty} 2^{-(2n+1)},$$

(7.91b)

where it is an exercise to verify

$$\sum_{n=0}^{\infty} 2^{-(2n+1)} = \frac{2}{3}.$$

Representations with base 16 (hexadecimal) and 8 (octal) are also of importance when working with digital computers. More generally, each natural number $b \geq 2$ can be used as a base.

Definition 7.96. Let $b \geq 2$ be a natural number.

(a) Given an integer $N \in \mathbb{Z}$ and a sequence $(d_N, d_{N-1}, d_{N-2}, \ldots)$ in \{0, \ldots, b-1\}, the series

$$\sum_{\nu=0}^{\infty} d_{N-\nu} b^{N-\nu}$$

(7.92)

is called a $b$-adic series. The number $b$ is called the base or the radix, and the numbers $d_\nu$ are called digits.

(b) If $x \in \mathbb{R}_0^+$ is the sum of the $b$-adic series given by (7.92), then one calls the $b$-adic series a $b$-adic representation or a $b$-adic expansion of $x$.

Theorem 7.97. Given a natural number $b \geq 2$ and a nonnegative real number $x \in \mathbb{R}_0^+$, there exists a $b$-adic series representing $x$, i.e. there is $N \in \mathbb{Z}$ and a sequence $(d_N, d_{N-1}, d_{N-2}, \ldots)$ in \{0, \ldots, b-1\} such that

$$x = \sum_{\nu=0}^{\infty} d_{N-\nu} b^{N-\nu}.$$  

(7.93)

If one introduces the additional requirement that $0 \neq d_N$, then each $x > 0$ has either a unique $b$-adic representation or precisely two $b$-adic representations. More precisely, for $0 \neq d_N$ and $x > 0$, the following statements are equivalent:

(i) The $b$-adic representation of $x$ is not unique.

(ii) There are precisely two $b$-adic representations of $x$.

(iii) There exists a $b$-adic representation of $x$ such that $d_n = 0$ for each $n \leq n_0$ for some $n_0 < N$.

(iv) There exists a $b$-adic representation of $x$ such that $d_n = b - 1$ for each $n \leq n_0$ for some $n_0 \leq N$. 

Proof. The proof is a bit lengthy and is provided in [Phi16, Sec. E.2]. ■

Example 7.98. Every natural number has precisely two decimal (i.e. 10-adic) representations. For instance,

\[ 2 = 2.\overline{0} = 1 + \sum_{n=1}^{\infty} 9 \cdot 10^{-n} = 1 + 9 \cdot \left(\frac{1}{10} - 1\right) = 1 + 9 \cdot \frac{1}{9}, \quad (7.94) \]

and analogously for all other natural numbers.

8 Convergence of \(\mathbb{K}\)-Valued Functions

8.1 Pointwise and Uniform Convergence

So far we have studied the convergence of sequences in \(\mathbb{K}\). We will now also need to study the convergence of sequences \((f_n)_{n \in \mathbb{N}}\), where each member \(f_n\) of the sequence is a function \(f_n : M \to \mathbb{K}\), \(M \subseteq \mathbb{C}\). Here, for the first time, we encounter the situation that there exist different useful notions of convergence for such sequences.

Definition 8.1. Let \((f_n)_{n \in \mathbb{N}}\) be a sequence of functions, \(f_n : M \to \mathbb{K}\), \(\emptyset \neq M \subseteq \mathbb{C}\).

(a) We say \((f_n)_{n \in \mathbb{N}}\) converges pointwise to \(f : M \to \mathbb{K}\) if, and only if, \(\lim_{n \to \infty} f_n(z) = f(z)\) for each \(z \in M\), i.e. if, and only if,

\[ \forall z \in M \quad \forall \epsilon > 0 \quad \exists N \in \mathbb{N} \quad \forall n > N \quad |f_n(z) - f(z)| < \epsilon. \quad (8.1) \]

So, in general, \(N\) in (8.1) depends on both \(z\) and \(\epsilon\).

(b) We say \((f_n)_{n \in \mathbb{N}}\) converges uniformly to \(f : M \to \mathbb{K}\) if, and only if,

\[ \forall \epsilon > 0 \quad \exists N \in \mathbb{N} \quad \forall n > N \quad \forall z \in M \quad |f_n(z) - f(z)| < \epsilon. \quad (8.2) \]

In (8.2), \(N\) is still allowed to depend on \(\epsilon\), but, in contrast to the situation of (8.1), not on \(z\) – in that sense, the convergence is uniform in \(z\).

Remark 8.2. It is immediate from Def. 8.1(a),(b) that uniform convergence implies pointwise convergence, but Ex. 8.3(b) below will show the converse is not true.

Example 8.3. (a) Let \(\emptyset \neq M \subseteq \mathbb{C}\) (for example \(M = [0,1]\) or \(M = B_1(0)\)), and \(f_n : M \to \mathbb{K}\), \(f_n(z) = 1/n\) for each \(n \in \mathbb{N}\). Then, clearly, \((f_n)_{n \in \mathbb{N}}\) converges uniformly to \(f \equiv 0\).

(b) The sequence \((f_n)_{n \in \mathbb{N}}\), where \(f_n : [0,1] \to \mathbb{R}\), \(f_n(x) := x^n\), converges pointwise, but not uniformly, to

\[ f : [0,1] \to \mathbb{R}, \quad f(x) := \begin{cases} 0 & \text{for } 0 \leq x < 1, \\ 1 & \text{for } x = 1. \end{cases} \quad (8.3) \]
For $x = 1$, $\lim_{n \to \infty} x^n = \lim_{n \to \infty} 1 = 1$, and, for $0 \leq x < 1$, $\lim_{n \to \infty} x^n = 0$ by Ex. 7.6. To see that the convergence is not uniform, consider $\epsilon := \frac{1}{2}$. Then, for every $n \in \mathbb{N}$, according to the intermediate value Th. 7.57, there exists $\xi_n \in [0, 1]$ such that $f_n(\xi_n) = \xi_n^2 = \frac{1}{2}$, i.e.

$$\forall\ n \in \mathbb{N} \ |f_n(\xi_n) - f(\xi_n)| = \xi_n^2 = \frac{1}{2} = \epsilon, \quad (8.4)$$

proving the convergence is not uniform.

**Theorem 8.4.** Let $(f_n)_{n \in \mathbb{N}}$ be a sequence of functions, $f_n : M \rightarrow \mathbb{K}$, $\emptyset \neq M \subseteq \mathbb{C}$. If $(f_n)_{n \in \mathbb{N}}$ converges uniformly to $f : M \rightarrow \mathbb{K}$ and all $f_n$ are continuous at $\zeta \in M$, then $f$ is also continuous at $\zeta$. In particular, if each $f_n$ is continuous, then so is $f$ (uniform limits of continuous functions are continuous).

**Proof.** Let $\epsilon > 0$. Due to the uniform convergence of $(f_n)_{n \in \mathbb{N}}$,

$$\exists\ m \in \mathbb{N} \ \forall\ z \in M \ |f_m(z) - f(z)| < \frac{\epsilon}{3}, \quad (8.5)$$

Due to the continuity of $f_m$ in $\zeta$,

$$\exists\ \delta > 0 \ \forall\ z \in M \cap B_{\delta}(\zeta) \ |f_m(z) - f_m(\zeta)| < \frac{\epsilon}{3}. \quad (8.6)$$

Thus,

$$\forall\ z \in M \cap B_{\delta}(\zeta) \ |f(z) - f(\zeta)| \leq |f(z) - f_m(z)| + |f_m(z) - f_m(\zeta)| + |f_m(\zeta) - f(\zeta)| < 3 \cdot \frac{\epsilon}{3} = \epsilon, \quad (8.7)$$

proving continuity of $f$ in $\zeta$. ■

### 8.2 Power Series

**Definition 8.5.** (a) In Def. 7.77, it was mentioned that series can be formed from each sequence in a set $A$, where an addition is defined. Letting $\emptyset \neq M \subseteq \mathbb{C}$, we now consider $A := \mathcal{F}(M, \mathbb{K})$, i.e. the set of functions from $M$ into $\mathbb{K}$. Then the addition on $A$ is defined according to (6.1a) and, given a sequence of functions $(f_n)_{n \in \mathbb{N}}$ in $A$, the series

$$\sum_{j=1}^{\infty} f_j := (s_n)_{n \in \mathbb{N}} \quad (8.8)$$

is defined as the sequence of partial sums $s_n := \sum_{j=1}^{n} f_j$.

(b) Given a sequence of functions $(f_n)_{n \in \mathbb{N}_0}$, where $f_n : \mathbb{K} \rightarrow \mathbb{K}$, $f_n(z) = a_n z^n$ with $a_n \in \mathbb{K}$, the series

$$\sum_{j=0}^{\infty} a_j z^j := \sum_{j=0}^{\infty} f_j \quad (8.9)$$
is called a power series and the \( a_j \) are called the coefficients of the power series. Note: The notation \( \sum_{j=0}^{\infty} a_j z^j \) introduced in (8.9) is very common, but not entirely correct, since one writes \( a_j z^j = f_j(z) \) for the summands, even though one actually means \( f_j \). Moreover, one uses the same notation if one actually does mean the series \( \sum_{j=0}^{\infty} f_j(z) \) in \( K \), so one has to see from the context if \( \sum_{j=0}^{\infty} a_j z^j \) means a series of \( K \)-valued functions or a series of numbers.

**Definition 8.6.** Consider a series of \( K \)-valued functions \( \sum_{j=1}^{\infty} f_j \) as in Def. 8.5(a), in particular, \( s_n := \sum_{j=1}^{n} f_j \) for each \( n \in \mathbb{N} \).

(a) The series converges pointwise to \( f : M \to K \) if, and only if, it (i.e. \( (s_n)_{n \in \mathbb{N}} \)) converges pointwise in the sense of Def. 8.1(a). In that case, we use the notation

\[
  f = \sum_{j=1}^{\infty} f_j. \tag{8.10}
\]

If (8.10) holds, then the series is sometimes called a series expansion of \( f \), in particular, a power series expansion if the series happens to be a power series.

Analogous to the situation of series in \( K \), the notation \( \sum_{j=1}^{\infty} f_j \) is also used with two different meanings – it can mean the sequence of partial sums as in (8.8) or, in the case of convergent series, the limit function as in (8.10) (cf. Caveat 7.79).

(b) The series converges uniformly to \( f : M \to K \) if, and only if, it converges uniformly in the sense of Def. 8.1(b).

**Corollary 8.7.** Consider a function series \( \sum_{j=1}^{\infty} f_j \) with \( f_j : M \to K \), \( \emptyset \neq M \subseteq \mathbb{C} \).

(a) The series converges uniformly to some \( f : M \to K \) if, and only if, for each \( n \in \mathbb{N} \) and each \( z \in M \), the remainder series \( \sum_{j=n+1}^{\infty} f_j(z) \) in \( K \) converges to some \( r_n(z) \in K \) such that

\[
  \forall \varepsilon \in \mathbb{R}^+ \exists N \in \mathbb{N} \forall n > N \forall z \in M \ |r_n(z)| < \varepsilon. \tag{8.11}
\]

(b) If \( \sum_{j=1}^{\infty} a_j \) is a convergent series in \( \mathbb{R}^+_0 \), then the condition

\[
  \forall z \in M \ \forall j \in \mathbb{N} \ |f_j(z)| \leq a_j \tag{8.12}
\]

implies uniform convergence of \( \sum_{j=1}^{\infty} f_j \).

(c) If each \( f_j \) is continuous in \( \zeta \in M \) and the series converges uniformly to \( f : M \to K \), then \( f \) is continuous in \( \zeta \). In particular, if each \( f_j \) is continuous, then \( f \) is continuous.

**Proof.** (a): If \( \sum_{j=1}^{\infty} f_j \) converges uniformly to \( f \), then \( f(z) = \sum_{j=1}^{\infty} f_j(z) \) holds for each \( z \in M \), \( r_n(z) = f(z) - s_n(z) \) for each \( n \in \mathbb{N} \), \( z \in M \) according to (7.76), where \( s_n(z) := \sum_{j=1}^{n} f_j(z) \). Then (8.11) is just (8.2), where the \( s_n \) now play the role of the \( f_n \) in (8.2). Conversely, if the remainder series converge for each \( z \in M \), then we can define
\( f: M \rightarrow \mathbb{K}, f(z) := f_1(z) + r_1(z) = \sum_{j=1}^{\infty} f_j(z) \). Then, once again, \( r_n(z) = f(z) - s_n(z) \) for each \( n \in \mathbb{N}, z \in M \), and (8.11) is just (8.2), yielding the uniform convergence of the series.

(b): First, (8.12) implies each remainder series \( \sum_{j=n+1}^{\infty} f_j(z) \) converges absolutely. Thus, with \( r_n(z) \) as in (a),

\[
\forall \ z \in M \ |r_n(z)| \leq \sum_{j=n+1}^{\infty} |f_j(z)| \leq \sum_{j=n+1}^{\infty} a_j \rightarrow 0 \text{ for } n \rightarrow \infty,
\]

such that (a) yields uniform convergence.

(c) is immediate from Th. 8.4. \( \square \)

**Remark 8.8.** Given a function series \( \sum_{j=1}^{\infty} f_j \) with \( f_j: M \rightarrow \mathbb{K}, \emptyset \neq M \subseteq \mathbb{C} \); for each \( z \in M \), \( \sum_{j=1}^{\infty} f_j(z) \) constitutes a series in \( \mathbb{K} \). Typically, one will only have convergence of \( \sum_{j=1}^{\infty} f_j(z) \) in \( \mathbb{K} \) on a subset \( C \subseteq M \). The series then converges pointwise in the sense of Def. 8.6(a) if all \( f_j \) are restricted to \( C \). It can be very difficult to determine if \( \sum_{j=1}^{\infty} f_j(z) \) converges or diverges for some \( z \in M \), and such investigations are often of particular interest in the context of function series. Even for power series, studying convergence can still be difficult, but the availability of the following Th. 8.9 does help to (at least partially) settle the question in many cases.

**Theorem 8.9.** For each power series \( \sum_{j=0}^{\infty} a_j z^j \), \( a_j \in \mathbb{K} \), there exists a number \( r \in [0, \infty] := \mathbb{R}^+_0 \cup \{\infty\} \), called the radius of convergence of the power series, such that

\[
\begin{align*}
\left( z \in \mathbb{K} \land |z| \leq r \right) & \Rightarrow \sum_{j=0}^{\infty} a_j z^j \text{ converges absolutely in } \mathbb{K}, \quad (8.13a) \\
\left( z \in \mathbb{K} \land |z| > r \right) & \Rightarrow \sum_{j=0}^{\infty} a_j z^j \text{ diverges in } \mathbb{K} \quad (8.13b)
\end{align*}
\]

(for \( r = \infty \), (8.13a) claims absolute convergence for each \( z \in \mathbb{K} \)). In particular, \( \sum_{j=0}^{\infty} a_j z^j \) converges pointwise in the sense of Def. 8.6(a) for each \( z \in B_r(0) \) (cf. Def. 7.7(a)). Moreover,

\[
\forall \ 0 < r_0 < r \ \left( \sum_{j=0}^{\infty} a_j z^j \text{ converges uniformly on } \overline{B}_{r_0}(0) \ (\text{cf. Ex. 7.47(a)}) \right).
\]

For the radius of convergence, one has the formula

\[
r = \frac{1}{L}, \quad \text{where } L := \limsup_{n \rightarrow \infty} \sqrt[n]{|a_n|}.
\]

In (8.15), \( \limsup \) denotes the so-called limit superior, which is defined as the largest cluster point of the sequence \( (\sqrt[n]{|a_n|})_{n \in \mathbb{N}} \) if the sequence is bounded (cf. Th. 7.27) and \( \infty \) if the sequence is unbounded. As the limit superior can be 0 or \( \infty \), we also define \( 1/0 := \infty \) and \( 1/\infty := 0 \) in (8.15).
One has the simpler formula
\[ r = \lim_{n \to \infty} \frac{|a_n|}{|a_{n+1}|}, \quad (8.16) \]
provided all \( a_n \) are nonzero and provided the limit in (8.16) either exists in \( \mathbb{R}_0^+ \) or is \( \infty \).

**Proof.** For the proof of (8.15), we apply the root test from Th. 7.89(b). Here, for the root test, we have to consider the sequence \( (\sqrt[n]{|a_n||z|^n})_{n \in \mathbb{N}} \). As a consequence of (7.11a) and Prop. 7.26, \( \limsup_{n \to \infty} (\lambda x_n) = \lambda \limsup_{n \to \infty} x_n \) for each \( \lambda > 0 \) and each sequence \( (x_n)_{n \in \mathbb{N}} \) in \( \mathbb{R} \) (with \( \lambda \cdot \infty := \infty \), this also holds if the limit superior is infinite). Thus,
\[ \limsup_{n \to \infty} \sqrt[n]{|a_n||z|^n} = |z| \limsup_{n \to \infty} \sqrt[n]{|a_n|} = |z| L. \]
If \( |z| > 1/L \), then \( |z| L > 1 \) and (7.82b) applies, i.e. (8.13b) holds for \( r = 1/L \). If \( |z| < 1/L \), then \( |z| L < 1 \), and, recalling the Bolzano-Weierstrass Th. 7.27, one sees that (7.82a) applies, i.e. (8.13a) holds for \( r = 1/L \).

Next, if \( 0 < r_0 < r \), then \( \sum_{j=0}^{\infty} |a_j r_0^j| \) converges according to (8.13a). Since, for each \( z \in B_{r_0}(0) \) and each \( j \in \mathbb{N} \), we have \( |a_j z^j| \leq |a_j r_0^j| \), (8.14) is a consequence of Cor. 8.7(b).

The validity of (8.16) follows from the ratio test of Th. 7.89(c): If all \( a_n \neq 0 \) and \( z \neq 0 \), then
\[ \lim_{n \to \infty} \left| \frac{a_{n+1} z^{n+1}}{a_n z^n} \right| = |z| \lim_{n \to \infty} \left| \frac{a_{n+1}}{a_n} \right| = \frac{|z|}{\lim_{n \to \infty} \left| \frac{a_n}{a_{n+1}} \right|}. \]
If \( |z| < l := \lim_{n \to \infty} \left| \frac{a_n}{a_{n+1}} \right| \), then \( |z|/l < 1 \), i.e. (7.83a) applies, proving (8.13a) for \( r = l \).
If \( |z| > l \), then \( |z|/l > 1 \), i.e. (7.83b) applies, proving (8.13b) for \( r = l \). \[ \square \]

**Corollary 8.10.** If \( \sum_{j=0}^{\infty} a_j z^j \), \( a_j \in \mathbb{K} \), is a power series with radius of convergence \( r \in ]0, \infty[ \), then the function
\[ f : B_r(0) \to \mathbb{K}, \quad f(z) := \sum_{j=0}^{\infty} a_j z^j, \quad (8.17) \]
is continuous. In particular, if \( r = \infty \), then \( f \) is continuous on \( \mathbb{K} \).

**Proof.** Each partial sum \( z \mapsto \sum_{j=0}^{n} a_j z^j \) is a polynomial, i.e. continuous on \( \mathbb{K} \). Moreover, if \( \zeta \in B_r(0) \), then the power series converges uniformly on \( M := B_{|\zeta|}(0) \) by (8.14), i.e. it is continuous at \( \zeta \in M \) by Th. 8.4. \[ \square \]

**Example 8.11.** (a) For each \( \alpha \in \mathbb{R} \), the radius of convergence of \( \sum_{n=1}^{\infty} n^\alpha z^n \) is \( r = 1 \), since
\[ \limsup_{n \to \infty} \sqrt[n]{|a_n|} = \lim_{n \to \infty} \sqrt[n]{n^\alpha} = 1, \quad (8.18) \]
which, for each \( \alpha \in \mathbb{Z} \), follows from (7.46) and Th. 7.13(a), and, then, for all \( \alpha \in \mathbb{R} \) from the Sandwich Th. 7.16.
Let us investigate what can happen for $|z| = r = 1$ for some cases: The series $\sum_{n=1}^{\infty} z^n$ ($\alpha = 0$) is divergent for each $z \in \mathbb{C}$ with $z = 1$ by the observation that $(z^n)_{n \in \mathbb{N}}$ does not converge to 0 for $n \to \infty$ (as $|z^n| = 1$ for each $n \in \mathbb{N}$); the series $\sum_{n=1}^{\infty} n^{-1} z^n$ ($\alpha = -1$) is the harmonic series, i.e. divergent, for $z = 1$, but convergent for $z = -1$ according to Ex. 7.86(a).

(b) The radius of convergence of both $\sum_{n=0}^{\infty} z^n$ and $\sum_{n=0}^{\infty} \frac{z^n}{n}$ is $r = \infty$ by (8.16) and (8.15), respectively, since

$$\lim_{n \to \infty} \frac{|a_n|}{|a_{n+1}|} = \lim_{n \to \infty} \frac{(n+1)!}{n!} = \lim_{n \to \infty} \frac{n+1}{n} = \infty,$$

$$\limsup_{n \to \infty} \sqrt[n]{|a_n|} = \lim_{n \to \infty} \frac{1}{n} = 0.$$  \hspace{1cm} (8.19a)

(c) The radius of convergence of $\sum_{n=0}^{\infty} n! z^n$ is $r = 0$ by (8.16), since

$$\lim_{n \to \infty} \frac{|a_n|}{|a_{n+1}|} = \lim_{n \to \infty} \frac{n!}{(n+1)!} = \lim_{n \to \infty} \frac{1}{n+1} = 0.$$  \hspace{1cm} (8.20)

Caveat 8.12. Theorem 8.9 does not claim the uniform convergence of $\sum_{j=0}^{\infty} a_j z^j$ on $B_r(0)$, which is usually not true (e.g., it is an exercise to show that $\sum_{j=0}^{\infty} z^j$ does not converge uniformly on $B_1(0)$). Theorem 8.9 also claims nothing about the convergence or divergence of $\sum_{j=0}^{\infty} a_j z^j$ for $|z| = r$, which has to be determined case by case (cf. Ex. 8.11(a)).

Definition and Remark 8.13. Given two power series $p := \sum_{j=0}^{\infty} a_j z^j$ and $q := \sum_{j=0}^{\infty} b_j z^j$ in $\mathbb{K}$, we define their Cauchy product

$$p * q := \sum_{j=0}^{\infty} c_j z^j, \quad \text{where} \quad c_j := \sum_{k=0}^{j} a_k b_{j-k} = a_0 b_j + a_1 b_{j-1} + \cdots + a_j b_0.$$  \hspace{1cm} (8.21)

Note that we have not assumed any convergence of the series so far, i.e. $p$, $q$, and $p * q$ are not $\mathbb{K}$-valued functions, but sequences of $\mathbb{K}$-valued functions according to Def. 8.5 (sequences of polynomials, actually). Sometimes one also calls the Cauchy product $p * q$ the convolution of $p$ and $q$.

Now, if we do assume $p$ and $q$ to have some nonzero radii of convergence, say $r_p, r_q \in [0, \infty]$, respectively, then, by (8.13a), both series are absolutely convergent for each $z \in B_r(0)$, where $r := \min\{r_p, r_q\}$. Thus, the functions

$$f : B_r(0) \to \mathbb{K}, \quad f(z) := \sum_{j=0}^{\infty} a_j z^j, \quad g : B_r(0) \to \mathbb{K}, \quad g(z) := \sum_{j=0}^{\infty} b_j z^j,$$

are well-defined, and (7.90) implies

$$\forall \quad z \in B_r(0) \quad f(z) g(z) = \sum_{j=0}^{\infty} c_j z^j \quad \text{with} \ c_j \text{ as in (8.21).}$$  \hspace{1cm} (8.23)
8.3 Exponential Functions

The notion of power series allows us to extend the definition of exponential functions to complex arguments:

**Definition and Remark 8.14.** We define the *exponential function*

\[
\exp : \mathbb{C} \rightarrow \mathbb{C}, \quad \exp(z) := \sum_{n=0}^{\infty} \frac{z^n}{n!} = 1 + z + \frac{z^2}{2!} + \frac{z^3}{3!} + \ldots \tag{8.24}
\]

From Ex. 8.11(b), we already know the radius of convergence of the power series in (8.24) is \(\infty\), such that the function in (8.24) is well-defined.

For the time being, we also redefine Euler’s number as \(e := \exp(1) > 1 > 0\) and, for each \(x \in \mathbb{R}^+\), \(\ln x := \log_{\exp(1)}(x)\). This, as well as calling the function of (8.24) exponential function, will be justified as soon as we will have proved

\[
\lim_{n \to \infty} \left(1 + \frac{1}{n}\right)^n = \sum_{n=0}^{\infty} \frac{1}{n!} = 1 + \frac{1}{1!} + \frac{1}{2!} + \frac{1}{3!} + \ldots \tag{8.25}
\]

and

\[
\forall x \in \mathbb{R} \quad e^x = \sum_{n=0}^{\infty} \frac{x^n}{n!} = 1 + x + \frac{x^2}{2!} + \frac{x^3}{3!} + \ldots \tag{8.26}
\]

in (8.36) of Th. 8.18 and in Th. 8.16(c) below, respectively.

**Proposition 8.15.** If a continuous function \(E : \mathbb{R} \rightarrow \mathbb{R}\) satisfies

\[
a := E(1) > 0 \quad \text{and} \quad \forall x, y \in \mathbb{R} \quad E(x + y) = E(x)E(y), \tag{8.27a, 8.27b}
\]

then \(f\) is an exponential function – more precisely

\[
\forall x \in \mathbb{R} \quad E(x) = a^x. \tag{8.28}
\]

**Proof.** First, \(a = E(1) = E(0 + 1) = E(0)E(1) = E(0)a\) and \(a > 0\) shows \(E(0) = 1\). Then, for each \(x \in \mathbb{R}, 1 = E(0) = E(x - x) = E(x)E(-x)\), i.e. \(E(-x) = (E(x))^{-1}\), showing \(E(x) \neq 0\) for each \(x \in \mathbb{R}\). Thus, \(E(1) > 0\), the continuity of \(E\), and the intermediate value Th. 7.57 imply \(E(x) > 0\) for each \(x \in \mathbb{R}\). Next, an induction shows

\[
\forall x \in \mathbb{R} \quad \forall n \in \mathbb{N} \quad E(n \cdot x) = \left(E(x)\right)^n : \tag{8.29}
\]

The base case is trivially true and the induction step is

\[
E((n + 1)x) = E(nx)E(x) \overset{\text{ind. hyp.}}{=} E(x)^n E(x) = (E(x))^{n+1}.
\]
Applying (8.29) with \( x = 1 \) shows \( E(n) = a^n \) for each \( n \in \mathbb{N} \). Applying (8.29) with \( x = 1/n, n \in \mathbb{N} \), shows \( a = E(1) = (E(1/n))^n \), i.e. \( E(1/n) = a^{1/n} \) since \( E(1/n) > 0 \). Next,

\[
\forall n,k \in \mathbb{N} \quad E(k/n) = (E(1/n))^k = (a^{1/n})^k = a^k,
\]

showing (8.28) holds for each \( x \in \mathbb{Q}^+ \). Then (8.28) also holds for each \( x \in \mathbb{R}^+ \), since, if \((q_n)_{n \in \mathbb{N}}\) is a sequence in \( \mathbb{Q}^+ \) with \( \lim_{n \to \infty} q_n = x \), then the continuity of \( E \) implies

\[
a^x = \lim_{n \to \infty} a^{q_n} = \lim_{n \to \infty} E(q_n) = E(x).
\]

Finally, if \( x \in \mathbb{R}^- \), then

\[
a^x = (a^{-x})^{-1} = (E(-x))^{-1} = E(x),
\]

completing the proof that (8.28) holds for each \( x \in \mathbb{R} \).

**Theorem 8.16.** We consider the exponential function \( \exp \) as defined in (8.24). The following holds:

(a) \( \exp \) is continuous on \( \mathbb{C} \).

(b) \( \exp(z + w) = \exp(z) \exp(w) \) is valid for all \( z, w \in \mathbb{C} \).

(c) With \( e := \exp(1) \) (cf. Def. and Rem. 8.14), it is

\[
\forall x \in \mathbb{R} \quad e^x = \exp(x) = \sum_{n=0}^{\infty} \frac{x^n}{n!}.
\]

**Proof.** (a) holds by Cor. 8.10; for (b), we compute (using (7.90)),

\[
\forall z,w \in \mathbb{C} \quad \exp(z) \exp(w) = \sum_{n=0}^{\infty} c_n,
\]

where

\[
c_n = \sum_{j=0}^{n} \frac{z^j w^{n-j}}{j! (n-j)!} = \frac{1}{n!} \sum_{j=0}^{n} \binom{n}{j} z^j w^{n-j} = \frac{(z+w)^n}{n!},
\]

and then (c) is an immediate consequence of (a), (b), and Prop. 8.15.

**Definition 8.17.** Let \( M \subseteq \mathbb{C} \). If \( \zeta \in \mathbb{C} \) is a cluster point of \( M \), then a function \( f : M \to \mathbb{K} \) is said to tend to \( \eta \in \mathbb{K} \) (or to have the limit \( \eta \in \mathbb{K} \)) for \( z \to \zeta \) (denoted by \( \lim_{z \to \zeta} f(z) = \eta \)) if, and only if, for each sequence \((z_k)_{k \in \mathbb{N}} \) in \( M \setminus \{ \zeta \} \) with \( \lim_{k \to \infty} z_k = \zeta \), the sequence \((f(z_k))_{k \in \mathbb{N}} \) converges to \( \eta \in \mathbb{K} \), i.e.

\[
\lim_{z \to \zeta} f(z) = \eta \iff \forall (z_k)_{k \in \mathbb{N}} \text{ in } M \setminus \{ \zeta \} \quad \left( \lim_{k \to \infty} z_k = \zeta \implies \lim_{k \to \infty} f(z_k) = \eta \right).
\]

(8.31)
Theorem 8.18. We consider the exponential function \( \exp \) as defined in (8.24). With \( e^z := \exp(z) \) for each \( z \in \mathbb{C} \) and \( \ln x := \log_{\exp(1)}(x) \) for each \( x \in \mathbb{R}^+ \) (cf. Th. 8.16(c) and Def. and Rem. 8.14), we have the following limits:

\[
\lim_{z \to 0} \frac{e^z - 1}{z} = 1 \quad (z \in M := \mathbb{C} \setminus \{0\}),
\]

\[
\lim_{x \to 0} \frac{\ln(1 + x)}{x} = 1 \quad (x \in M := \mathbb{R}^+) - 1, \mathbb{R} \setminus \{0\}),
\]

\[
\forall \xi \in \mathbb{R} \lim_{x \to 0} \frac{\ln(1 + \xi x)^{1/x}}{x} = \xi \quad (x \in M := \{x \in \mathbb{R} : 1 + \xi x > 0\} \setminus \{0\}),
\]

\[
\forall x \in \mathbb{R} \lim_{n \to \infty} \left(1 + \frac{x}{n}\right)^n = e^x = \sum_{n=0}^{\infty} \frac{x^n}{n!},
\]

Proof. (8.32): From (8.24) and \( e^z = \exp(z) \), we obtain

\[
\forall z \neq 0 \quad e^z - 1 = \sum_{n=0}^{\infty} \frac{z^n}{(n+1)!} = 1 + \frac{z}{2!} + \frac{z^2}{3!} + \ldots,
\]

which, since \( z \mapsto \sum_{n=0}^{\infty} \frac{z^n}{(n+1)!} \) is continuous on \( \mathbb{C} \) by Cor. 8.10, implies (8.32).

(8.33): Consider the auxiliary function \( f :] - 1, \infty[ \to \mathbb{R}, f(x) := \ln(x + 1), \) with \( f^{-1}(x) = e^x - 1 \). Now, given a sequence \((x_k)_{k \in \mathbb{N}}\) in \( ] - 1, \infty[ \setminus \{0\} \) with \( \lim_{k \to \infty} x_k = 0 \), one obtains

\[
\lim_{k \to \infty} \ln(1 + x_k) = \lim_{k \to \infty} \ln \left(1 + f^{-1}(f(x_k))\right) = \lim_{k \to \infty} \ln \left(1 + e^{f(x_k)} - 1\right)
\]

\[
= \lim_{k \to \infty} \frac{f(x_k)}{e^{f(x_k)} - 1} \quad (8.32)
\]

where, in the last step, it was used that \( \lim_{k \to \infty} x_k = 0 \) and the continuity of \( f \) implies \( \lim_{k \to \infty} f(x_k) = \ln 1 = 0 \).

Similarly, but simpler, one obtains (8.34) and (8.35) (exercise). Finally, for the sequence \((x_n)_{n \in \mathbb{N}}\) with \( x_n := 1/n \), (8.35) implies (8.36).

Definition 8.19 (Exponentiation with Complex Exponents). For each \((a, z) \in \mathbb{R}^+ \times \mathbb{C}\), we define

\[
a^z := \exp(z \ln a),
\]

where \( \exp \) is the function defined in (8.24). For \( a = e \), (8.37) yields \( e^z = \exp(z) \), i.e. (8.37) is consistent with (8.26).

Theorem 8.20. (a) The first two exponentiation rules of (7.54) still hold for each \( a, b > 0 \) and each \( z, w \in \mathbb{C} \):

\[
a^{z+w} = a^z a^w, \quad (8.38a)
\]

\[
a^z b^w = (ab)^z. \quad (8.38b)
\]
(b) For each \( a \in \mathbb{R}^+ \), the exponential function
\[
f : \mathbb{C} \rightarrow \mathbb{C}, \quad f(z) := a^z,
\]
is continuous, and, for each \( \zeta \in \mathbb{C} \), the power function
\[
g : \mathbb{R}^+ \rightarrow \mathbb{C}, \quad g(x) := x^\zeta,
\]
is continuous.

(c) The limit in (8.36) extends to complex numbers:
\[
\forall \ z \in \mathbb{C} \lim_{n \to \infty} \left(1 + \frac{z}{n}\right)^n = e^z = \sum_{n=0}^{\infty} \frac{z^n}{n!}.
\]

Proof. (a): We compute
\[
a^z a^w \overset{(8.37)}{=} \exp((z + w) \ln a) = \exp(z \ln a + w \ln a) \overset{\text{Th. 8.16(b)}}{=} \exp(z \ln a) \exp(w \ln a) \overset{(8.37)}{=} a^z a^w,
\]
proving (8.38a), and
\[
a^z b^z \overset{(8.37)}{=} \exp(z \ln a) \exp(z \ln b) \overset{\text{Th. 8.16(b)}}{=} \exp(z \ln a + z \ln b) \overset{(7.65e)}{=} \exp(z \ln(ab)) \overset{(8.37)}{=} (ab)^z,
\]
proving (8.38b).

(b): The continuity of both functions follows from the continuity of \( \exp \) (according to Th. 8.16(a)) and from the fact that continuity is preserved by compositions (according to Th. 7.41): The exponential function \( f \), given by \( f(z) = e^{z \ln a} \), is the composition of the continuous functions \( z \mapsto z \ln a \) and \( w \mapsto e^w \), whereas (analogous to Ex. 7.76(a)), the power function \( g \), given by \( g(x) = e^{\zeta \ln x} \), is the composition \( g = \exp \circ (\zeta \ln) \), where \( \ln \) is continuous by Cor. 7.74.

(c): We have to show that
\[
\lim_{n \to \infty} \left|\left(1 + \frac{z}{n}\right)^n - \sum_{k=0}^{\infty} \frac{z^k}{k!}\right| = 0.
\]
Given \( \epsilon > 0 \), choose \( K \in \mathbb{N} \) such that
\[
\forall \ z \in \mathbb{C} \quad \sum_{k=n}^{\infty} \frac{(|z| + 1)^k}{k!} < \epsilon.
\]
We continue by using (5.23) to estimate
\[
\forall \ n \in \mathbb{N} \quad A_n := \left|\left(1 + \frac{z}{n}\right)^n - \sum_{k=0}^{\infty} \frac{z^k}{k!}\right| \leq R_n + S_n + T,
\]
where
\[
\forall n \in \mathbb{N} \quad R_n := \sum_{k=0}^{K-1} \left| \binom{n}{k} \frac{z^k}{n^k} - \frac{z^k}{k!} \right|, \quad S_n := \sum_{k=K}^{n} \binom{n}{k} \frac{|z|^k}{n^k}, \quad T := \sum_{k=K}^{\infty} \frac{|z|^k}{k!}.
\]

We proceed to estimate each of the three terms \( R_n, S_n, \) and \( T \), starting with the last:
\[
T = \sum_{k=K}^{\infty} \frac{|z|^k}{k!} < \sum_{k=K}^{\infty} \frac{(|z| + 1)^k}{k!} < \frac{\epsilon}{3}.
\]

To estimate \( S_n \), we first estimate
\[
\forall n \in \mathbb{N} \quad \forall 1 \leq k \leq n \quad \binom{n}{k} \frac{1}{n^k} = \frac{n!}{k!(n-k)!n^k} = \frac{1}{k!} \prod_{j=1}^{k} \frac{n-k+j}{n} \leq \frac{1}{k!},
\]

We then obtain
\[
\forall n \geq K \quad S_n = \sum_{k=K}^{n} \binom{n}{k} \frac{|z|^k}{n^k} \leq \sum_{k=K}^{\infty} \frac{|z|^k}{k!} = T < \frac{\epsilon}{3}.
\]

To estimate \( R_n \), we first compute the limit
\[
\lim_{n \to \infty} \binom{n}{k} \frac{1}{n^k} = \frac{1}{k!} \lim_{n \to \infty} \prod_{j=1}^{k} \frac{n-k+j}{n} = \frac{1}{k!} \prod_{j=1}^{k} 1 = \frac{1}{k!},
\]
implies \( \lim_{n \to \infty} R_n = 0 \) and
\[
\exists N \geq K \quad \forall n > N \quad R_n < \frac{\epsilon}{3}.
\]

Combining the three estimates shows
\[
\forall n > N \geq K \quad A_n \leq R_n + S_n + T < \frac{\epsilon}{3} + \frac{\epsilon}{3} + \frac{\epsilon}{3} = \epsilon,
\]
completing the proof. ■

### 8.4 Trigonometric Functions

The first “definition” of the trigonometric functions sine and cosine is the one based on geometric visualization usually given in high school: \( \cos x \) and \( \sin x \) are the coordinates of the point \( p = (p_1, p_2) \in \mathbb{R}^2 \) on the unit circle, such that \( x \) is the angle measured in radian between the line segment between \((0, 0)\) and \((1, 0)\) and the line segment between \((0, 0)\) and \(p\).

While this “definition” allows to obtain many important properties of sine and cosine using geometric arguments, it is not mathematically rigorous, and, for example, provides no clue how to compute values like \( \sin 1 \). The problem is related to the fact that the angle measured in radian between the line segment between \((0, 0)\) and \((1, 0)\) and the line
segment between \((0,0)\) and \(p\) is supposed to be the length of the segment of the unit
circle between \((1,0)\) and \(p\) (taken in the counter-clockwise direction).

In the following Def. and Rem. 8.21, we will provide a mathematically rigorous definition
of sine and cosine using power series, and we will then verify that the functions have the
familiar properties one learns in high school. However, as the computation of lengths
of curved paths is actually beyond the scope of this lecture, we will not be able to see
that our sine and cosine functions are precisely the same we visualized in high school
(the interested reader is referred to Ex. 1 in Sec. 5.14 of [Wal02]).

**Definition and Remark 8.21.** We define the sine function, denoted \(\sin\), and the
cosine function, denoted \(\cos\) by

\[
\sin : \mathbb{C} \rightarrow \mathbb{C}, \quad \sin z := \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n+1}}{(2n+1)!} = z - \frac{z^3}{3!} + \frac{z^5}{5!} - + \ldots, \tag{8.41a}
\]

\[
\cos : \mathbb{C} \rightarrow \mathbb{C}, \quad \cos z := \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n}}{(2n)!} = 1 - \frac{z^2}{2!} + \frac{z^4}{4!} - + \ldots. \tag{8.41b}
\]

(a) \(\sin\) and \(\cos\) are well-defined and continuous: For both series and each \(z \in \mathbb{C}\), we can
estimate the absolute value of the \(n\)th summand by the \(n\)th summand of the series
for the exponential function \(e^{|z|}\) (cf. (8.36)), which we know to be convergent from
Ex. 8.11(b). Thus, by Th. 8.9, both series in (8.41) have radius of convergence \(\infty\)
and are continuous by Cor. 8.10.

(b) \(\cos : \mathbb{R} \rightarrow \mathbb{R}\) (i.e. \(\cos |_{\mathbb{R}}\)) has a smallest positive zero \(\alpha \in \mathbb{R}^+\). We define \(\pi := 2\alpha\).
One can show \(\pi\) is an irrational number (see [Phi16, Sec. H.2]) and its first digits
are \(\pi = 3.14159\ldots\)

To see \(\cos\) has a smallest positive zero and to obtain a first (very coarse) estimate, note

\[
\forall x \in \mathbb{R}^+ \quad \forall k \in \mathbb{N} \quad \left( \frac{x^k}{k!} > \frac{x^{k+1}}{(k+1)!} \right) \iff 1 > \frac{x}{k+1} \iff k + 1 > x,
\]

showing \(\frac{x^k}{k!} > \frac{x^{k+1}}{(k+1)!}\) holds for each \(k \geq 2\) and each \(x \in ]0,3[\). In particular, the
summands of the series in (8.41) converge monotonically to 0 (for \(k \geq 2\)) and, since
the series are alternating for \(x \neq 0\), Th. 7.85 applies and (7.79) yields

\[
\forall 0 < x < 3 \quad \begin{cases} f(x) := 1 - \frac{x^2}{2} < \cos x < 1 - \frac{x^2}{2} + \frac{x^4}{24} =: g(x), \\
x - \frac{x^3}{6} < \sin x < x - \frac{x^3}{6} + \frac{x^5}{120}. \end{cases} \tag{8.42}
\]

The zeros of \(x \mapsto f(x)\) are \(-\sqrt{2}, \sqrt{2}\), i.e. \(\sqrt{2}\) is its smallest positive zero; the zeros
of \(x \mapsto g(x)\) are \(-\sqrt{6 - 2\sqrt{3}}, -\sqrt{6 + 2\sqrt{3}}, \sqrt{6 - 2\sqrt{3}}, \sqrt{6 + 2\sqrt{3}}\), i.e. \(\sqrt{6 - 2\sqrt{3}}\)
is its smallest positive zero. Thus, as \(f(0) = g(0) = 1\), the intermediate value Th.
7.57 implies \(\cos\) has a smallest positive zero \(\alpha\) and

\[
1.4 < \frac{\sqrt{2}}{2} < \frac{\pi}{2} := \alpha < \sqrt{6 - 2\sqrt{3}} < 1.6 \tag{8.43}
\]
Theorem 8.22. We have the following identities:

\begin{align}
\forall z \in \mathbb{C} & \quad \sin z = -\sin(-z), \quad \cos z = \cos(-z), \\
\forall z \in \mathbb{C} & \quad \sin(z + w) = \sin z \cos w + \cos z \sin w, \\
\forall z \in \mathbb{C} & \quad \cos(z + w) = \cos z \cos w - \sin z \sin w, \\
\forall z \in \mathbb{C} & \quad (\sin z)^2 + (\cos z)^2 = 1, \\
\forall z \in \mathbb{C} & \quad \sin z = -\sin(\pi/2 + z), \\
\forall z \in \mathbb{C} & \quad \cos z = -\cos(\pi/2 + z), \\
\forall z \in \mathbb{C} & \quad \sin(2\pi + z) = \sin z, \\
\forall z \in \mathbb{C} & \quad \cos(2\pi + z) = \cos z, \\
\lim_{z \to 0} \frac{\sin z}{z} = 1, \quad \lim_{z \to 0} \frac{\cos z - 1}{z^2} = -\frac{1}{2}.
\end{align}

Identities (8.44i) can be restated as sine and cosine being periodic functions with period $2\pi$.

Proof. (8.44a) is immediate from (8.41) since, for $z = 0$, all summands of the sine series are 0 and all summands of the cosine series are 0, except the first one, which is $\frac{(-1)^0 0!}{0!} = 1$.

(8.44b) is also immediate from (8.41), since $(-z)^{2n+1} = (-1)^{2n+1}z^{2n+1} = -z^{2n+1}$ and $(-z)^{2n} = (-1)^{2n}z^{2n} = z^{2n}$.

(8.44c) and (8.44d) can be verified using the Cauchy product: According to (7.90),

\[
\forall z, w \in \mathbb{C} \quad \begin{pmatrix}
\sin z \cos w = \sum_{n=0}^{\infty} c_n, \\
\cos z \sin w = \sum_{n=0}^{\infty} d_n,
\end{pmatrix}
\]

where

\[
c_n = \sum_{j=0}^{n} \frac{(-1)^j z^{2j+1} (-1)^{n-j} w^{2(n-j)}}{(2j+1)! (2(n-j))!},
\]

\[
d_n = \sum_{j=0}^{n} \frac{(-1)^j z^{2j} (-1)^{n-j} w^{2(n-j)+1}}{(2j)! (2(n-j)+1)!},
\]

and

\[
\forall z \in \mathbb{C} \quad \sin z = -\sin(\pi/2 + z), \\
\forall z \in \mathbb{C} \quad \cos z = -\cos(\pi/2 + z), \\
\forall z \in \mathbb{C} \quad \sin(2\pi + z) = \sin z, \\
\forall z \in \mathbb{C} \quad \cos(2\pi + z) = \cos z, \\
\lim_{z \to 0} \frac{\sin z}{z} = 1, \quad \lim_{z \to 0} \frac{\cos z - 1}{z^2} = -\frac{1}{2}.
\]
that means, for each \( z, w \in \mathbb{C} \),

\[
c_n + d_n = \sum_{j=0}^{n} \frac{(-1)^n z^{2j+1}}{(2j+1)!} \cdot \frac{w^{2(n-j)}}{(2(n-j))!} + \sum_{j=0}^{n} \frac{(-1)^n z^{2j}}{(2j)!} \cdot \frac{w^{2(n-j)+1}}{2(n-j) + 1} \tag{8.44c}
\]

\[
= \sum_{j=0}^{n} \frac{(-1)^n z^{2j+1}}{(2j+1)!} \cdot \frac{w^{2n+1-(2j+1)}}{(2n + 1 - (2j + 1))!} + \sum_{j=0}^{n} \frac{(-1)^n z^{2j}}{(2j)!} \cdot \frac{w^{2n+1-2j}}{(2n + 1 - 2j)!} \tag{8.44d}
\]

\[
= (-1)^n \sum_{j=0}^{2n+1} \frac{z^{j}}{j!} \cdot \frac{w^{2n+1-j}}{(2n+1)!} = \left(\frac{-1)^n}{(2n+1)!}\right) \sum_{j=0}^{2n+1} \binom{2n+1}{j} z^{j} w^{2n+1-j}
\]

proving (8.44c). Similarly, according to (7.90),

\[
\forall \ z, w \in \mathbb{C} \quad \begin{cases} 
\cos z \cos w = \sum_{n=0}^{\infty} c_n, \\
\sin z \sin w = \sum_{n=0}^{\infty} d_n,
\end{cases}
\]

where

\[
c_n = \sum_{j=0}^{n} \frac{(-1)^j z^{2j} (-1)^{n-j} w^{2(n-j)}}{(2j)!} \cdot \frac{w^{2(n-j)+1}}{(2n-j)!},
\]

\[
d_n = \sum_{j=0}^{n} \frac{(-1)^j z^{2j+1} (-1)^{n-j} w^{2(n-j)+1}}{(2j+1)!} \cdot \frac{w^{2n-(2j+1)}}{(2n-(2j+1))!}
\]

that means, for each \( z, w \in \mathbb{C} \),

\[
\forall \ n \in \mathbb{N} \quad c_n - d_{n-1} = \sum_{j=0}^{n} \frac{(-1)^n z^{2j}}{(2j)!} \cdot \frac{w^{2(n-j)}}{(2(n-j))!} - \sum_{j=0}^{n-1} \frac{(-1)^n z^{2j+1}}{(2j+1)!} \cdot \frac{w^{2(n-j)+1}}{(2n-1-j) + 1} \tag{8.44d}
\]

\[
= \sum_{j=0}^{n} \frac{(-1)^n z^{2j}}{(2j)!} \cdot \frac{w^{2n-2j}}{(2n-2j)!} + \sum_{j=0}^{n-1} \frac{(-1)^n z^{2j+1}}{(2j+1)!} \cdot \frac{w^{2n-(2j+1)}}{(2n-(2j+1))!}
\]

\[
= (-1)^n \sum_{j=0}^{2n} \frac{z^{j}}{j!} \cdot \frac{w^{2n-j}}{(2n-j)!} = \left(\frac{-1)^n}{(2n)!}\right) \sum_{j=0}^{2n} \binom{2n}{j} z^{j} w^{2n-j}
\]

proving (8.44d).

(8.44e): One computes for each \( z \in \mathbb{C} \):

\[
(sin z)^2 + (cos z)^2 = cos z \cos(-z) - sin z \sin(-z) \quad (8.44d)
\]

\[
= \cos(z-z) = \cos 0 = 1.
\]

(8.44f): \(\cos \frac{\pi}{2} = 0\) and \(\cos x > 0\) for \(0 \leq x < \frac{\pi}{2}\) hold according to the definition of \(\pi\) in Def. and Rem. 8.21(b). Then

\[
\left(\sin \frac{\pi}{2}\right)^2 - \left(\cos \frac{\pi}{2}\right)^2 = 1 \quad \text{and} \quad \sin \frac{\pi}{2} > \frac{\pi}{2} - \left(\frac{\pi/2}{6}\right)^3 \quad (8.43) \quad > 1.4 - \left(\frac{1.6}{6}\right)^3 \quad > 0.7 > 0.
\]
(8.44g) is immediate from (8.44c), (8.44d), and (8.44f).

(8.44h): One obtains

\[
\sin \pi = \sin \left( \frac{\pi}{2} + \frac{\pi}{2} \right) = 1 \cdot 0 + 0 \cdot 1 = 0,
\]
\[
\cos \pi = \cos \left( \frac{\pi}{2} + \frac{\pi}{2} \right) = 0 \cdot 0 - 1 \cdot 1 = -1,
\]

\[
\forall z \in \mathbb{C} \quad \sin(z + \pi) = -\sin z + 0 = -\sin z,
\]
\[
\forall z \in \mathbb{C} \quad \cos(z + \pi) = -\cos z + 0 = -\cos z.
\]

(8.44i): One obtains

\[
\sin(2\pi) = \sin(\pi + \pi) = 0 + 0 = 0,
\]
\[
\cos(2\pi) = \cos(\pi + \pi) = (-1)(-1) - 0 = 1,
\]

\[
\forall z \in \mathbb{C} \quad \sin(z + 2\pi) = \sin z + 0 = \sin z,
\]
\[
\forall z \in \mathbb{C} \quad \cos(z + 2\pi) = \cos z - 0 = \cos z.
\]

(8.44j): One obtains

\[
\forall z \in \mathbb{C} \setminus \{0\} \quad \begin{pmatrix}
\frac{\sin z}{z} = \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n}}{(2n+1)!} = 1 - \frac{z^2}{3!} + \frac{z^4}{5!} - + \ldots,
\frac{\cos z - 1}{z^2} = \sum_{n=0}^{\infty} \frac{(-1)^{n+1} z^{2n}}{(2(n+1))!} = -\frac{1}{2!} + \frac{z^2}{4!} - \frac{z^4}{6!} + + \ldots
\end{pmatrix}.
\]

For both series on the right-hand side and each \( z \in \mathbb{C} \), we can estimate the absolute value of each summand by the corresponding summand of the exponential series for \( e^{|z|} \) (cf. (8.36)), showing they have radius of convergence \( \infty \) and are continuous by Cor. 8.10. In particular, their continuity in \( z = 0 \) proves (8.44j).

**Theorem 8.23.** One has \( \sin(\mathbb{R}) = \cos(\mathbb{R}) = [-1, 1] \), i.e. the range of both sine and cosine is \([-1, 1]\). Moreover, for each \( k \in \mathbb{Z} \):

\[
\sin \text{ is strictly increasing on } \left[ -\frac{\pi}{2} + 2k\pi, \frac{\pi}{2} + 2k\pi \right], \quad (8.45a)
\]
\[
\sin \text{ is strictly decreasing on } \left[ \frac{\pi}{2} + 2k\pi, \frac{3\pi}{2} + 2k\pi \right], \quad (8.45b)
\]
\[
\cos \text{ is strictly increasing on } \left[ (2k-1)\pi, 2k\pi \right], \quad (8.45c)
\]
\[
\cos \text{ is strictly decreasing on } \left[ 2k\pi, (2k+1)\pi \right], \quad (8.45d)
\]

which, due to (8.44e), can be summarized (and visualized) by saying that, if \( x \) runs from \( 2k\pi \) to \( 2(k+1)\pi \), then \( (\cos x, \sin x) \) runs once counterclockwise through the unit circle, starting at \((1,0)\).
Proof. From (8.44e), we know \( \sin(\mathbb{R}) \subseteq [-1,1] \) and \( \cos(\mathbb{R}) \subseteq [-1,1] \). As 

\[
\sin \left( \frac{\pi}{2} \right) = 1, \quad \sin \left( -\frac{\pi}{2} \right) = -1, \quad \cos 0 = 1, \quad \cos \pi = -1, \quad \cos 0 = -1,
\]

the continuity of sine and cosine together with the intermediate value Th. 7.57 implies \( \sin(\mathbb{R}) = \cos(\mathbb{R}) = [-1,1] \).

From (8.42), we know \( 0 < x - \frac{x^3}{6} < \sin x < 1 - \frac{x^2}{2} + \frac{x^4}{24} < 1 \) for each \( x \in ]0, \frac{\pi}{2}] \), implying

\[
\forall 0 \leq x < x + y \leq \frac{\pi}{2} \quad \cos(x + y) = \cos x \cos y - \sin x \sin y \leq \cos x \cos y < \cos x,
\]

showing \( \cos \) is strictly decreasing on \( [0, \frac{\pi}{2}] \). Then \( \cos \) is strictly increasing on \( [\frac{\pi}{2}, 0] \) by (8.44b), sin is strictly increasing on \( [0, \frac{\pi}{2}] \) and strictly decreasing on \( [\frac{\pi}{2}, \pi] \) by (8.44g), implying sin is strictly increasing on \( [-\frac{\pi}{2}, 0] \) and strictly decreasing on \( [-\pi, -\frac{\pi}{2}] \) by (8.44b), i.e. sin is strictly increasing on \( [\frac{3\pi}{2}, 2\pi] \) and strictly decreasing on \( [\pi, \frac{3\pi}{2}] \) by (8.44i), implying \( \cos \) is strictly decreasing on \( [\frac{\pi}{2}, \pi] \) and strictly increasing on \( [-\pi, -\frac{\pi}{2}] \) by (8.44g). Since this fixes the monotonicity properties of both sine and cosine over more than one period, the general statements in (8.45) are provided by (8.44i).

We now come to important complex number relations between sine, cosine, and the exponential function.

**Theorem 8.24.** One has the following formulas, relating the (complex) sine, cosine, and exponential function:

\[
\forall z \in \mathbb{C} \quad e^{iz} = \cos z + i \sin z \quad \text{(Euler formula),} \quad (8.46a)
\]

\[
\forall z \in \mathbb{C} \quad \cos z = \frac{e^{iz} + e^{-iz}}{2}, \quad (8.46b)
\]

\[
\forall z \in \mathbb{C} \quad \sin z = \frac{e^{iz} - e^{-iz}}{2i}. \quad (8.46c)
\]

**Proof.** Let \( z \in \mathbb{C} \). For (8.46a), one computes

\[
e^{iz} \overset{(8.37),(8.24)}{=} \sum_{n=0}^{\infty} \frac{(iz)^n}{n!} = \sum_{n=0}^{\infty} \left( \frac{(-1)^n z^{2n}}{(2n)!} + i \frac{(-1)^n z^{2n+1}}{(2n+1)!} \right) \overset{\text{Th. 7.93(a)}}{=} \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n}}{(2n)!} + i \sum_{n=0}^{\infty} \frac{(-1)^n z^{2n+1}}{(2n+1)!} \overset{(8.41)}{=} \cos z + i \sin z.
\]

Then

\[
e^{iz} + e^{-iz} \overset{(8.46a)}{=} \cos z + i \sin z + \cos(-z) + i \sin(-z) = 2 \cos z
\]

proves (8.46b), and

\[
e^{iz} - e^{-iz} \overset{(8.46a)}{=} \cos z + i \sin z - \cos(-z) - i \sin(-z) = 2i \sin z
\]

proves (8.46c).
As a first application of (8.46), we can now determine all solutions to the equation $e^z = 1$ and all zeros (if any) of exp, sin, and cos:

**Theorem 8.25.** The set of (complex) solutions to the equation $e^z = 1$ consists precisely of all integer multiples of $2\pi i$, the exponential function has no zeros (neither in $\mathbb{R}$ nor in $\mathbb{C}$), and the set of all (real or complex) zeros of sine and cosine consists of a discrete set of real numbers. More precisely:

\[
\begin{align*}
\exp^{-1}\{1\} & = \{2k\pi i : k \in \mathbb{Z}\}, \\
\exp^{-1}\{0\} & = \emptyset, \\
\sin^{-1}\{0\} & = \{k\pi : k \in \mathbb{Z}\}, \\
\cos^{-1}\{0\} & = \{ (2k + 1) \frac{n}{2} : k \in \mathbb{Z} \}.
\end{align*}
\]

**Proof.** We start by considering the zeros of the functions $\cos, \sin : \mathbb{R} \rightarrow \mathbb{R}$: Due to (8.44f), $\cos x > 0$ for each $x \in \left[0, \frac{\pi}{2}\right]$ such that $\cos(-x) = \cos x$ (by (8.44b)) implies $\frac{\pi}{2}$ to be the only zero of $\cos$ in the interval $\left[-\frac{\pi}{2}, \frac{\pi}{2}\right]$. Then, since $\cos(x + \pi) = -\cos x$ for each $x \in \mathbb{R}$ by (8.44g), $\frac{\pi}{2}$ and $\frac{\pi}{2} + \pi$ are the only zeros of $\cos$ in the interval $\left[-\frac{\pi}{2}, \frac{\pi}{2} + \pi\right]$, and, thus, using that $\cos$ has period $2\pi$ according to (8.44i), adding integer multiples of $2\pi$ to $\frac{\pi}{2}$ and $\frac{\pi}{2} + \pi$ must generate precisely all zeros of $\cos : \mathbb{R} \rightarrow \mathbb{R}$, i.e.

$$\mathbb{R} \cap \cos^{-1}\{0\} = \left\{ \frac{\pi}{2} + k\pi : k \in \mathbb{Z} \right\} = \left\{ (2k + 1) \frac{\pi}{2} : k \in \mathbb{Z} \right\}.$$  

Since, by (8.44g), $\sin x = -\cos(x + \frac{\pi}{2})$ for each $x \in \mathbb{R}$, we also obtain

$$\mathbb{R} \cap \sin^{-1}\{0\} = \left\{ -\frac{\pi}{2} + x : x \in \mathbb{R} \cap \cos^{-1}\{0\} \right\} = \{k\pi : k \in \mathbb{Z}\}.$$  

We consider (8.47a) next. If $k \in \mathbb{Z}$, then

$$e^{2k\pi i} = \left(e^{2\pi i}\right)^k = \cos(2\pi) + i\sin(2\pi) = 1 = 1^k = 1,$$

proving “$\supseteq$”. For the remaining inclusion, assume $z \in \exp^{-1}\{1\}$, i.e. $e^z = 1$, and write $z = x + iy$ with $x, y \in \mathbb{R}$. Then

\[
1 = |e^z| = e^x |e^{iy}| = e^x \sqrt{\cos^2 y + \sin^2 y} = e^x,
\]

first implying $x = 0$ and, then, using (8.46a) once again, $1 = e^z = e^{iy} = \cos y + i\sin y$ implies $\cos y = 1$ and $\sin y = 0$, i.e. $y \in \{2k\pi : k \in \mathbb{Z}\}$, proving “$\subseteq$”.

To finish the proof of (8.47c), assume $\sin z = 0$. Then $e^{iz} = \cos z = \cos(-z) = e^{-iz}$, implying $e^{2iz} = 1$ and, by (8.47a), there is $k \in \mathbb{Z}$ such that $2iz = 2k\pi i$, i.e. $z = k\pi$, proving (8.47c). Since, by (8.44g), $\cos z = \sin(z + \frac{\pi}{2})$ for each $z \in \mathbb{C}$, we also obtain (8.47d):

$$\cos^{-1}\{0\} = \left\{ -\frac{\pi}{2} + z : z \in \sin^{-1}\{0\} \right\} = \left\{ (2k + 1) \frac{\pi}{2} : k \in \mathbb{Z} \right\}.$$  

Finally, if $z = x + iy$ with $x, y \in \mathbb{R}$, then $|e^z| = e^x |e^{iy}| = e^x \neq 0$ proves (8.47b).  

8 CONVERGENCE OF $\mathbb{K}$-VALUED FUNCTIONS

**Definition and Remark 8.26.** We define *tangent* and *cotangent* by

\[
\tan : \mathbb{C} \setminus \cos^{-1}\{0\} \to \mathbb{C}, \quad \tan z := \frac{\sin z}{\cos z},
\]

\[
(8.48a)
\]

\[
\cot : \mathbb{C} \setminus \sin^{-1}\{0\} \to \mathbb{C}, \quad \cot z := \frac{\cos z}{\sin z},
\]

\[
(8.48b)
\]

respectively. Since sine and cosine are both continuous, tangent and cotangent are also both continuous on their respective domains. Both functions have period $\pi$, since, for each $z$ in the respective domains,

\[
\tan(z + \pi) = \frac{\sin(z + \pi)}{\cos(z + \pi)} = \frac{-\sin z}{-\cos z} = \tan z, \quad \cot(z + \pi) = \frac{-\cos z}{-\sin z} = \cot z.
\]

\[
(8.49)
\]

Since

\[
\lim_{n \to \infty} \sin \left(\frac{\pi}{2} - \frac{1}{n}\right) = \sin \frac{\pi}{2} = 1 \quad \land \quad \lim_{n \to \infty} \cos \left(\frac{\pi}{2} - \frac{1}{n}\right) = \cos \frac{\pi}{2} = 0
\]

\[
\land \quad \cos \left(\frac{\pi}{2} - \frac{1}{n}\right) > 0 \quad \Rightarrow \quad \lim_{n \to \infty} \tan \left(\frac{\pi}{2} - \frac{1}{n}\right) = \infty,
\]

\[
\lim_{n \to \infty} \sin \left(-\frac{\pi}{2} + \frac{1}{n}\right) = \sin \left(-\frac{\pi}{2}\right) = -1 \quad \land \quad \lim_{n \to \infty} \cos \left(-\frac{\pi}{2} + \frac{1}{n}\right) = \cos \left(-\frac{\pi}{2}\right) = 0
\]

\[
\land \quad \cos \left(-\frac{\pi}{2} + \frac{1}{n}\right) > 0 \quad \Rightarrow \quad \lim_{n \to \infty} \tan \left(-\frac{\pi}{2} + \frac{1}{n}\right) = -\infty,
\]

we obtain $\tan(\mathbb{R} \setminus \cos^{-1}\{0\}) = \cot(\mathbb{R} \setminus \sin^{-1}\{0\}) = \mathbb{R}$.

For each $k \in \mathbb{Z}$,

\[
\tan \text{ is strictly increasing on } \left] -\frac{\pi}{2} + k\pi, \frac{\pi}{2} + k\pi \right[, \quad (8.50a)
\]

\[
\cot \text{ is strictly decreasing on } \left] k\pi, (k + 1)\pi \right[, \quad (8.50b)
\]
On \( ]0, \frac{\pi}{2}[ \), \( \sin \) is strictly increasing and \( \cos \) is strictly decreasing, i.e. \( \tan \) is strictly increasing and \( \cot \) is strictly decreasing. Since \( \tan(-x) = \sin(-x)/\cos(-x) = -\tan(x) \), on \( ]-\frac{\pi}{2}, 0[ \), \( \tan \) is strictly increasing and \( \cot \) is strictly decreasing. Taking into account the signs of \( \tan \) and \( \cot \) on the respective intervals and their \( \pi \)-periodicity according to (8.49) proves (8.50).

Definition and Remark 8.27. Since we have seen \( \sin \) to be strictly increasing on \( ]-\frac{\pi}{2}, \frac{\pi}{2}[ \) with range \( ]-1, 1[ \), \( \cos \) to be strictly decreasing one \( [0, \pi[ \) with range \( ]-1, 1[ \), \( \tan \) to be strictly increasing on \( ]-\frac{\pi}{2}, \frac{\pi}{2}[ \) with range \( \mathbb{R} \), and \( \cot \) to be strictly decreasing one \( ]0, \pi[ \) with range \( \mathbb{R} \); and since all four functions are continuous, Th. 7.60 implies the existence of inverse functions, denoted by

\[
\text{arcsin} : ]-1, 1[ \longrightarrow ]-\pi/2, \pi/2[ , \quad (8.51a) \\
\text{arccos} : ]-1, 1[ \longrightarrow [0, \pi) \\
\text{arctan} : \mathbb{R} \longrightarrow ]-\pi/2, \pi/2[ , \quad (8.51c) \\
\text{arccot} : \mathbb{R} \longrightarrow ]0, \pi[ , \quad (8.51d)
\]

respectively, where all four inverse functions are continuous, \( \text{arcsin} \) is strictly increasing, \( \text{arccos} \) is strictly decreasing, \( \text{arctan} \) is strictly increasing, and \( \text{arccot} \) is strictly decreasing.

Of course, using (8.45) and (8.50), respectively, one can also obtain the inverse functions on different intervals, and, in the literature, such inverse functions are, indeed, considered as well. Somewhat confusingly, it is common to denote all these different functions by the same symbols, namely the ones introduced in (8.51). Here, we will not need to pursue this any further, i.e. we will only consider the inverse functions precisely as defined in (8.51), which are also known as the principle inverse functions of \( \sin \), \( \cos \), \( \tan \), and \( \cot \), respectively.

8.5 Polar Form of Complex Numbers, Fundamental Theorem of Algebra

Theorem 8.28. For each complex number \( z \in \mathbb{C} \), there exist real numbers \( r \geq 0 \) and \( \varphi \in \mathbb{R} \) such that

\[
z = r e^{i\varphi} . \quad (8.52)
\]

Moreover, if (8.52) holds with \( r \geq 0 \) and \( \varphi \in \mathbb{R} \), then \( r \) is the modulus of \( z \) and, for \( z \neq 0 \), \( \varphi \) is uniquely determined up to addition of an integer multiple of \( 2\pi \), i.e.

\[
\forall z \in \mathbb{C}\setminus\{0\} \quad \left( z = r e^{i\varphi_1} = r e^{i\varphi_2} \land r \geq 0 \Rightarrow r = |z| \land \exists_{k \in \mathbb{Z}} \varphi_1 - \varphi_2 = 2\pi k \right) . \quad (8.53)
\]

Proof. For \( z = 0 \), there is nothing to prove, so we assume \( z \neq 0 \) and set \( r := |z| \). We write \( z = x + iy \) with \( x, y \in \mathbb{R} \), first assuming \( y \geq 0 \). Then

\[
\frac{z}{r} = \xi + i\eta , \quad \text{where} \quad \xi = \frac{x}{r} , \quad \eta = \frac{y}{r} \geq 0 , \quad \xi^2 + \eta^2 = 1 . \quad (8.54)
\]
In particular, \(-1 \leq \xi \leq 1\). Thus, letting 
\[
\varphi := \arccos \xi,
\]
we obtain \(\varphi \in [0, \pi]\), \(\xi = \cos \varphi\), and \(\sin \varphi \geq 0\), yielding 
\[
\sin \varphi = \sqrt{1 - (\cos \varphi)^2} = \sqrt{1 - \xi^2} \quad (8.54) \equiv \eta.
\]
In consequence, 
\[
\bar{z} = \xi + i\eta = \cos \varphi + i \sin \varphi \quad (8.46a) \equiv e^{i\varphi},
\]
as desired. If \(y \leq 0\), then the above shows the existence of \(\psi \in \mathbb{R}\) such that \(\bar{z} = x - iy = re^{i\psi} = r \cos \psi + ir \sin \psi\). Letting \(\varphi := -\psi\), we, once again, have \(z = r \cos \psi - ir \sin \psi = re^{-i\psi} = re^{i\varphi}\), as desired, completing the existence proof for the representation (8.52). Now assume (8.52) holds with \(r \geq 0\). Then 
\[
|z| = r|e^{i\varphi}| = r \sqrt{(\sin \varphi)^2 + (\cos \varphi)^2} = r.
\]
Finally, if \(re^{i\varphi_1} = re^{i\varphi_2}\) with \(r > 0\), then \(e^{i(\varphi_1 - \varphi_2)} = 1\), i.e. \(i(\varphi_1 - \varphi_2) \in \{2k\pi i : k \in \mathbb{Z}\}\) by (8.47a).

**Definition and Remark 8.29.** The representation of \(z \in \mathbb{C}\) given by (8.52) is called its **polar form**, where \((r, \varphi)\) are also called **polar coordinates** of \(z\), \(\varphi\) is called an **argument** of \(z\). For \(z \neq 0\), one can fix the argument uniquely by the additional requirement \(\varphi \in [0, 2\pi[\) (but one also finds other choices, for example \(\varphi \in ]-\pi, \pi]\), in the literature). The above terminology is consistent with the common use of calling \((r, \varphi)\) **polar coordinates** of the vector \(z = (x, y) \in \mathbb{R}^2(= \mathbb{C})\) (in contrast to the *Cartesian coordinates* \((x, y))\), where \(r\) constitutes the distance of the point \(z = (x, y)\) from the origin \((0, 0)\) and \(\varphi\) is the angle between the vector \(z = (x, y)\) and the \(x\)-axis (cf. the three introductory paragraphs of the previous Sec. 8.4). As promised, we can now better understand the geometric interpretation of complex multiplication already described in Rem. 5.12: If \(z_1 = r_1e^{i\varphi_1}\) and \(z_2 = r_2e^{i\varphi_2}\), then \(z_1z_2 = r_1r_2e^{i(\varphi_1 + \varphi_2)}\), i.e. complex multiplication, indeed, means multiplying absolute values and adding arguments.

**Corollary 8.30.** If \(z \in \mathbb{C}\), then \(|z| = 1\) holds if, and only if, there exists \(\varphi \in \mathbb{R}\) such that \(z = e^{i\varphi}\) – in other words, the map 
\[
f : \mathbb{R} \to \{z \in \mathbb{C} : |z| = 1\}, \quad f(\varphi) := e^{i\varphi}, \quad (8.55)
\]
is surjective. Moreover \(f(\varphi_1) = f(\varphi_2)\) holds if, and only if, \(\varphi_1 - \varphi_2 = 2\pi k\) for some \(k \in \mathbb{Z}\).

**Proof.** Everything is immediate from Th. 8.28. \(\blacksquare\)

**Corollary 8.31 (Roots of Unity).** For each \(n \in \mathbb{N}\), the equation \(z^n = 1\) has precisely \(n\) distinct solutions \(\zeta_1, \ldots, \zeta_n \in \mathbb{C}\), where 
\[
\forall \ z^{\frac{k2\pi}{n}} = \cos \frac{k2\pi}{n} + i \sin \frac{k2\pi}{n} = \zeta_1^k. \quad (8.56)
\]
The numbers \(\zeta_1, \ldots, \zeta_n\) defined in (8.56) are called the \(n\)th roots of unity.
Proof. It is \( \zeta_k^n = e^{k2\pi i} = 1 \) for each \( k \in \{1, \ldots, n\} \) and the \( \zeta_1, \ldots, \zeta_n \) are all distinct by Cor. 8.30, since, for \( k, l \in \{1, \ldots, n\} \) with \( k \neq l, (k - l)/n \notin \mathbb{Z} \). As \( \zeta_1, \ldots, \zeta_n \) are \( n \) distinct zeros of the polynomial \( P : \mathbb{C} \to \mathbb{C}, P(z) := z^n - 1 \), and \( P \) has at most \( n \) zeros by Th. 6.6(a), \( \zeta_1, \ldots, \zeta_n \) constitute all solutions to \( z^n = 1 \). \( \blacksquare \)

We are now in a position to prove one of the central results of analysis and algebra, namely the fundamental theorem of algebra. The following proof does not need any tools beyond the ones provided by this class – it is actually mainly founded on continuous functions attaining a min and a max on compact sets according to Th. 7.54 and the existence of \( n \)th roots of unity according to Cor. 8.31.

**Theorem 8.32** (Fundamental Theorem of Algebra). Every polynomial \( P : \mathbb{C} \to \mathbb{C}, P(z) := \sum_{j=0}^{n} a_j z^j \), of degree \( n \geq 1 \) (i.e. \( a_0, \ldots, a_n \in \mathbb{C} \) with \( a_n \neq 0 \)) has at least one zero \( z_0 \in \mathbb{C} \).

Proof. Dividing the equation \( P(z) = 0 \) by \( a_n \neq 0 \), it suffices to consider the case \( a_n = 1 \). We therefore assume

\[
\forall \ z \in \mathbb{C} \quad P(z) = z^n + a_{n-1}z^{n-1} + \cdots + a_1z + a_0.
\]

**Claim 1.** The function \( |P| \) attains its global min on \( \mathbb{C} \), i.e. there exists \( z_0 \in \mathbb{C} \) such that \( |P| \) is minimal in \( z_0 \).

Proof. We first note

\[
\forall \ z \neq 0 \quad P(z) = z^n(1 + r(z)), \quad \text{where} \quad r(z) := \frac{a_{n-1}}{z} + \cdots + \frac{a_0}{z^n}.
\]

Set \( M := |a_0| + \cdots + |a_{n-1}| \) and \( R := \max\{1, 2M\} \).

Then

\[
\forall \ |z| \geq R \quad |r(z)| \leq \frac{M}{|z|} \quad \text{and} \quad \frac{1}{2} \leq \frac{|z|^{2M}}{2} \leq |z^n|/2 \geq M.
\]

This estimate together with \( |P(0)| = |a_0| \leq M \) shows that the min of \( |P| \) on the compact disk \( \mathcal{B}_R(0) \) (see Ex. 7.47(a)) (such a min \( z_0 \in \mathcal{B}_R(0) \) exists due to Th. 7.54) must be the global min of \( |P| \) on \( \mathbb{C} \). \( \blacktriangle \)

**Claim 2.** If \( |P| \) has a min in \( z_0 \in \mathbb{C} \), then \( P(z_0) = 0 \).

Proof. Proceeding by contraposition, we assume \( P(z_0) \neq 0 \) and show that \( |P| \) does not have a min in \( z_0 \). We need to construct \( z_1 \in \mathbb{C} \) such that \( |P(z_1)| < |P(z_0)| \). To this end, define

\[
p : \mathbb{C} \to \mathbb{C}, \quad p(z) := \frac{P(z_0 + z)}{P(z_0)}.
\]
Then $p$ is still a polynomial of degree $n$. Since $p(0) = 1$,
\[ \exists k \in \{1, \ldots, n\} \quad \exists b_0, \ldots, b_n \in \mathbb{C} \quad \forall z \in \mathbb{C} \quad p(z) = 1 + \sum_{j=k}^{n} b_j z^j, \quad b_k \neq 0. \]

Write $-b_k^{-1}$ in polar form, i.e. $-b_k^{-1} = r e^{i\varphi}$ with $r \in \mathbb{R}^+$ and $\varphi \in \mathbb{R}$. Define
\[ \beta := \sqrt[n]{r} e^{i\varphi/n} \quad (i.e. \beta^k = r e^{i\varphi} = -b_k^{-1}) \]
and
\[ q : \mathbb{C} \to \mathbb{C}, \quad q(z) := p(\beta z) = 1 + b_k \beta^k z^k + \sum_{j=k+1}^{n} b_j \beta^j z^j = 1 - z^k + z^{k+1} S(z), \]
where $S$ is the polynomial
\[ S : \mathbb{C} \to \mathbb{C}, \quad S(z) := \sum_{j=0}^{n-k-1} b_{k+1+j} \beta^{k+1+j} z^j \quad (S \equiv 0 \text{ in case } k = n). \]

Then, according to Th. 7.54,
\[ \exists C \in \mathbb{R}^+ \quad \forall z \in B_1(0) \quad |S(z)| \leq C. \]
Letting
\[ c := \min\{1, C^{-1}\}, \]
one obtains
\[ \forall 0 < |z| < c \quad |z^{k+1} S(z)| \leq C |z|^{k+1} < |z|^k \]
and, thus,
\[ \forall x \in [0, c] \quad |q(x)| \leq 1 - x^k + |x^{k+1} S(x)| < 1 - x^k + x^k = 1. \]
Thus, finally,
\[ \forall x \in [0, c] \quad \left| \frac{P(z_0 + \beta x)}{P(z_0)} \right| = \frac{|p(\beta x)|}{|p(0)|} = |q(x)| < 1, \]
showing $|P|$ does not have a min in $z_0$. \hfill \triangle

Combining Claims 1 and 2 completes the proof of the theorem. \hfill \blacksquare

**Corollary 8.33.** For every polynomial $P : \mathbb{C} \to \mathbb{C}$ of degree $n \geq 1$, there exist numbers $c, \zeta_1, \ldots, \zeta_n \in \mathbb{C}$ such that
\[ P(z) = c \prod_{j=1}^{n} (z - \zeta_j) = c(z - \zeta_1)(z - \zeta_2) \cdots (z - \zeta_n) \quad (8.57) \]
(the $\zeta_1, \ldots, \zeta_n$ are precisely all the zeros of $P$, some or all of which might be identical).

**Proof.** One just combines Th. 8.32 with Rem. 6.7. \hfill \blacksquare
9 Differential Calculus

9.1 Definition of Differentiability and Rules

The basic idea of differential calculus is to locally approximate nonlinear functions \( f \) by linear functions. In our case, \( f \) will be defined on a subset \( M \) of \( \mathbb{R} \) and, given \( \xi \in M \) and \( \mathbb{R} \)-valued \( f \), we will investigate the question if we can define a number \( f'(\xi) \in \mathbb{R} \) that represents the slope of the graph of \( f \) at \( \xi \) such that the line through \( \xi \) with slope \( f'(\xi) \) (called the tangent of \( f \) in \( \xi \)) can be considered as a local approximation of the graph of \( f \).

If such a local approximation of \( f \) in \( \xi \) is at all reasonable, then, for \( x \neq \xi \),

\[
\frac{f(x) - f(\xi)}{x - \xi}
\]

should provide “good” approximations of \( f'(\xi) \) if \( x \) tends to \( \xi \). This leads to the following Def. 9.1, where we also allow \( \mathbb{C} \)-valued functions (while the above-described geometric interpretation only works for \( \mathbb{R} \)-valued functions, it can be applied to both the real and the imaginary parts of a \( \mathbb{C} \)-valued function, cf. Rem. 9.2 below); but note that we do not consider differentiability of functions \( f : \mathbb{C} \rightarrow \mathbb{C} \), which would lead to the notion of complex differentiability or holomorphicity, which is studied in the field of Complex Analysis and is beyond the scope of this class.

**Definition 9.1.** Let \( a < b, f : ]a, b[ \rightarrow \mathbb{K} \) (\( a = -\infty, b = \infty \) is admissible), and \( \xi \in ]a, b[ \). Then \( f \) is said to be differentiable at \( \xi \) if, and only if, the following limit in (9.1) exists in the sense of Def. 8.17 (where \( x \mapsto \frac{f(x) - f(\xi)}{x - \xi} \) plays the role of \( x \mapsto f(x) \) in Def. 8.17). The limit is then called the derivative of \( f \) in \( \xi \). Many symbols are used in the literature to denote derivatives, the following provides a selection:

\[
f'(\xi) := \partial_x f(\xi) := \lim_{x \to \xi} \frac{f(x) - f(\xi)}{x - \xi} = \lim_{h \to 0} \frac{f(\xi + h) - f(\xi)}{h}.
\]

(9.1)

Note both limits occurring in (9.1) are, indeed, identical, since the sequence \((x_k)_{k \in \mathbb{N}}\) in \( ]a, b[ \) converges to \( \xi \) if, and only if, the sequence \((h_k)_{k \to \infty}\) with \( h_k := x_k - \xi \) converges to 0. The number in (9.1) (if it exists) is also called a differential quotient, whereas \( \frac{f(x) - f(\xi)}{x - \xi} \) is known as a difference quotient.

\( f \) is called differentiable if, and only if, it is differentiable at each \( \xi \in ]a, b[ \). In that case, one calls the function

\[
f' : ]a, b[ \rightarrow \mathbb{K}, \quad x \mapsto f'(x),
\]

(9.2)

the derivative of \( f \).

**Remark 9.2.** In the situation of Def. 9.1, the complex-valued function \( f : ]a, b[ \rightarrow \mathbb{C} \) is differentiable at \( \xi \in ]a, b[ \) if, and only if, both functions \( \text{Re} f, \text{Im} f : ]a, b[ \rightarrow \mathbb{R} \) are differentiable, and, in that case

\[
f'(\xi) = (\text{Re} f)'(\xi) + i (\text{Im} f)'(\xi).
\]

(9.3)
Indeed, we merely have to note
\[
\forall_{x,\xi\in[a,b], x\neq\xi} \frac{f(x) - f(\xi)}{x - \xi} = \frac{\text{Re} \ f(x) - \text{Re} \ f(\xi)}{x - \xi} + i \frac{\text{Im} \ f(x) - \text{Im} \ f(\xi)}{x - \xi} \quad (9.4)
\]
and that, by (7.2) a sequence \((z_n)_{n\in\mathbb{N}}\) in \(\mathbb{C}\) converges to \(\zeta \in \mathbb{C}\) if, and only if, both \(\lim_{n\to\infty} \text{Re} \ z_n = \text{Re} \ \zeta\) and \(\lim_{n\to\infty} \text{Im} \ z_n = \text{Im} \ \zeta\) hold.

**Definition 9.3.** If \(f : ]a, b[ \to \mathbb{R}\) as in Def. 9.1 is differentiable at \(\xi \in ]a, b[\), then the graph of the affine function
\[
L : \mathbb{R} \to \mathbb{R}, \quad L(x) := f(\xi) + f'(\xi)(x - \xi),
\]
i.e. the line through \((\xi, f(\xi))\) with slope \(f'(\xi)\) is called the tangent to the graph of \(f\) at \(\xi\).

**Theorem 9.4.** If \(f : ]a, b[ \to \mathbb{K}\) as in Def. 9.1 is differentiable at \(\xi \in ]a, b[\), then it is continuous at \(\xi\). In particular, if \(f\) is everywhere differentiable, then it is everywhere continuous.

**Proof.** Let \((x_k)_{k\in\mathbb{N}}\) be a sequence in \(]a, b[\setminus\{\xi\}\) such that \(\lim_{k\to\infty} x_k = \xi\). Then
\[
\lim_{k\to\infty} \left( f(x_k) - f(\xi) \right) = \lim_{k\to\infty} \frac{(x_k - \xi)(f(x_k) - f(\xi))}{x_k - \xi} = 0 \cdot f'(\xi) = 0, \quad (9.6)
\]
proving the continuity of \(f\) in \(\xi\).

**Example 9.5.** (a) For each \(a, b \in \mathbb{K}\), the affine function \(f : \mathbb{R} \to \mathbb{K}, f(x) := ax + b\), is differentiable with \(f'(x) = a\) for each \(x \in \mathbb{R}\): If \(x \in \mathbb{R}\) and \((h_k)_{k\in\mathbb{N}}\) is a sequence with \(h_k \neq 0\) such that \(\lim_{k\to\infty} h_k = 0\), then
\[
\lim_{k\to\infty} \frac{f(x + h_k) - f(x)}{h_k} = \lim_{k\to\infty} \frac{a(x + h_k) + b - ax - b}{h_k} = \lim_{k\to\infty} \frac{a h_k}{h_k} = a. \quad (9.7)
\]
In particular, each constant function \(f \equiv b\) has derivative \(f' \equiv 0\).

(b) For each \(c \in \mathbb{K}\), the function \(f : \mathbb{R} \to \mathbb{K}, f(x) := e^{cx}\), is differentiable with \(f'(x) = ce^{cx}\) for each \(x \in \mathbb{R}\) (in particular, \(c = 1\) yields \(f'(x) = e^x\) for \(f(x) = e^x\), and \(c = \ln a\) yields \(f'(x) = (\ln a) a^x\) for \(f(x) = a^x = e^{x \ln a}, a \in \mathbb{R}^+\)): The case \(c = 0\) was treated in (a). Thus, let \(c \neq 0\). If \(x \in \mathbb{R}\) and \((h_k)_{k\in\mathbb{N}}\) is a sequence with \(h_k \neq 0\) such that \(\lim_{k\to\infty} h_k = 0\), then
\[
\lim_{k\to\infty} \frac{f(x + h_k) - f(x)}{h_k} = \lim_{k\to\infty} \frac{e^{cx + ch_k} - e^{cx}}{h_k} = c e^{cx} \lim_{k\to\infty} \frac{e^{ch_k} - 1}{ch_k} = c e^{cx}. \quad (8.32)
\]
The sine and the cosine function $f, g : \mathbb{R} \to \mathbb{R}$, $f(x) := \sin x$, $g(x) := \cos x$, are differentiable with $f'(x) = \cos x$ and $g'(x) = -\sin x$ for each $x \in \mathbb{R}$: If $x \in \mathbb{R}$ and $(h_k)_{k \in \mathbb{N}}$ is a sequence with $h_k \neq 0$ such that $\lim_{k \to \infty} h_k = 0$, then

$$
\lim_{k \to \infty} \frac{f(x + h_k) - f(x)}{h_k} = \lim_{k \to \infty} \frac{\sin(x + h_k) - \sin x}{h_k} = \lim_{k \to \infty} \frac{\sin x \cos h_k + \cos x \sin h_k - \sin x}{h_k} = \sin x \lim_{k \to \infty} \frac{h_k(\cos h_k - 1)}{h_k^2} + \cos x \lim_{k \to \infty} \frac{\sin h_k}{h_k} = (\sin x) \cdot 0 \cdot \left( -\frac{1}{2} \right) + (\cos x) \cdot 1 = \cos x. \quad (9.9)
$$

The proof of $g'(x) = -\sin x$ is left as an exercise.

The absolute value function $f : \mathbb{R} \to \mathbb{R}$, $f(x) := |x|$, is not differentiable at $\xi = 0$:

$$
\lim_{n \to \infty} \frac{f(0 + \frac{1}{n}) - f(0)}{\frac{1}{n}} = \lim_{n \to \infty} 1 = 1, \quad \text{(9.10a)}
$$

$$
\lim_{n \to \infty} \frac{f(0 - \frac{1}{n}) - f(0)}{-\frac{1}{n}} = \lim_{n \to \infty} -\frac{1}{n} = -1, \quad \text{(9.10b)}
$$

showing that $\frac{f(0+h) - f(0)}{h}$ does not have a limit for $h \to 0$.

**Theorem 9.6.** Let $a < b$, $f, g : ]a, b[ \to \mathbb{K}$ ($a = -\infty$, $b = \infty$ is admissible), and $\xi \in ]a, b[$. Assume $f$ and $g$ are differentiable at $\xi$.

(a) For each $\lambda \in \mathbb{K}$, $\lambda f$ is differentiable at $\xi$ and $(\lambda f)'(\xi) = \lambda f'(\xi)$.

(b) $f + g$ is differentiable at $\xi$ and $(f + g)'(\xi) = f'(\xi) + g'(\xi)$.

(c) Product Rule: $fg$ is differentiable at $\xi$ and $(fg)'(\xi) = f'(\xi)g(\xi) + f(\xi)g'(\xi)$.

(d) Quotient Rule: If $g(\xi) \neq 0$, then $f/g$ is differentiable at $\xi$ and

$$
(f/g)'(\xi) = \frac{f'(\xi)g(\xi) - f(\xi)g'(\xi)}{(g(\xi))^2}, \quad \text{in particular} \quad (1/g)'(\xi) = -\frac{g'(\xi)}{(g(\xi))^2}.
$$

**Proof.** Let $(h_k)_{k \in \mathbb{N}}$ be a sequence with $h_k \neq 0$ such that $\lim_{k \to \infty} h_k = 0$.

For (a), one computes

$$
\lim_{k \to \infty} \frac{(\lambda f)(\xi + h_k) - (\lambda f)(\xi)}{h_k} = \lim_{k \to \infty} \frac{\lambda f(\xi + h_k) - \lambda f(\xi)}{h_k} = \lambda \lim_{k \to \infty} \frac{f(\xi + h_k) - f(\xi)}{h_k} = \lambda f'(\xi).
$$
For (b), one computes
\[
\lim_{k \to \infty} \frac{(f + g)(\xi + h_k) - (f + g)(\xi)}{h_k} = \lim_{k \to \infty} \frac{f(\xi + h_k) - f(\xi) + g(\xi + h_k) - g(\xi)}{h_k}
\]
which implies the general case using (c):
\[
\lim_{k \to \infty} \frac{f(\xi + h_k) - f(\xi)}{h_k} + \lim_{k \to \infty} \frac{g(\xi + h_k) - g(\xi)}{h_k} = f'(\xi) + g'(\xi).
\]

For (c), one computes
\[
\lim_{k \to \infty} \frac{(fg)(\xi + h_k) - (fg)(\xi)}{h_k}
\]
\[
= \lim_{k \to \infty} \frac{f(\xi + h_k)g(\xi + h_k) - f(\xi)g(\xi + h_k) + f(\xi)g(\xi + h_k) - f(\xi)g(\xi)}{h_k}
\]
\[
= \lim_{k \to \infty} g(\xi + h_k) \lim_{k \to \infty} \frac{f(\xi + h_k) - f(\xi)}{h_k} + f(\xi) \lim_{k \to \infty} \frac{g(\xi + h_k) - g(\xi)}{h_k}
\]
\[
= f'(\xi)g(\xi) + f(\xi)g'(\xi),
\]
where, in the last equality, we used the continuity of \(g\) in \(\xi\) according to Th. 9.4.

For (d), one first proves the special case \(f \equiv 1\) by
\[
\lim_{k \to \infty} \frac{(1/g)(\xi + h_k) - (1/g)(\xi)}{h_k} = \lim_{k \to \infty} \frac{g(\xi) - g(\xi + h_k)}{g(\xi + h_k)g(\xi)h_k} = -\frac{g'(\xi)}{(g(\xi))^2},
\]
which implies the general case using (c):
\[
(f/g)'(\xi) = \left(f \cdot \frac{1}{g}\right)'(\xi) = \frac{f'(\xi)}{g(\xi)} - \frac{f(\xi)g'(\xi)}{(g(\xi))^2} = \frac{f'(\xi)g(\xi) - f(\xi)g'(\xi)}{(g(\xi))^2},
\]
completing the proof.

**Example 9.7.** (a) Each polynomial is differentiable and the derivative is, again, a polynomial. More precisely,
\[
P : \mathbb{R} \to \mathbb{K}, \quad P(x) = \sum_{j=0}^{n} a_j x^j, \quad a_j \in \mathbb{K}
\]
\[\Rightarrow \quad P' : \mathbb{R} \to \mathbb{K}, \quad P'(x) = \sum_{j=1}^{n} j a_j x^{j-1} ;
\]
(9.11)
The cases \(n = 0, 1\) are provided by Ex. 9.5(a). To complete the induction proof of (9.11), we carry out the induction step for each \(n \in \mathbb{N}\): Writing \(P(x) = \sum_{j=0}^{n} a_j x^j + a_{n+1} x x^n\) and applying the induction hypothesis as well as the rules of Th. 9.6 yields
\[
P'(x) = \sum_{j=1}^{n} j a_j x^{j-1} + a_{n+1} (1 \cdot x^n + x \cdot n \cdot x^{n-1}) = \sum_{j=1}^{n+1} j a_j x^{j-1},
\]
which establishes the case.
(b) Clearly, the derivatives of rational functions \( P/Q \) with polynomials \( P \) and \( Q \) can be computed from (9.11) and the quotient rule of Th. 9.6(d).

(c) The functions \( \tan \) and \( \cot \) as defined in (8.48) and restricted to \( \mathbb{R} \setminus \cos^{-1}\{0\} \) and \( \mathbb{R} \setminus \sin^{-1}\{0\} \), respectively, are differentiable and one obtains

\[
\begin{align*}
\tan' &: \mathbb{R} \setminus \cos^{-1}\{0\} \to \mathbb{R}, \quad \tan' x = \frac{1}{(\cos x)^2} = 1 + (\tan x)^2, \quad (9.12a) \\
\cot' &: \mathbb{R} \setminus \sin^{-1}\{0\} \to \mathbb{R}, \quad \cot' x = -\frac{1}{(\sin x)^2} = -(1 + (\cot x)^2). \quad (9.12b)
\end{align*}
\]

One merely needs the derivatives of \( \sin \) and \( \cos \) from Ex. 9.5(c) and the quotient rule of Th. 9.6(d):

\[
\begin{align*}
\tan' x &= \frac{\cos x \cos x - \sin x(-\sin x)}{(\cos x)^2} = \frac{1}{(\cos x)^2} = 1 + (\tan x)^2, \\
\cot' x &= \frac{-\sin x \sin x - \cos x \cos x}{(\sin x)^2} = -\frac{1}{(\sin x)^2} = -(1 + (\cot x)^2).
\end{align*}
\]

**Theorem 9.8** (Derivative of Inverse Functions). Let \( a < b, I := ]a,b[ \) (\( a = -\infty, b = \infty \) is admissible). If \( f : I \to \mathbb{R} \) is differentiable and strictly increasing (resp. decreasing), then \( f \) has a continuous, strictly increasing (resp. decreasing) inverse function \( f^{-1} \) defined on the interval \( J := f(I) \), i.e. \( f^{-1} : J \to I \), and, for each \( \xi \in I \) with \( f'(\xi) \neq 0 \), \( f^{-1} \) is differentiable at \( \eta := f(\xi) \) with

\[
(f^{-1})'(\eta) = \frac{1}{f'(\xi)} = \frac{1}{f'(f^{-1}(\eta))}. \quad (9.13)
\]

**Proof.** As a differentiable function, \( f \) is continuous by Th. 9.4, i.e. Th. 7.60 provides all the present assertions, except differentiability at \( \eta \) and (9.13). Let \( (y_k)_{k \in \mathbb{N}} \) be a sequence in \( J \setminus \{\eta\} \) such that \( \lim_{k \to \infty} y_k = \eta \). Then, as \( f^{-1} \) is bijective and continuous, \( (f^{-1}(y_k))_{k \in \mathbb{N}} \) is a sequence in \( I \setminus \{\xi\} \) such that \( \lim_{k \to \infty} f^{-1}(y_k) = \xi \), and one obtains

\[
\lim_{k \to \infty} \frac{f^{-1}(y_k) - f^{-1}(\eta)}{y_k - \eta} = \lim_{k \to \infty} \frac{f^{-1}(y_k) - f^{-1}(\eta)}{f(f^{-1}(y_k)) - f(f^{-1}(\eta))} = \frac{1}{f'(f^{-1}(\eta))}, \quad (9.14)
\]

establishing the case. \( \blacksquare \)

**Example 9.9.** (a) The function \( \ln : \mathbb{R}^+ \to \mathbb{R} \) is differentiable and, for each \( x \in \mathbb{R}^+ \), \( \ln' x = 1/x \): If \( f(x) = e^x \), then \( f'(x) = e^x \neq 0 \) for each \( x \in \mathbb{R} \), \( \ln x = f^{-1}(x) \), and (9.13) yields

\[
\ln' x = \frac{1}{f'(\ln x)} = \frac{1}{e^{\ln x}} = \frac{1}{x}.
\]
(b) The function \( \arcsin : [-1,1] \rightarrow [-\pi/2,\pi/2] \) is differentiable and, for each \( x \in [-1,1] \), \( \arcsin' x = 1/\sqrt{1-x^2} \): if \( f(x) = \sin x \), then \( f'(x) = \cos x \neq 0 \) for each \( x \in [-\pi/2,\pi/2] \), \( \arcsin x = f^{-1}(x) \), and (9.13) yields

\[
\arcsin' x = \frac{1}{f'(\arcsin x)} = \frac{1}{\cos \arcsin x} = \frac{1}{\sqrt{1 - (\sin \arcsin x)^2}} = \frac{1}{\sqrt{1 - x^2}},
\]

where, at (*), it was used that \( \cos^2 = 1 - \sin^2 \) and \( \cos t > 0 \) for each \( t \in [-\pi/2,\pi/2] \).

(c) The function \( \arccos : [-1,1] \rightarrow [0,\pi] \) is differentiable and, for each \( x \in [-1,1] \), \( \arccos' x = -1/\sqrt{1-x^2} \): if \( f(x) = \cos x \), then \( f'(x) = -\sin x \neq 0 \) for each \( x \in [0,\pi] \), \( \arccos x = f^{-1}(x) \), and (9.13) yields

\[
\arccos' x = \frac{1}{f'(\arccos x)} = \frac{1}{-\sin \arccos x} = \frac{1}{\sqrt{1 - (\cos \arccos x)^2}} = \frac{1}{\sqrt{1 - x^2}},
\]

where, at (*), it was used that \( \sin^2 = 1 - \cos^2 \) and \( \sin t > 0 \) for each \( t \in [0,\pi] \).

(d) The function \( \arctan : \mathbb{R} \rightarrow (-\pi/2,\pi/2] \) is differentiable and, for each \( x \in \mathbb{R} \), \( \arctan' x = 1/(1 + x^2) \): Apply Th. 9.8 with \( f(x) = \tan x \) as an exercise.

(e) The function \( \arccot : \mathbb{R} \rightarrow [0,\pi] \) is differentiable and, for each \( x \in \mathbb{R} \), \( \arccot' x = -1/(1 + x^2) \): Apply Th. 9.8 with \( f(x) = \cot x \) as an exercise.

**Theorem 9.10 (Chain Rule).** Let \( a < b, c < d, f : [a,b] \rightarrow \mathbb{R}, g : [c,d] \rightarrow \mathbb{K}, f([a,b]) \subseteq [c,d], (a,c = -\infty, b,d = \infty \text{ is admissible}) \). If \( f \) is differentiable in \( \xi \in [a,b] \) and \( g \) is differentiable in \( f(\xi) \in [c,d] \), then \( g \circ f : [a,b] \rightarrow \mathbb{K} \) is differentiable in \( \xi \) and

\[
(g \circ f)'(\xi) = f'(\xi)g'(f(\xi)).
\]

**Proof.** Let \( \eta := f(\xi) \) and define the auxiliary function

\[
\tilde{g} : [c,d] \rightarrow \mathbb{K}, \quad \tilde{g}(x) := \begin{cases} \frac{g(x) - g(\eta)}{x - \eta} & \text{for } x \neq \eta, \\ g'(x) & \text{for } x = \eta. \end{cases}
\]

Then

\[
\forall \ x \in [c,d], \quad g(x) - g(\eta) = \tilde{g}(x)(x - \eta).
\]

Let \( (x_k)_{k \in \mathbb{N}} \) be a sequence in \( [a,b] \setminus \{\xi\} \) such that \( \lim_{k \to \infty} x_k = \xi \). One obtains

\[
\lim_{k \to \infty} \frac{g(f(x_k)) - g(f(\xi))}{x_k - \xi} = \lim_{k \to \infty} \frac{\tilde{g}(f(x_k))(f(x_k) - f(\xi))}{x_k - \xi} = \lim_{k \to \infty} \tilde{g}(f(x_k)) \lim_{k \to \infty} \frac{f(x_k) - f(\xi)}{x_k - \xi} = f'(\xi)g'(f(\xi)),
\]

establishing the case.
Example 9.11. (a) According to the chain rule of Th. 9.10, the function \( h : \mathbb{R} \rightarrow \mathbb{R}, h(x) := \sin(-x^3) \) is differentiable and, for each \( x \in \mathbb{R}, h'(x) = -3x^2 \cos(-x^3) \).

(b) According to the chain rule of Th. 9.10, each power function \( h : \mathbb{R}^+ \rightarrow \mathbb{R}, h(x) := x^\alpha = e^{\alpha \ln x}, \alpha \in \mathbb{R}, \) is differentiable and, for each \( x \in \mathbb{R}^+, h'(x) = \frac{\alpha}{x} e^{\alpha \ln x} = \alpha x^{\alpha - 1}. \)

Indeed, \( h = g \circ f \), where \( f : \mathbb{R}^+ \rightarrow \mathbb{R}, f(x) := \ln x \) with \( f' : \mathbb{R}^+ \rightarrow \mathbb{R}, f'(x) := \frac{1}{x}, \)

according to Ex. 9.5(b), and \( g : \mathbb{R} \rightarrow \mathbb{R}, g(x) := e^{\alpha x}, \) with \( g' : \mathbb{R} \rightarrow \mathbb{R}, g'(x) := \alpha e^{\alpha x} \) according to Ex. 9.9(a).

9.2 Higher Order Derivatives and the Sets \( C^k \)

Definition 9.12. Let \( a < b, I := [a, b[, f : I \rightarrow \mathbb{R} \) \((a = -\infty, b = \infty \text{ is admissible})\). If \( f \) is differentiable, then \( f' \) might or might not itself be differentiable. If \( f' \) is differentiable, then its derivative is denoted by \( f'' \) and is called the second derivative of \( f \). Clearly, this process can be iterated, leading to the following general recursive definition of higher-order derivatives:

Let \( f^{(0)} := f \). For \( k \in \mathbb{N}_0 \) assume the \( k \)th derivative of \( f \), denoted by \( f^{(k)} \) exists on \( I \). Then \( f \) is said to have a derivative of order \( k + 1 \) at \( \xi \in I \) if, and only if, \( f^{(k)} \) is differentiable at \( \xi \). In that case, define

\[
  f^{(k+1)}(\xi) := (f^{(k)})'(\xi). \tag{9.19}
\]

If \( f^{(k+1)}(\xi) \) exists for all \( \xi \in I \), then \( f \) is said to be \((k+1)\)-times differentiable and the function \( f^{(k+1)} : I \rightarrow \mathbb{R}, x \mapsto f^{(k+1)}(x) \), is called the \((k+1)\)st derivative of \( f \). It is common to write \( f' := f^{(1)}, f'' := f^{(2)}, f''' := f^{(3)}, \) but \( f^{(k)} \) if \( k \geq 4 \).

If \( f^{(k)} \) exists, it might or might not be continuous (cf. Ex. 9.13(c) below). One defines

\[
  \forall \; k \in \mathbb{N}_0, \quad C^k(I, \mathbb{K}) := \left\{ f \in \mathcal{F}(I, \mathbb{K}) : f^{(k)} \text{ exists and is continuous on } I \right\}; \tag{9.20}
\]

\[
  C^\infty(I, \mathbb{K}) := \bigcap_{k \in \mathbb{N}_0} C^k(I, \mathbb{K}) \tag{9.21}
\]

(note \( C^0(I, \mathbb{K}) = C(I, \mathbb{K}) \) and \( C(I, \mathbb{K}) \supseteq C^1(I, \mathbb{K}) \supseteq C^2(I, \mathbb{K}) \supseteq \ldots \)). Finally, we define the notation \( C^k(I) := C^k(I, \mathbb{R}) \) for \( k \in \mathbb{N}_0 \cup \{ \infty \} \).

Example 9.13. (a) One has \( \sin \in C^\infty(\mathbb{R}) \) with \( \sin' = \cos, \sin'' = -\sin, \sin''' = -\cos, \)

\( \sin^{(4)} = \sin, \ldots \).

(b) A simple induction shows, for each polynomial \( P : \mathbb{R} \rightarrow \mathbb{K}, P(x) = \sum_{j=0}^n a_j x^j,\) \( a_j \in \mathbb{K}, n \in \mathbb{N}_0 \), that \( P^{(n)}(x) = n! a_n. \) In particular, \( P \in C^\infty(\mathbb{R}, \mathbb{K}) \).

(c) It is an exercise to show the following function \( f \) is differentiable, but \( f' \) is not continuous, i.e. \( f \notin C^1(\mathbb{R}) \):

\[
  f : \mathbb{R} \rightarrow \mathbb{R}, \quad f(x) := \begin{cases} x^2 \cos \left( \frac{1}{x} \right) & \text{for } x \neq 0, \\ 0 & \text{for } x = 0. \end{cases} \]
9.3 Mean Value Theorem, Monotonicity, and Extrema

**Theorem 9.14.** Let \( a < b \). If \( f : [a, b] \rightarrow \mathbb{R} \) is differentiable in \( \xi \in ]a, b[ \) and \( f \) has a local min or max in \( \xi \), then \( f'(\xi) = 0 \).

**Proof.** Suppose \( f \) has a local max at \( \xi \). Then there exists \( \epsilon > 0 \) such that \( |h| < \epsilon \) implies \( f(\xi + h) - f(\xi) \leq 0 \). Now let \( (h_k)_{k \in \mathbb{N}} \) be a sequence in \( ]0, \epsilon[ \) with \( \lim_{k \to \infty} h_k = 0 \). Then \( f(\xi + h_k) - f(\xi) \leq 0 \) for all \( k \in \mathbb{N} \) implies

\[
f'(\xi) = \lim_{k \to \infty} \frac{f(\xi + h_k) - f(\xi)}{h_k} \leq 0,
\]

showing \( f'(\xi) = 0 \). Now, if \( f \) has a local min at \( \xi \), then \( -f \) has a local max at \( \xi \), and \( f'(\xi) = -(-f'(\xi)) = 0 \) establishes the case. \( \square \)

**Remark 9.15.** For \( f : \mathbb{R} \rightarrow \mathbb{R}, f(x) := x^3 \), it is \( f'(0) = 0 \), but \( f \) does not have a local min or max at \( 0 \), showing that, while being necessary for an differentiable function \( f \) to have a local extremum at \( \xi \), \( f'(\xi) = 0 \) is not a sufficient condition for such an extremum at \( \xi \). Points \( \xi \) with \( f'(\xi) = 0 \) are sometimes called stationary or critical points of \( f \).

Now, we first prove an important special case of the mean value theorem:

**Theorem 9.16** (Rolle’s Theorem). Let \( a < b \). If \( f : [a, b] \rightarrow \mathbb{R} \) is continuous on the compact interval \([a, b]\), differentiable on the open interval \(]a, b[\), and \( f(a) = f(b) \), then there exists \( \xi \in ]a, b[ \) such that \( f'(\xi) = 0 \).

**Proof.** If \( f \) is constant, then \( f'(\xi) = 0 \) holds for each \( \xi \in ]a, b[ \). If \( f \) is nonconstant, then there exists \( x \in ]a, b[ \) with \( f(x) \neq f(a) \). If \( f(x) > f(a) \), then Th. 7.54 implies the existence of \( \xi \in ]a, b[ \) such that \( f \) attains its (global and, thus, local) max in \( \xi \). Then Th. 9.14 yields \( f'(\xi) = 0 \). The case \( f(x) < f(a) \) is treated analogously. \( \square \)

**Theorem 9.17** (Mean Value Theorem). Let \( a < b \). If \( f, g : [a, b] \rightarrow \mathbb{R} \) are continuous on the compact interval \([a, b]\), differentiable on the open interval \(]a, b[\), and \( g'(x) \neq 0 \) for each \( x \in ]a, b[ \), then there exists \( \xi \in ]a, b[ \) such that

\[
\frac{f(b) - f(a)}{g(b) - g(a)} = \frac{f'(\xi)}{g'(
\xi)}.
\]

(9.23a)

In the special case that \( g : [a, b] \rightarrow \mathbb{R}, g(x) = x \), one obtains the standard form

\[
\frac{f(b) - f(a)}{b - a} = f'(\xi).
\]

(9.23b)

**Proof.** First note that Rolle’s Th. 9.16 and \( g' \neq 0 \) imply \( g(b) - g(a) \neq 0 \). Next, one applies Rolle’s Th. 9.16 to the auxiliary function

\[
h : [a, b] \rightarrow \mathbb{R}, \quad h(x) := f(x) - (g(x) - g(a)) \frac{f(b) - f(a)}{g(b) - g(a)}.
\]

(9.24)
Since \( f \) and \( g \) are continuous on \([a, b]\) and differentiable on \(]a, b[\), so is \( h \). Moreover, \( h(a) = f(a) = h(b) \), i.e. Rolle’s Th. 9.16 applies and yields \( \xi \in ]a, b[ \) satisfying \( h'(\xi) = 0 \). However, (9.24) implies \( h'(\xi) = 0 \) is equivalent to (9.23a).

**Corollary 9.18.** Let \( c < d \) and \( f : ]c, d[\to \mathbb{R} \) be differentiable (\( c = -\infty \), \( d = \infty \) is admissible).

(a) If \( f' \geq 0 \) (resp. \( f' \leq 0 \)), then \( f \) is increasing (resp. decreasing). Moreover, if the inequalities are strict, then the monotonicity of \( f \) is strict as well.

(b) If \( f' \equiv 0 \), then \( f \) is constant.

**Proof.** If \( c < a < b < d \) and \( f' \geq 0 \) (resp. \( f' \leq 0 \), resp. \( f' \equiv 0 \)), then (9.23b) implies \( f(b) \geq f(a) \) (resp. \( f(b) \leq f(a) \), resp. \( f(b) = f(a) \)). Moreover, strict inequalities for \( f' \) yield strict inequality between \( f(b) \) and \( f(a) \).

**Lemma 9.19.** Let \( a < b \), \( f : ]a, b[\to \mathbb{R} \), \( \xi \in ]a, b[ \), and assume \( f \) is differentiable at \( \xi \). If \( f'(\xi) > 0 \) (resp. \( f'(\xi) < 0 \)), then there exists \( \epsilon > 0 \) such that \( \xi - \epsilon, \xi + \epsilon \subset ]a, b[ \) and

\[
\forall a_1 \in ]\xi - \epsilon, \xi[ \quad \forall b_1 \in ]\xi, \xi + \epsilon[ \quad f(a_1) < f(\xi) < f(b_1) \quad (\text{resp. } f(a_1) > f(\xi) > f(b_1)).
\]

**Proof.** If there does not exist \( \epsilon > 0 \) such that \( f(a_1) < f(\xi) < f(b_1) \) for each \( a_1 \in ]\xi - \epsilon, \xi[ \) and each \( b_1 \in ]\xi, \xi + \epsilon[ \), then there exists a sequence \( (x_k)_{k \in \mathbb{N}} \) in \( ]a, b[ \) such that \( \lim_{k \to \infty} x_k = \xi \) and

\[
\forall k \in \mathbb{N} \quad \frac{f(x_k) - f(\xi)}{x_k - \xi} \leq 0,
\]

showing \( f'(\xi) \leq 0 \). Analogously, one obtains that \( f'(\xi) \geq 0 \) provided there does not exist \( \epsilon > 0 \) such that \( f(a_1) > f(\xi) > f(b_1) \) for each \( a_1 \in ]\xi - \epsilon, \xi[ \) and each \( b_1 \in ]\xi, \xi + \epsilon[ \).

**Theorem 9.20** (Sufficient Conditions for Extrema). Let \( c < d \), let \( f : ]c, d[\to \mathbb{R} \) be differentiable, and assume \( f'(\xi) = 0 \) for some \( \xi \in ]c, d[ \).

(a) If \( f'(x) > 0 \) for each \( x \in ]\xi, d[ \) and \( f'(x) < 0 \) for each \( x \in ]c, \xi[ \), then \( f \) has a strict max at \( \xi \). Likewise, if \( f''(\xi) \) exists and is negative, then \( f \) has a strict max at \( \xi \).

(b) If \( f'(x) < 0 \) for each \( x \in ]\xi, d[ \) and \( f'(x) > 0 \) for each \( x \in ]c, \xi[ \), then \( f \) has a strict min at \( \xi \). Likewise, if \( f''(\xi) \) exists and is positive, then \( f \) has a strict min at \( \xi \).

**Proof.** We just present the proof for (a); (b) is proved analogously. If \( f'(x) > 0 \) for each \( x \in ]\xi, d[ \), then (9.23b) shows \( f(\xi) - f(a) > 0 \) for each \( c < a < \xi \); analogously, if \( f'(x) < 0 \) for each \( x \in ]\xi, d[ \), then (9.23b) shows \( f(\xi) - f(b) > 0 \) for each \( \xi < b < d \). Altogether, we have shown \( f \) to have a strict max at \( \xi \). If \( f''(\xi) \) exists and is negative, then Lem. 9.19 yields the existence of \( \epsilon > 0 \) such that \( f' \) is positive on \( ]\xi - \epsilon, \xi[ \) and negative on \( ]\xi, \xi + \epsilon[ \). Applying what we have already proved with \( c := \xi - \epsilon \) and \( d := \xi + \epsilon \) establishes the case.
Example 9.21. One obtains
\[
\begin{align*}
    f : \mathbb{R} &\to \mathbb{R}, \quad f(x) := x e^x, \quad (9.25a) \\
    f' : \mathbb{R} &\to \mathbb{R}, \quad f'(x) = e^x + x e^x = (1 + x) e^x, \quad (9.25b) \\
    f'' : \mathbb{R} &\to \mathbb{R}, \quad f''(x) = 2e^x + xe^x = (2 + x) e^x. \quad (9.25c)
\end{align*}
\]

From Th. 9.14, we know that \( f \) can have at most one extremum, namely at \( \xi = -1 \), where \( f'(\xi) = 0 \). Since \( f''(\xi) = e^{-x} > 0 \), Th. 9.20(b) implies that \( f \) has a strict min at \(-1\).

9.4 L'Hôpital's Rule

L'Hôpital's rule is a result that can help to determine (function) limits (cf. Def. 8.17).

Theorem 9.22 (L'Hôpital's Rule). Let \( \xi \in \mathbb{R} \) and either \( I = ]a,\xi[ \) with \( a < \xi \) or \( I := ]\xi,b[ \) with \( \xi < b \). Moreover, assume \( f,g : I \to \mathbb{R} \) are differentiable, \( g'(x) \neq 0 \) for each \( x \in I \), and one of the following two conditions (a), (b) is satisfied:

(a) \( \lim_{x \to \xi} f(x) = \lim_{x \to \xi} g(x) = 0 \).

(b) \( \lim_{x \to \xi} g(x) = \infty \) or \( \lim_{x \to \xi} g(x) = -\infty \), where Def. 8.17 is extended to the case \( \eta \in \{-\infty, \infty\} \) in the obvious way.

Then
\[
\lim_{x \to \xi} \frac{f'(x)}{g'(x)} = \eta \implies \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta. \quad (9.26)
\]

The above statement also holds for \( \xi \in \{-\infty, \infty\} \) and/or \( \eta \in \{-\infty, \infty\} \) if, as in (b), one extends Def. 8.17 to these cases in the obvious way.

Proof. First, we assume (a). Consider the case \( \xi \in \mathbb{R} \). Since \( f \) and \( g \) are continuous, (a) implies \( f \) and \( g \) remain continuous, if we extend them to \( \xi \) by letting \( f(\xi) := g(\xi) = 0 \). This extension will now allow us to apply Th. 9.17 to \( f \) and \( g \). To prove (9.26), let \( (x_k)_{k \in \mathbb{N}} \) be a sequence in \( I \) with \( \lim_{k \to \infty} x_k = \xi \). Then (9.23a) yields, for each \( k \in \mathbb{N} \), some \( \xi_k \in ]x_k, \xi[ \) if \( x_k < \xi \) and some \( \xi_k \in ]\xi, x_k[ \) if \( \xi < x_k \), satisfying
\[
\frac{f(x_k)}{g(x_k)} = \frac{f(x_k) - f(\xi)}{g(x_k) - g(\xi)} = \frac{f'(\xi_k)}{g'(\xi_k)}. \quad (9.27)
\]

From the Sandwich Th. 7.16, we obtain \( \lim_{k \to \infty} \xi_k = \xi \), i.e. (9.27) and \( \lim_{x \to \xi} \frac{f'(x)}{g'(x)} = \eta \) imply \( \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta \) (also for \( \eta \in \{-\infty, \infty\} \)). Now consider the case \( \xi \in \{-\infty, \infty\} \) and let \( (x_k)_{k \in \mathbb{N}} \) be as before. If \( \xi = \infty \), then choose \( 1 \leq c \in I \) and set \( \tilde{I} := ]0, e^{-1}[ \); if \( \xi = -\infty \), then choose \( -1 \geq c \in I \) and set \( \tilde{I} := ]e^{-1}, 0[ \). We apply what we have already proved above to the auxiliary functions
\[
\tilde{f} : \tilde{I} \to \mathbb{R}, \quad \tilde{f}(x) := f(1/x), \quad \tilde{g} : \tilde{I} \to \mathbb{R}, \quad \tilde{g}(x) := g(1/x)
\]
at \( \tilde{x} := 0 \). From the chain rule (9.15), we know \( \tilde{f}'(x) = -f'(1/x^2) \) and \( \tilde{g}'(x) = -g'(1/x^2) \) for each \( x \in \tilde{I} \). Thus, \( \lim_{x \to \xi} \frac{f'(x)}{g'(x)} = \eta \) implies,

\[
\eta = \lim_{k \to \infty} \frac{f'(x_k)}{g'(x_k)} = \frac{\lim_{k \to \infty} -x_k^2 f'(x_k)}{\lim_{k \to \infty} -x_k^2 g'(x_k)} = \frac{\lim_{k \to \infty} \tilde{f}'(1/x_k)}{\lim_{k \to \infty} \tilde{g}'(1/x_k)} = \lim_{k \to \infty} \frac{\tilde{f}'(1/x_k)}{\tilde{g}'(1/x_k)} = \lim_{k \to \infty} f(x_k),
\]

proving \( \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta \).

We now assume (b), still letting \((x_k)_{k \in \mathbb{N}}\) be as before. Note that \( g' \neq 0 \) implies \( g \) is injective by Rolle’s Th. 9.16. Then the intermediate value theorem implies \( g \) is either strictly increasing or strictly decreasing. We proceed with the proof for the case \( I = [a, \xi[ \), the proof for \( I = ]\xi, b] \) can be done completely analogous. We first consider the case where \( g \) is strictly increasing, i.e. \( \lim_{x \to \xi} g(x) = \infty \). Assume \( \eta \in \mathbb{R} \) and \( \epsilon > 0 \). Then \( \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta \) and \( \lim_{x \to \xi} g(x) = \infty \) imply

\[
\exists \ c \in [a, \xi[ \quad \forall \ c \in [c, \xi[ \left( g(x) > 0 \land \eta - \frac{\epsilon}{2} < \frac{f(x)}{g(x)} < \eta + \frac{\epsilon}{2} \right).
\]

Since \( \lim_{k \to \infty} x_k = \xi \), there exists \( N_0 \in \mathbb{N} \) such that, for each \( k > N_0 \), \( c < x_k < \xi \). Next, according to Th. 9.17,

\[
\forall \ k > N_0 \quad \exists \ \xi_k \in ]c, x_k[ \quad \eta - \frac{\epsilon}{2} < \frac{f(x_k) - f(c)}{g(x_k) - g(c)} = \frac{f'(\xi_k)}{g'(\xi_k)} < \eta + \frac{\epsilon}{2}.
\]

In consequence, using \( g(x_k) > g(c) \), as \( g \) is strictly increasing,

\[
\forall \ k > N_0 \quad \left( \eta - \frac{\epsilon}{2} \right) f(x_k) - g(c) < f(x_k) - f(c) < \left( \eta + \frac{\epsilon}{2} \right) g(x_k) - g(c)
\]

and

\[
\forall \ k > N_0 \quad \left( \eta - \frac{\epsilon}{2} \right) f(c) - \left( \eta - \frac{\epsilon}{2} \right) g(c) < f(x_k) - f(c) < \left( \eta + \frac{\epsilon}{2} \right) f(c) - \left( \eta + \frac{\epsilon}{2} \right) g(c).
\]

Since \( \lim_{k \to \infty} g(x_k) = \infty \),

\[
\exists \ N \geq N_0 \quad \forall \ k > N \left( \frac{f(c) - (\eta - \frac{\epsilon}{2}) g(c)}{g(x_k)} < \frac{\epsilon}{2} \land \frac{f(c) - (\eta + \frac{\epsilon}{2}) g(c)}{g(x_k)} < \frac{\epsilon}{2} \right),
\]

that means

\[
\forall \ k > N \quad \eta - \epsilon < \frac{f(x_k)}{g(x_k)} < \eta + \epsilon,
\]

proving \( \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta \). For \( \eta = \infty \) and given \( n \in \mathbb{N} \), the argument is similar:

\( \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta \) and \( \lim_{x \to \xi} g(x) = \infty \) imply

\[
\exists \ c \in [a, \xi[ \quad \forall \ x \in ]c, \xi[ \left( g(x) > 0 \land n < \frac{f'(x)}{g'(x)} \right).
\]
As before, since \( \lim_{k \to \infty} x_k = \xi \), there exists \( N_0 \in \mathbb{N} \) such that, for each \( k > N_0 \), 
\( c < x_k < \xi \). Again, according to Th. 9.17,
\[
\forall \ k > N_0 \quad \exists \ \xi_k \in ]c,x_k[ \quad n < \frac{f(x_k) - f(c)}{g(x_k) - g(c)} = \frac{f'(\xi_k)}{g'(\xi_k)}. \]
In consequence, using \( g(x_k) > g(c) \), as \( g \) is strictly increasing,
\[
\forall \ k > N_0 \quad n (g(x_k) - g(c)) < f(x_k) - f(c)
\]
and
\[
\forall \ k > N_0 \quad n + \frac{f(c) - n g(c)}{g(x_k)} < \frac{f(x_k)}{g(x_k)}.
\]
Since \( \lim_{k \to \infty} g(x_k) = \infty \),
\[
\exists \ N \geq N_0 \quad \forall \ k > N \quad \left| \frac{f(c) - n g(c)}{g(x_k)} \right| < 1,
\]
that means
\[
\forall \ k > N \quad n - 1 < \frac{f(x_k)}{g(x_k)},
\]
proving \( \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta \). If \( \eta = -\infty \), then, using what we have already shown,
\[
\lim_{x \to \xi} \frac{f'(x)}{g'(x)} = \eta \quad \Rightarrow \quad \lim_{x \to \xi} -\frac{f'(x)}{g'(x)} = \infty = \lim_{x \to \xi} -\frac{f(x)}{g(x)} \quad \Rightarrow \quad \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta.
\]
Finally, if \( g \) strictly decreasing, then \(-g\) is strictly increasing and we obtain
\[
\lim_{x \to \xi} \frac{f'(x)}{g'(x)} = \eta \quad \Rightarrow \quad \lim_{x \to \xi} -\frac{f'(x)}{g'(x)} = -\eta = \lim_{x \to \xi} \frac{f(x)}{g(x)} \quad \Rightarrow \quad \lim_{x \to \xi} \frac{f(x)}{g(x)} = \eta,
\]
concluding the proof. \( \square \)

**Example 9.23.** (a) Applying L'Hôpital’s rule to \( f : ]-\pi/2, \pi/2[ \to \mathbb{R}, f(x) := \tan x, \) 
\( g : ]-\pi/2, \pi/2[ \to \mathbb{R}, g(x) := e^x - 1 \), with \( \xi = 0 \) yields
\[
\lim_{x \to 0} \frac{\tan x}{e^x - 1} = \lim_{x \to 0} \frac{\frac{\tan x}{x} - \frac{1}{x^2}}{e^x \frac{1}{e^x} - \frac{1}{x^2}} = \frac{1}{1} = 1
\] (9.28)
(note \( g'(x) = e^x \neq 0 \) for each \( x \in ]-\pi/2, \pi/2[ \)).

(b) It can happen that a single application of L'Hôpital’s rule does not, yet, yield a useful result, but that a repeated application does. An example is provided by considering \( \alpha > 0, n \in \mathbb{N} \), and \( f : \mathbb{R}^+ \to \mathbb{R}, f(x) := e^{\alpha x}, g : \mathbb{R}^+ \to \mathbb{R}, \) 
g(\( x \)) := x^n, \( \xi := \infty \). Applying L'Hôpital’s rule \( n \) times yields
\[
\lim_{x \to \infty} \frac{e^{\alpha x}}{x^n} = \lim_{x \to \infty} \frac{\alpha^n x^n e^{\alpha x}}{n!} = \infty
\] (9.29)
(note \( g^{(k)}(x) = n(n - 1) \cdots (n - k + 1)x^{n-k} \neq 0 \) for each \( k \in \{1, \ldots, n\} \) and each \( x \in \mathbb{R}^+ \)).
(c) It can also happen that even repeated applications of L'Hôpital's rule do not help at all, even though \( \lim_{x \to \xi} \frac{f(x)}{g(x)} \) does exist and the hypotheses of Th. 9.22 are all satisfied. A simple example is given by \( f : \mathbb{R} \to \mathbb{R} \), \( f(x) := e^x \), \( g : \mathbb{R} \to \mathbb{R} \), \( g(x) := 2e^x \), and \( \xi = -\infty \). Even though \( \lim_{x \to -\infty} \frac{f(x)}{g(x)} = \frac{1}{2} \), one has
\[
\lim_{x \to -\infty} f^{(n)}(x) = \lim_{x \to -\infty} g^{(n)}(x) = 0 \text{ for every } n \in \mathbb{N}.
\]

\section{The Riemann Integral on Intervals in \( \mathbb{R} \)}

\subsection{Definition and Simple Properties}

Given a nonnegative function \( f : M \to \mathbb{R}_0^+ \), \( M \subseteq \mathbb{R} \), we aim to compute the area \( \int_M f \) of the set “under the graph” of \( f \), i.e. of the set
\[
\{(x, y) \in \mathbb{R}^2 : x \in M \text{ and } 0 \leq y \leq f(x)\}.
\] (10.1)

This area \( \int_M f \) (if it exists) will be called the integral of \( f \) over \( M \). Moreover, for functions \( f : M \to \mathbb{R} \) that are not necessarily nonnegative, we would like to count areas of sets of the form (10.1) (which are below the graph of \( f \) and above the set \( M \approx \{(x, 0) \in \mathbb{R}^2 : x \in M\} \subseteq \mathbb{R}^2 \)) with a positive sign, and whereas we would like to count areas of sets above the graph of \( f \) and below the set \( M \) with a negative sign. In other words, making use of the positive and negative parts \( f^+ \) and \( f^- \) of \( f = f^+ - f^- \) as defined in (6.1i) and (6.1j), respectively, we would like our integral to satisfy
\[
\int_M f = \int_M f^+ - \int_M f^-.
\] (10.2)

Difficulties arise from the fact that both the function \( f \) and the set \( M \) can be extremely complicated. To avoid dealing with complicated sets \( M \), we restrict ourselves to the situation of integrals over compact intervals, i.e. to integrals over sets of the form \( M = [a, b] \). Moreover, we will also restrict ourselves to bounded functions \( f \), which we now define:

\begin{definition}
Let \( \emptyset \neq M \subseteq \mathbb{R} \) and \( f : M \to \mathbb{R} \). Then \( f \) is called bounded if, and only if, the set \( \{|f(x)| : x \in M\} \subseteq \mathbb{R}_0^+ \) is bounded, i.e. if, and only if,
\[
\|f\|_{\sup} := \sup\{|f(x)| : x \in M\} \in \mathbb{R}_0^+.
\] (10.3)
\end{definition}

The basic idea for the definition of the Riemann integral \( \int_M f \) is rather simple: Decompose the set \( M \) into small pieces \( I_1, \ldots, I_N \) and approximate \( \int_M f \) by the finite sum
\[
\sum_{j=1}^N f(x_j)|I_j|,
\]
where \( x_j \in I_j \) and \( |I_j| \) denotes the size of the set \( I_j \). Define \( \int_M f \) as the limit of such sums as the size of the \( I_j \) tends to zero (if the limit exists). However, to carry out this idea precisely and rigorously does require some work.

As stated before, we will assume that \( M \) is a closed bounded interval, and we will choose the \( I_j \) to be closed bounded intervals as well. To emphasize we are dealing with intervals, in the following, we will prefer to use the symbol \( I \) instead of \( M \).
Definition 10.2. If \( a, b \in \mathbb{R}, a \leq b \), and \( I := [a, b] \), then we call
\[
|I| := b - a = |a - b|,
\]
the length or the (1-dimensional) size, volume, or measure of \( I \).

Definition 10.3. Given a real interval \( I := [a, b] \subseteq \mathbb{R}, a, b \in \mathbb{R}, a < b \), the \((N+1)\)-tuple
\[
\Delta := (x_0, \ldots, x_N) \in \mathbb{R}^{N+1},
\]
\( N \in \mathbb{N} \), is called a partition of \( I \) if, and only if, \( a = x_0 < x_1 < \cdots < x_N = b \). We call \( x_0, \ldots, x_N \) the nodes of \( \Delta \), and let \( \nu(\Delta) := \{x_0, \ldots, x_N\} \)
be the set of all nodes. A tagged partition of \( I \) is a partition together with an \( N \)-tuple \((t_1, \ldots, t_N) \in \mathbb{R}^N\) such that \( t_j \in [x_{j-1}, x_j] \) for each \( j \in \{1, \ldots, N\} \). Given a partition \( \Delta \)
(with or without tags) of \( I \) as above and letting \( I_j := [x_{j-1}, x_j] \), the number
\[
|\Delta| := \max \{|I_j| : j \in \{1, \ldots, N\}\},
\]
is called the mesh size of \( \Delta \). It is sometimes convenient, if we extend our definitions to trivial intervals, consisting of just one point: For \( a = b \), we have \( I = [a, a] = \{a\} \). We then define \( \Delta = x_0 = a \) to be a partition of \( I \), \( \nu(\Delta) = \{x_0\} \), and \( a \) is then the only tag
that makes \( \Delta \) into a tagged partition. We also set \( I_0 := I = \{a\} \), and the mesh size in this case is \( |\Delta| = 0 \).

Definition 10.4. Let \( \Delta \) be a partition of \( I = [a, b] \subseteq \mathbb{R}, a \leq b \), as in Def. 10.3. Given
a function \( f : I \longrightarrow \mathbb{R} \) that is bounded according to Def. 10.1, define
\[
m_j := m_j(f) := \inf \{f(x) : x \in I_j\}, \quad M_j := M_j(f) := \sup \{f(x) : x \in I_j\},
\]
and
\[
r(\Delta, f) := \sum_{j=1}^{N} m_j \ |I_j| = \sum_{j=1}^{N} m_j (x_j - x_{j-1}),
\]
\[
R(\Delta, f) := \sum_{j=1}^{N} M_j \ |I_j| = \sum_{j=1}^{N} M_j (x_j - x_{j-1}),
\]
where \( r(\Delta, f) \) is called the lower Riemann sum and \( R(\Delta, f) \) is called the upper Riemann sum
associated with \( \Delta \) and \( f \). If \( \Delta \) is tagged by \( \tau := (t_1, \ldots, t_N) \), then we also define
the intermediate Riemann sum
\[
\rho(\Delta, f) := \sum_{j=1}^{N} f(t_j) \ |I_j| = \sum_{j=1}^{N} f(t_j) (x_j - x_{j-1}).
\]
Note that, for \( a = b \), all the above sums are empty and we have \( r(\Delta, f) = R(\Delta, f) = \rho(\Delta, f) = 0 \).

Definition 10.5. Let \( I = [a, b] \subseteq \mathbb{R} \) be an interval, \( a \leq b \), and suppose \( f : I \longrightarrow \mathbb{R} \)
is bounded.
(a) Define

\[ J_*(f, I) := \sup \{ r(\Delta, f) : \Delta \text{ is a partition of } I \}, \quad (10.8a) \]

\[ J^*(f, I) := \inf \{ R(\Delta, f) : \Delta \text{ is a partition of } I \}. \quad (10.8b) \]

We call \( J_*(f, I) \) the \textit{lower Riemann integral} of \( f \) over \( I \) and \( J^*(f, I) \) the \textit{upper Riemann integral} of \( f \) over \( I \).

(b) The function \( f \) is called \textit{Riemann integrable} over \( I \) if, and only if, \( J_*(f, I) = J^*(f, I) \). If \( f \) is Riemann integrable over \( I \), then

\[ \int_a^b f(x) \, dx := \int_I f(x) \, dx := \int_a^b f := \int_I f = J_*(f, I) = J^*(f, I) \quad (10.9) \]

is called the \textit{Riemann integral} of \( f \) over \( I \). The set of all functions \( f : I \rightarrow \mathbb{R} \) that are Riemann integrable over \( I \) is denoted by \( \mathcal{R}(I, \mathbb{R}) \) or just by \( \mathcal{R}(I) \).

(c) The function \( g : I \rightarrow \mathbb{C} \) is called \textit{Riemann integrable} over \( I \) if, and only if, both \( \text{Re} \, g \) and \( \text{Im} \, g \) are Riemann integrable. The set of all Riemann integrable functions \( g : I \rightarrow \mathbb{C} \) is denoted by \( \mathcal{R}(I, \mathbb{C}) \). If \( g \in \mathcal{R}(I, \mathbb{C}) \), then

\[ \int_I g := \left( \int_I \text{Re} \, g, \int_I \text{Im} \, g \right) = \int_I \text{Re} \, g + i \int_I \text{Im} \, g \in \mathbb{C} \quad (10.10) \]

is called the Riemann integral of \( g \) over \( I \).

\textbf{Remark 10.6.} If \( I = [a, b] \subseteq \mathbb{R} \), \( \Delta \) is a partition of \( I \), and \( f : I \rightarrow \mathbb{R} \) is bounded, then (10.6) implies

\[ m_j(f) \overset{(4.9c)}{=} -M_j(-f) \quad \text{and} \quad m_j(-f) \overset{(4.9d)}{=} -M_j(f), \quad (10.11a) \]

(10.7) implies

\[ r(\Delta, f) = -R(\Delta, -f) \quad \text{and} \quad r(\Delta, -f) = -R(\Delta, f), \quad (10.11b) \]

and (10.8) implies

\[ J_*(f, I) = -J^*(-f, I) \quad \text{and} \quad J_*(-f, I) = -J^*(f, I). \quad (10.11c) \]

\textbf{Example 10.7.} (a) If \( I = [a, b] \subseteq \mathbb{R} \) as before and \( f : I \rightarrow \mathbb{R} \) is constant, i.e. \( f \equiv c \) with \( c \in \mathbb{R} \), then \( f \in \mathcal{R}(I) \) and

\[ \int_a^b f = c (b - a) = c |I| : \quad (10.12) \]

We have, for each partition \( \Delta \) of \( I \),

\[ r(\Delta, f) = \sum_{j=1}^N m_j |I_j| = c \sum_{j=1}^N |I_j| = c |I| = c (b - a) = \sum_{j=1}^N M_j |I_j| = R(\Delta, f), \]

proving \( J_*(f, I) = c (b - a) = J^*(f, I) \).
(b) An example of a function that is not Riemann integrable for \(a < b\) is given by the \textit{Dirichlet function}

\[
 f : [a, b] \rightarrow \mathbb{R}, \quad f(x) := \begin{cases} 
 0 & \text{for } x \text{ irrational,} \\
 1 & \text{for } x \text{ rational,}
\end{cases} \quad a < b.
\]

(10.14)

Since \(r(\Delta, f) = 0\) and \(R(\Delta, f) = \sum_{j=1}^{N} |I_j| = b - a\) for every partition \(\Delta\) of \(I\), one obtains \(J^*(f, I) = 0\) for every partition \(\Delta\) of \([a, b]\), showing that \(f \notin R(I)\).

**Definition 10.8.** (a) If \(\Delta\) is a partition of \([a, b] \subseteq \mathbb{R}\) as in Def. 10.3, then another partition \(\Delta'\) of \([a, b]\) is called a \textit{refinement} of \(\Delta\) if, and only if, \(\nu(\Delta) \subseteq \nu(\Delta')\), i.e. if, and only if, the nodes of \(\Delta'\) include all the nodes of \(\Delta\).

(b) If \(\Delta\) and \(\Delta'\) are partitions of \([a, b]\), then the \textit{superposition} of \(\Delta\) and \(\Delta'\), denoted \(\Delta + \Delta'\), is the unique partition of \([a, b]\) having \(\nu(\Delta) \cup \nu(\Delta')\) as its set of nodes. Note that the superposition of \(\Delta\) and \(\Delta'\) is always a common refinement of \(\Delta\) and \(\Delta'\).

**Lemma 10.9.** Let \(a, b \in \mathbb{R}, a < b, I := [a, b]\), and suppose \(f : I \rightarrow \mathbb{R}\) is bounded with \(M := \|f\|_{\sup} \in \mathbb{R}_0^+\). Let \(\Delta'\) be a partition of \(I\) and assume

\[
 \alpha := \#\left(\nu(\Delta') \setminus \{a, b\}\right) \geq 1
\]

(10.15)

is the number of interior nodes that occur in \(\Delta'\). Then, for each partition \(\Delta\) of \(I\), the following holds:

\[
 r(\Delta, f) \leq r(\Delta + \Delta', f) \leq r(\Delta, f) + 2 \alpha M |\Delta|,
\]

(10.16a)

\[
 R(\Delta, f) \geq R(\Delta + \Delta', f) \geq R(\Delta, f) - 2 \alpha M |\Delta|.
\]

(10.16b)

**Proof.** We carry out the proof of (10.16a) – the proof of (10.16b) can be conducted completely analogous. Consider the case \(\alpha = 1\) and let \(\xi\) be the single element of \(\nu(\Delta') \setminus \{a, b\}\). If \(\xi \in \nu(\Delta)\), then \(\Delta + \Delta' = \Delta\), and (10.16a) is trivially true. If \(\xi \notin \nu(\Delta)\), then \(x_{k-1} < \xi < x_k\) for a suitable \(k \in \{1, \ldots, N\}\). Define

\[
 I' := [x_{k-1}, \xi], \quad I'' := [\xi, x_k]
\]

(10.17)

and

\[
 m' := \inf\{f(x) : x \in I'\}, \quad m'' := \inf\{f(x) : x \in I''\}.
\]

(10.18)

Then we obtain

\[
 r(\Delta + \Delta', f) - r(\Delta, f) = m' |I'| + m'' |I''| - m_k |I_k| = (m' - m_k) |I'| + (m'' - m_k) |I''|.
\]

(10.19)

Together with the observation

\[
 0 \leq m' - m_k \leq 2M, \quad 0 \leq m'' - m_k \leq 2M,
\]

(10.20)

(10.19) implies

\[
 0 \leq r(\Delta + \Delta', f) - r(\Delta, f) \leq 2M (|I'| + |I''|) \leq 2M |\Delta|.
\]

(10.21)

The general form of (10.16a) follows by an induction on \(\alpha\).
Theorem 10.10. Let \( a, b \in \mathbb{R}, \ a \leq b, \ I := [a, b], \) and let \( f : I \to \mathbb{R} \) be bounded.

(a) Suppose \( \Delta \) and \( \Delta' \) are partitions of \( I \) such that \( \Delta' \) is a refinement of \( \Delta \). Then
\[
r(\Delta, f) \leq r(\Delta', f), \quad R(\Delta, f) \geq R(\Delta', f).
\] (10.22)

(b) For arbitrary partitions \( \Delta \) and \( \Delta' \), the following holds:
\[
r(\Delta, f) \leq R(\Delta', f).
\] (10.23)

(c) \( J_*(f, I) \leq J^*(f, I) \).

(d) For each sequence of partitions \( (\Delta_n)_{n \in \mathbb{N}} \) of \( I \) such that \( \lim_{n \to \infty} |\Delta_n| = 0 \), one has
\[
\lim_{n \to \infty} r(\Delta_n, f) = J_*(f, I), \quad \lim_{n \to \infty} R(\Delta_n, f) = J^*(f, I).
\] (10.24)

In particular, if \( f \in \mathcal{R}(I) \), then
\[
\lim_{n \to \infty} r(\Delta_n, f) = \lim_{n \to \infty} R(\Delta_n, f) = \int_I f,
\] (10.25a)
and if \( f \in \mathcal{R}(I) \) and the \( \Delta_n \) are tagged, then also
\[
\lim_{n \to \infty} \rho(\Delta_n, f) = \int_I f.
\] (10.25b)

Proof. (a): If \( \Delta' \) is a refinement of \( \Delta \), then \( \Delta' = \Delta + \Delta' \). Thus, (10.22) is immediate from (10.16).

(b): This also follows from (10.16):
\[
r(\Delta, f) \overset{(10.16a)}{\leq} r(\Delta + \Delta', f) \overset{(10.7)}{\leq} R(\Delta + \Delta', f) \overset{(10.16b)}{\leq} R(\Delta', f).
\] (10.26)

(c): One just combines (10.8) with (b).

(d): For \( a = b \), there is nothing to show. For \( a < b \), let \( (\Delta_n)_{n \in \mathbb{N}} \) be a sequence of partitions of \( I \) such that \( \lim_{n \to \infty} |\Delta_n| = 0 \), and let \( \Delta' \) be an arbitrary partition of \( I \) with numbers \( \alpha \) and \( M \) defined as in Lem. 10.9. Then, according to (10.16a):
\[
r(\Delta_n, f) \leq r(\Delta_n + \Delta', f) \leq r(\Delta_n, f) + 2\alpha M |\Delta_n| \quad \text{for each } n \in \mathbb{N}.
\] (10.27)

From (b), we conclude the sequence \( (r(\Delta_n, f))_{n \in \mathbb{N}} \) is bounded. According to the Bolzano-Weierstrass Th. 7.27, if we can show that the sequence has \( J_*(f, I) \) as its only cluster point, then the first equality of (10.24) must hold. Thus, according to Prop. 7.26, it suffices to show that every converging subsequence of \( (r(\Delta_n, f))_{n \in \mathbb{N}} \) converges to \( J_*(f, I) \).

To this end, suppose \( (r(\Delta_n, f))_{k \in \mathbb{N}} \) is a converging subsequence of \( (r(\Delta_n, f))_{n \in \mathbb{N}} \) with \( \beta := \lim_{k \to \infty} r(\Delta_n, f) \). First note \( \beta \leq J_*(f, I) \) due to the definition of \( J_*(f, I) \). Moreover, (10.27) implies \( \lim_{k \to \infty} r(\Delta_n + \Delta', f) = \beta \). Since \( r(\Delta', f) \leq r(\Delta_n + \Delta', f) \) and \( \Delta' \) is arbitrary, we obtain \( J_*(f, I) \leq \beta \), i.e. \( J_*(f, I) = \beta \). Thus, we have shown that, indeed, every subsequence of \( (r(\Delta_n, f))_{n \in \mathbb{N}} \) converges to \( \beta = J_*(f, I) \). In the same manner, one conducts the proof of \( J^*(f, I) = \lim_{n \to \infty} R(\Delta_n, f) \). Then (10.25a) is immediate from the definition of Riemann integrability, and (10.25b) follows from (10.25a), since (10.7) implies \( r(\Delta, f) \leq \rho(\Delta, f) \leq R(\Delta, f) \) for each tagged partition \( \Delta \) of \( I \).
Theorem 10.11. Let $a, b \in \mathbb{R}$, $a \leq b$, $I := [a, b]$.

(a) The integral is linear: More precisely, if $f, g \in \mathcal{R}(I, \mathbb{K})$ and $\lambda, \mu \in \mathbb{K}$, then $\lambda f + \mu g \in \mathcal{R}(I, \mathbb{K})$ and

$$
\int_I (\lambda f + \mu g) = \lambda \int_I f + \mu \int_I g.
$$

(b) Let $\tilde{\Delta} = (y_0, \ldots, y_M)$, $M \in \mathbb{N}$, be a partition of $I$, $J_k := [y_{k-1}, y_k]$. Then $f \in \mathcal{R}(I, \mathbb{K})$ if, and only if, $f \in \mathcal{R}(J_k, \mathbb{K})$ for each $k \in \{1, \ldots, M\}$. If $f \in \mathcal{R}(I, \mathbb{K})$, then

$$
\int_a^b f = \int_I f = \sum_{k=1}^M \int_{J_k} f = \sum_{k=1}^M \int_{y_{k-1}}^{y_k} f.
$$

(c) Monotonicity of the Integral: If $f, g : I \rightarrow \mathbb{R}$ are bounded and $f \leq g$ (i.e. $f(x) \leq g(x)$ for each $x \in I$), then $J_*(f, I) \leq J_*(g, I)$ and $J^*(f, I) \leq J^*(g, I)$. In particular, if $f, g \in \mathcal{R}(I)$ and $f \leq g$, then

$$
\int_I f \leq \int_I g.
$$

(d) Triangle Inequality: For each $f \in \mathcal{R}(I, \mathbb{C})$, one has

$$
\left| \int_I f \right| \leq \int_I |f|.
$$

(e) Mean Value Theorem for Integration: If $f \in \mathcal{R}(I)$, then, for each $m, M \in \mathbb{R}$ with $m \leq f \leq M$:

$$
m (b - a) = m |I| \leq \int_a^b f = \int_I f \leq M |I| = M (b - a).
$$

The theorem’s name comes from the fact that, for $a < b$, $|I|^{-1} \int_I f$ is sometimes referred to as the mean value of $f$ on $I$.

Proof. (a): First, consider $\mathbb{K} = \mathbb{R}$, i.e. $f, g : I \rightarrow \mathbb{R}$ and $\lambda, \mu \in \mathbb{R}$. For $a = b$, there is nothing to prove, so let $a < b$. Let $(\Delta_n)_{n \in \mathbb{N}}$ be a sequence of partitions of $I$, $\Delta_n = (x_{n,0}, \ldots, x_{n,N_n})$, $I_{n,j} := [x_{n,j-1}, x_{n,j}]$, satisfying $\lim_{n \to \infty} |\Delta_n| = 0$. Note that, for
each \( n \in \mathbb{N} \) and each \( j \in \{1, \ldots, N_n\} \),

\[
m_{n,j}(f + g) = \inf \{f(x) + g(x) : x \in I_{n,j}\} \\
\geq \inf \{f(x) : x \in I_{n,j}\} + \inf \{g(x) : x \in I_{n,j}\} \\
= m_{n,j}(f) + m_{n,j}(g), \quad \text{(10.33a)}
\]

\[
M_{n,j}(f + g) = \sup \{f(x) + g(x) : x \in I_{n,j}\} \\
\leq \sup \{f(x) : x \in I_{n,j}\} + \sup \{g(x) : x \in I_{n,j}\} \\
= M_{n,j}(f) + M_{n,j}(g), \quad \text{(10.33b)}
\]

\[
\forall \lambda \in \mathbb{R} \quad m_{n,j}(\lambda f) = \inf \{\lambda f(x) : x \in I_{n,j}\} \quad \text{(4.9d)}
\]

\[
\forall \lambda \in \mathbb{R} \quad M_{n,j}(\lambda f) = \sup \{\lambda f(x) : x \in I_{n,j}\} \quad \text{(4.9c)}
\]

Thus,

\[
J_*(f + g, I) = \lim_{n \to \infty} r(\Delta_n, f + g) = \lim_{n \to \infty} \sum_{j=1}^{N_n} m_{n,j}(f + g) |I_{n,j}| \quad \text{(10.24)}
\]

\[
J^*(f + g, I) = \lim_{n \to \infty} R(\Delta_n, f + g) = \lim_{n \to \infty} \sum_{j=1}^{N_n} M_{n,j}(f + g) |I_{n,j}| \quad \text{(10.34b)}
\]

Thus, if \( f \) and \( g \) are both Riemann integrable over \( I \), then we obtain \( J_*(f + g, I) \geq J_*(f, I) + J_*(g, I) = J^*(f, I) + J^*(g, I) \geq J^*(f + g, I) \), i.e., by Th. 10.10(c), \((f + g) \in \mathcal{R}(I)\);
and $J_\lambda(\lambda f, I) = \lambda J_\lambda(f, I) = \lambda J^*(f, I)$ for $\lambda \geq 0$, $J_\lambda(\lambda f, I) = \lambda J^*(f, I) = \lambda J_\lambda(f, I)$ for $\lambda < 0$, i.e. $(\lambda f) \in \mathcal{R}(I)$ in each case. In particular, for each $\lambda, \mu \in \mathbb{R}$,

$$\int_I (\lambda f + \mu g) = J_\lambda(\lambda f + \mu g, I) = \lambda J_\lambda(f, I) + \mu J_\lambda(g, I) = \lambda \int_I f + \mu \int_I g,$$

proving (10.28) for $\mathbb{K} = \mathbb{R}$. It remains to consider $f, g \in \mathcal{R}(I, \mathbb{C})$ and $\lambda, \mu \in \mathbb{C}$. One computes, using the real-valued case,

$$\int_I (\lambda f) = \left( \int_I (\text{Re} \lambda \text{Re} f - \text{Im} \lambda \text{Im} f) \right) = \left( \text{Re} \lambda \int_I f - \text{Im} \lambda \int_I f \right) = \lambda \int_I f$$

and

$$\int_I (f + g) = \left( \int_I f + \int_I g \right) = \left( \int_I f + \int_I g \right) = \int_I f + \int_I g.$$

(b): Once again, consider $\mathbb{K} = \mathbb{R}$ first. For $a = b$, there is nothing to prove, so let $a < b$. For $M = 1$, there is still nothing to prove. For $M = 2$, we have $a = y_0 < y_1 < y_2 = b$. Consider a sequence $(\Delta_n)_{n \in \mathbb{N}}$ of partitions of $I$, $\Delta_n = (x_{n,0}, \ldots, x_{n,N_n})$, such that $\lim_{n \to \infty} \Delta_n = 0$ and $y_1 \in \nu(\Delta_n)$ for each $n \in \mathbb{N}$. Define $\Delta' := (x_{n,0}, \ldots, y_1)$, $\Delta'' := (y_1, \ldots, x_{n,N_n})$. Then $\Delta'$ and $\Delta''$ are partitions of $J_1$ and $J_2$, respectively, and $\lim_{n \to \infty} \Delta'_n = \lim_{n \to \infty} \Delta''_n = 0$. Moreover,

$$\forall n \in \mathbb{N} \quad \left( r(\Delta_n, f) = r(\Delta'_n, f) + r(\Delta''_n, f), \quad R(\Delta_n, f) = R(\Delta'_n, f) + R(\Delta''_n, f) \right),$$

implying $J_\lambda(f, I) = J_\lambda(f, J_1) + J_\lambda(f, J_2)$ and $J^*(f, I) = J^*(f, J_1) + J^*(f, J_2)$. This proves $\int_I f = \int_{J_1} f + \int_{J_2} f$ provided $f \in \mathcal{R}(I) \cap \mathcal{R}(J_1) \cap \mathcal{R}(J_2)$. So it just remains to show the claimed equivalence between $f \in \mathcal{R}(I)$ and $f \in \mathcal{R}(J_1) \cap \mathcal{R}(J_2)$. If $f \in \mathcal{R}(J_1) \cap \mathcal{R}(J_2)$, then $J_\lambda(f, I) = J_\lambda(f, J_1) + J_\lambda(f, J_2) = J^*(f, J_1) + J^*(f, J_2) = J^*(f, I)$, showing $f \in \mathcal{R}(I)$. Conversely, $J_\lambda(f, I) = J^*(f, I)$ implies $J_\lambda(f, J_1) = J^*(f, J_1) + J^*(f, J_2) - J_\lambda(f, J_2) \geq J^*(f, J_1)$, showing $J_\lambda(f, J_1) = J^*(f, J_1)$ and $f \in \mathcal{R}(J_1)$; $f \in \mathcal{R}(J_2)$ follows completely analogous. The general case now follows by induction on $M$. If, $f \in \mathcal{R}(I, \mathbb{C})$, then one computes, using the real-valued case,

$$\int_I f = \left( \int_I f, \int_I g \right) = \left( \sum_{k=1}^M \int_{J_k} f, \sum_{k=1}^M \int_{J_k} g \right) = \sum_{k=1}^M \int_{J_k} f.$$
(d): We will see in Th. 10.17(c) below, that \( f \in \mathcal{R}(I, \mathbb{R}) \) implies \( |f| \in \mathcal{R}(I) \). Let \( \Delta \) be an arbitrary partition of \( I \), tagged by \((t_1, \ldots, t_N)\). Then, using the same notation as in Def. 10.3 and Def. 10.4,

\[
\left| (\rho(\Delta, \text{Re } f), \rho(\Delta, \text{Im } f)) \right| := \left( \sum_{j=1}^N \text{Re } f(t_j) |I_j|, \sum_{j=1}^N \text{Im } f(t_j) |I_j| \right)
\leq \sum_{j=1}^N \left( \text{Re } f(t_j), \text{Im } f(t_j) \right) |I_j|
= \sum_{j=1}^N |f(t_j)| |I_j| =: \rho(\Delta, |f|).
\tag{10.35}
\]

Since the intermediate Riemann sums in (10.35) converge to the respective integrals by (10.25b), one obtains

\[
\left| \int_I f \right| = \lim_{|\Delta| \to 0} \left| (\rho(\Delta, \text{Re } f), \rho(\Delta, \text{Im } f)) \right| \leq \lim_{|\Delta| \to 0} \rho(\Delta, |f|) = \int_I |f|,
\]
proving (10.31).

(e): We compute

\[
m |I| \overset{(10.12)}{=} \int_I m \overset{(c)}{=} \int_I f \overset{(c)}{=} \int_I M = M |I|, \tag{10.36}
\]

thereby establishing the case.

\[\Box\]

**Theorem 10.12** (Riemann’s Integrability Criterion). Let \( I = [a, b] \subseteq \mathbb{R} \) and suppose \( f : I \to \mathbb{R} \) is bounded. Then \( f \) is Riemann integrable over \( I \) if, and only if, for each \( \epsilon > 0 \), there exists a partition \( \Delta \) of \( I \) such that

\[
R(\Delta, f) - r(\Delta, f) < \epsilon. \tag{10.37}
\]

**Proof.** Suppose, for each \( \epsilon > 0 \), there exists a partition \( \Delta \) of \( I \) such that (10.37) is satisfied. Then

\[
J^*(f, I) - J_*(f, I) \leq R(\Delta, f) - r(\Delta, f) < \epsilon, \tag{10.38}
\]
showing \( J^*(f, I) \leq J_*(f, I) \). As the opposite inequality always holds, we have \( J^*(f, I) = J_*(f, I) \), i.e. \( f \in \mathcal{R}(I) \) as claimed. Conversely, if \( f \in \mathcal{R}(I) \) and \((\Delta_n)_{n \in \mathbb{N}}\) is a sequence of partitions of \( I \) with \( \lim_{n \to \infty} |\Delta_n| = 0 \), then (10.25a) implies that, for each \( \epsilon > 0 \), there is \( N \in \mathbb{N} \) such that \( R(\Delta_n, f) - r(\Delta_n, f) < \epsilon \) for each \( n > N \).

The previous theorem will allow us to prove that every continuous function on \([a, b]\) is Riemann integrable. However, we will also need to make use of the following result:

**Proposition 10.13.** Let \( I = [a, b] \subseteq \mathbb{R}, \ a \leq b, \ f : I \to \mathbb{R} \). If \( f \) is continuous, then \( f \) is even uniformly continuous, i.e.

\[
\forall \epsilon \in \mathbb{R}^+ \exists \delta \in \mathbb{R}^+ \forall x, y \in I \ (|x - y| < \delta \implies |f(x) - f(y)| < \epsilon). \tag{10.39}
\]
Proof. Arguing by contraposition, we assume \( f \) not to be uniformly continuous on \( I \). Then the negation of (10.39) must hold, i.e.

\[
\exists \epsilon_0 \in \mathbb{R}^+ \quad \forall \epsilon \in \mathbb{R}^+ \quad \exists x, y \in I \quad (|x - y| < \delta \wedge |f(x) - f(y)| \geq \epsilon_0).
\] (10.40)

In particular, for each \( n \in \mathbb{N} \), there exist \( x_n, y_n \in I \) such that

\[
|x_n - y_n| < \delta := 1/n
\] (10.41)

and \(|f(x_n) - f(y_n)| \geq \epsilon_0\). Then the sequence \((x_n)_{n \in \mathbb{N}}\) is bounded and the Bolzano-Weierstrass Th. 7.27 provides a convergent subsequence \((x_{n_k})_{k \in \mathbb{N}}\), i.e. there is \( \xi \in \mathbb{R} \) with \( \lim_{n \to \infty} x_{n_k} = \xi \). Clearly, \( \xi \in [a, b] \) and (10.41) implies \( \lim_{n \to \infty} y_{n_k} = \xi \) as well. However, due to \(|f(x_{n_k}) - f(y_{n_k})| \geq \epsilon_0 > 0\), the sequences \((f(x_{n_k}))_{n \in \mathbb{N}}\) and \((f(y_{n_k}))_{n \in \mathbb{N}}\) can not both converge to \( f(\xi) \), showing that \( f \) can not be continuous.

Caveat 10.14. It is important in Prop. 10.13 that \( f \) is defined on a compact interval \( I \). The examples \( f : ]0, 1[ \to \mathbb{R}, f(x) := 1/x \), and \( f : \mathbb{R} \to \mathbb{R}, f(x) := x^2 \) are examples of continuous functions that are not uniformly continuous.

**Theorem 10.15.** Let \( I = [a, b] \subseteq \mathbb{R}, a \leq b \).

(a) If \( f : I \to \mathbb{C} \) is continuous, then \( f \) is Riemann integrable over \( I \).

(b) If \( f : I \to \mathbb{R} \) is increasing or decreasing, then \( f \) is Riemann integrable over \( I \).

**Proof.** (a): As \( f \) is continuous if, and only if, \( \text{Re} f \) and \( \text{Im} f \) are both continuous, it suffices to consider a real-valued continuous \( f \). For \( a = b \), there is nothing to prove, so let \( a < b \). First note that, if \( f \) is continuous on \( I = [a, b] \), then \( f \) is bounded by Th. 7.54. Moreover, \( f \) is uniformly continuous due to Prop. 10.13. Thus, given \( \epsilon > 0 \), there is \( \delta > 0 \) such that \(|x - y| < \delta \) implies \(|f(x) - f(y)| < \epsilon/|I| \) for each \( x, y \in I \). Then, for each partition \( \Delta \) of \( I \) satisfying \(|\Delta| < \delta\), we obtain

\[
R(\Delta, f) - r(\Delta, f) = \sum_{j=1}^{N} (M_j - m_j)|I_j| \leq \frac{\epsilon}{|I|} \sum_{j=1}^{N} |I_j| = \epsilon,
\] (10.42)

as \(|\Delta| < \delta \) implies \(|x - y| < \delta \) for each \( x, y \in I_j \) and each \( j \in \{1, \ldots, N\} \). Finally, (10.42) implies \( f \in \mathcal{R}(I) \) due to Riemann’s integrability criterion of Th. 10.12.

(b): Suppose \( f : [a, b] \to \mathbb{R} \) is increasing. Then \( f \) is bounded, as \( f(a) \leq f(x) \leq f(b) \) for each \( x \in [a, b] \). If \( f(a) = f(b) \), then \( f \) is constant. Thus, assume \( f(a) < f(b) \). Moreover, if \( \Delta = (x_0, \ldots, x_N) \) is a partition of \( I \) as in Def. 10.3, then

\[
R(\Delta, f) - r(\Delta, f) = \sum_{j=1}^{N} (M_j - m_j)|I_j| = \sum_{j=1}^{N} (f(x_j) - f(x_{j-1}))[I_j] \leq |\Delta|(f(b) - f(a)).
\] (10.43)

Thus, given \( \epsilon > 0 \), we have \( R(\Delta, f) - r(\Delta, f) < \epsilon \) for each partition \( \Delta \) of \( I \) satisfying \(|\Delta| < \epsilon/(f(b) - f(a))\). In consequence, \( f \in \mathcal{R}(I) \), once again due to Riemann’s integrability criterion of Th. 10.12. If \( f \) is decreasing, then \( -f \) is increasing, and Th. 10.11(a) establishes the case.
Definition and Remark 10.16. Let $M \subseteq \mathbb{C}$. A function $f : M \rightarrow \mathbb{C}$ is called Lipschitz continuous in $M$ with Lipschitz constant $L$ if, and only if,

$$
\exists L \in \mathbb{R}^+_0 \forall x,y \in M \quad |f(x) - f(y)| \leq L|x - y|. \quad (10.44)
$$

Every Lipschitz continuous function is, indeed, continuous, since, if $\xi \in M$ and $(y_n)_{n \in \mathbb{N}}$ is a sequence in $M$ with $\lim_{n \to \infty} y_n = \xi$, then (10.44) implies

$$
\forall n \in \mathbb{N} \quad |f(\xi) - f(y_n)| \leq L|\xi - y_n|, \quad (10.45)
$$

proving $\lim_{n \to \infty} f(y_n) = f(\xi)$. Moreover, it is not too much harder to prove Lipschitz continuous functions are even uniformly continuous, but we will not pursue this right now. On the other hand, $f : \mathbb{R}^+_0 \rightarrow \mathbb{R}, f(x) := \sqrt{x}$, is an example of a continuous function (actually, even uniformly continuous) that is not Lipschitz continuous.

Theorem 10.17. Let $a, b \in \mathbb{R}$, $a \leq b$, $I := [a, b]$.

(a) If $f \in \mathcal{R}(I, \mathbb{R})$ and $\phi : f(I) \rightarrow \mathbb{C}$ is Lipschitz continuous, then $\phi \circ f \in \mathcal{R}(I, \mathbb{C})$.

(b) If $f \in \mathcal{R}(I, \mathbb{C})$ and $\phi : f(I) \rightarrow \mathbb{R}$ is Lipschitz continuous, then $\phi \circ f \in \mathcal{R}(I, \mathbb{R})$.

(c) If $f \in \mathcal{R}(I)$, then $|f|, f^2, f^+, f^- \in \mathcal{R}(I)$. In particular, we, indeed, have (10.2) from the introduction (with $M$ replaced by $I$). If, in addition, there exists $\delta > 0$ such that $f(x) \geq \delta$ for each $x \in I$, then $1/f \in \mathcal{R}(I)$. Moreover, $|f| \in \mathcal{R}(I)$ also holds for $f \in \mathcal{R}(I, \mathbb{C})$.

(d) If $f, g \in \mathcal{R}(I)$, then $\max(f, g), \min(f, g) \in \mathcal{R}(I)$. If $f, g \in \mathcal{R}(I, \mathbb{R})$, then $f, g \in \mathcal{R}(I, \mathbb{R})$. If, in addition, there exists $\delta > 0$ such that $|g(x)| \geq \delta$ for each $x \in I$, then $f/g \in \mathcal{R}(I, \mathbb{R})$.

Proof. (a),(b): First, we carry out the proof for the case $f \in \mathcal{R}(I)$, $\phi : f(I) \rightarrow \mathbb{R}$. Thus, let $f \in \mathcal{R}(I)$ and let $\phi : f(I) \rightarrow \mathbb{R}$ be Lipschitz continuous. Then there exists $L \geq 0$ such that

$$
|\phi(x) - \phi(y)| \leq L|x - y| \quad \text{for each } x, y \in f(I). \quad (10.46)
$$

As $f \in \mathcal{R}(I)$, given $\epsilon > 0$, Th. 10.12 provides a partition $\Delta$ of $I$ such that $R(\Delta, f) - r(\Delta, f) < \epsilon/L$, and we obtain

$$
R(\Delta, \phi \circ f) - r(\Delta, \phi \circ f) = \sum_{j=1}^{N} (M_j(\phi \circ f) - m_j(\phi \circ f))|I_j|
$$

$$
\leq \sum_{j=1}^{N} L(M_j(f) - m_j(f))|I_j|
$$

$$
= L \left( R(\Delta, f) - r(\Delta, f) \right) < \epsilon. \quad (10.47)
$$
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Thus, \( \phi \circ f \in \mathcal{R}(I) \) by another application of Th. 10.12. We now prove (a) as stated, i.e. with \( \phi \) \( \mathbb{C} \)-valued: Assume \( \phi \) to be \( L \)-Lipschitz, \( L \geq 0 \). For each \( x, y \in f(I) \), one has

\[
|\text{Re} \phi(x) - \text{Re} \phi(y)| \leq \frac{L |x - y|}{\epsilon} \nonumber
\]

and

\[
|\text{Im} \phi(x) - \text{Im} \phi(y)| \leq \frac{L |x - y|}{\epsilon} \nonumber
\]

showing \( \text{Re} \phi \) and \( \text{Im} \phi \) are \( L \)-Lipschitz, such that \( \text{Re}(\phi \circ f) \) and \( \text{Im}(\phi \circ f) \) are Riemann integrable by the real-valued case treated above. It remains to show (b) as stated, i.e. with \( f \) \( \mathbb{C} \)-valued and \( \phi \) \( \mathbb{R} \)-valued: Assume \( \phi : f(I) \rightarrow \mathbb{R} \) to be \( L \)-Lipschitz, \( L \geq 0 \). If \( f \in \mathcal{R}(I, \mathbb{C}) \), then \( \text{Re} f, \text{Im} f \in \mathcal{R}(I, \mathbb{R}) \), and, given \( \epsilon > 0 \), Riemann’s integrability criterion of Th. 10.12 provides partitions \( \Delta_1, \Delta_2 \) of \( I \) such that

\[
R(\Delta, \text{Re} f) - r(\Delta, \text{Re} f) < \epsilon/2L, \quad R(\Delta, \text{Im} f) - r(\Delta, \text{Im} f) < \epsilon/2L.
\]

Recalling that, for each \( g : I \rightarrow \mathbb{R} \) and \( \Delta = (x_0, \ldots, x_N) \in \mathbb{R}^{N+1}, N \in \mathbb{N}, a = x_0 < x_1 < \cdots < x_N = b, I_j := [x_{j-1}, x_j] \), it is

\[
r(\Delta, g) = \sum_{j=1}^{N} m_j |I_j| = \sum_{j=1}^{N} m_j(g)(x_j - x_{j-1}),
\]

\[
R(\Delta, g) = \sum_{j=1}^{N} M_j |I_j| = \sum_{j=1}^{N} M_j(g)(x_j - x_{j-1}),
\]

where

\[
m_j(g) := \inf\{g(x) : x \in I_j\}, \quad M_j(g) := \sup\{g(x) : x \in I_j\},
\]

we obtain, for each \( \xi_j, \eta_j \in I_j \),

\[
|\phi \circ f(\xi_j) - \phi \circ f(\eta_j)| \leq L |f(\xi_j) - f(\eta_j)| \leq L |\text{Re} f(\xi_j) - \text{Re} f(\eta_j)| + L |\text{Im} f(\xi_j) - \text{Im} f(\eta_j)|
\]

and, thus,

\[
R(\Delta, \phi \circ f) - r(\Delta, \phi \circ f) = \sum_{j=1}^{N} (M_j(\phi \circ f) - m_j(\phi \circ f)) |I_j|
\]

\[
\leq \sum_{j=1}^{N} L (M_j(\text{Re} f) - m_j(\text{Re} f)) |I_j| + \sum_{j=1}^{N} L (M_j(\text{Im} f) - m_j(\text{Im} f)) |I_j|
\]

\[
= L (R(\Delta, \text{Re} f) - r(\Delta, \text{Re} f)) + L (R(\Delta, \text{Im} f) - r(\Delta, \text{Im} f)) \tag{10.49} \leq \epsilon. \tag{10.52}
\]
Thus, $\phi \circ f \in \mathcal{R}(I, \mathbb{R})$ by Th. 10.12.

(c): $|f|, f^2, f^+, f^- \in \mathcal{R}(I)$ follows from (b) (for $|f|$ also for $f \in \mathcal{R}(I, \mathbb{C})$), since each of the maps $x \mapsto |x|, x \mapsto x^2, x \mapsto \max\{x, 0\}, x \mapsto -\min\{x, 0\}$ is Lipschitz continuous on the bounded set $f(I)$ (recall that $f \in \mathcal{R}(I)$ implies that $f$ is bounded). Since $f = f^+ - f^-$, (10.2) is implied by (10.28). Finally, if $f(x) \geq \delta > 0$, then $x \mapsto x^{-1}$ is Lipschitz continuous on the bounded set $f(I)$, and $f^{-1} \in \mathcal{R}(I)$ follows from (b).

(d): For $f, g \in \mathcal{R}(I)$, we note that, due to

$$fg = \frac{1}{4} (f + g)^2 - (f - g)^2,$$

(10.53a)

$$\max(f, g) = f + (g - f)^+, \quad \min(f, g) = g - (f - g)^-,$$

(10.53b) (10.53c)

everything is a consequence of (c). For $f, g \in \mathcal{R}(I, \mathbb{C})$, due to

$$\bar{f} = (\Re f, -\Im f),$$

(10.53d)

$$fg = (\Re f \Re g - \Im f \Im g, \Re f \Im g + \Im f \Re g),$$

(10.53e)

$$1/g = (\Re g/|g|^2, -\Im g/|g|^2),$$

(10.53f)

everything follows from the real-valued case together with (c) and Th. 10.11(a), where $|g| \geq \delta > 0$ guarantees $|g|^2 \geq \delta^2 > 0$.

10.2 Important Theorems

This section compiles a number of important theorems on Riemann integrals, which, in particular, provide powerful tools to actually evaluate such integrals.

10.2.1 Fundamental Theorem of Calculus

We provide two variants of the fundamental theorem with slightly different flavors: In the first variant, Th. 10.19(a), we start with a function $f$, obtain another function $F$ by means of integrating $f$, and recover $f$ by taking the derivative of $F$. In the second variant, Th. 10.19(b), one first differentiates the given function $F$, obtaining $f := F'$, followed by integrating $f$, recovering $F$ up to an additive constant.

Notation 10.18. If $a, b \in \mathbb{R}, a \leq b$, $I := [a, b], f : I \rightarrow \mathbb{C}$, then denote

$$\int_a^b f := \int_I f,$$

$$\int_a^b f := -\int_b^a f,$$

(10.54a)

$$[f(t)]_a^b := [f]_a^b := f(b) - f(a), \quad [f(t)]_a^b := [f]_a^b := f(a) - f(b),$$

(10.54b)

where $f \in \mathcal{R}(I, \mathbb{C})$ for (10.54a).

Theorem 10.19. Let $a, b \in \mathbb{R}, a < b, I := [a, b]$. 

(a) If \( f \in \mathcal{R}(I, \mathbb{K}) \) is continuous in \( \xi \in I \), then, for each \( c \in I \), the function
\[
F_c : I \rightarrow \mathbb{K}, \quad F_c(x) := \int_c^x f(t) \, dt ,
\]
is differentiable in \( \xi \) with \( F'_c(\xi) = f(\xi) \). In particular, if \( f \in C(I, \mathbb{K}) \), then \( F_c \in C^1(I, \mathbb{K}) \) and \( F'_c(x) = f(x) \) for each \( x \in I \).

(b) If \( F \in C^1(I, \mathbb{K}) \) or, alternatively, \( F : I \rightarrow \mathbb{K} \) is differentiable with integrable derivative \( F' \in \mathcal{R}(I, \mathbb{K}) \), then
\[
F(b) - F(a) = [F(t)]^b_a = \int_a^b F'(t) \, dt ,
\]
and
\[
F(x) = F(c) + \int_c^x F'(t) \, dt \quad \text{for each } c, x \in I .
\]

Proof. It suffices to prove the case \( \mathbb{K} = \mathbb{R} \), since the case \( \mathbb{K} = \mathbb{C} \) then follows by applying the case \( \mathbb{K} = \mathbb{R} \) to \( \text{Re} \, F_c \) and \( \text{Im} \, F_c \) (for (a)) and to \( \text{Re} \, F \) and \( \text{Im} \, F \) (for (b)). Thus, for the rest of the proof, we assume \( \mathbb{K} = \mathbb{R} \).

(a): We need to show that
\[
\lim_{h \to 0} A(h) = 0 , \quad \text{where} \quad A(h) := \frac{F_c(\xi + h) - F_c(\xi)}{h} - f(\xi) .
\]
One computes
\[
A(h) = \frac{1}{h} \int_\xi^{\xi+h} f(t) \, dt - \frac{1}{h} f(\xi) \int_\xi^{\xi+h} \, dt = \frac{1}{h} \int_\xi^{\xi+h} (f(t) - f(\xi)) \, dt .
\]
Now, given \( \epsilon > 0 \), the continuity of \( f \) in \( \xi \) allows us to find \( \delta > 0 \) such that \( |(f(t) - f(\xi))| < \epsilon/2 \) for each \( t \in I \) with \( |t - \xi| < \delta \). Thus, for each \( h \) with \( |h| < \delta \), we obtain
\[
|A(h)| \leq \frac{1}{h} \int_\xi^{\xi+h} |f(t) - f(\xi)| \, dt \leq \frac{\epsilon h}{2h} < \epsilon ,
\]
thereby proving \( \lim_{h \to 0} A(h) = 0 \), i.e. \( f(\xi) = F'_c(\xi) \).

(b): First assume \( F \in C^1(I) \). Then \( F' \) is continuous on \( I \), and we can apply (a) to the function
\[
G : I \rightarrow \mathbb{R}, \quad G(x) := \int_a^x F'(t) \, dt ,
\]
to obtain \( G' = F' \). Thus, for \( H := F - G \), we obtain \( H' \equiv 0 \), showing that \( H \) must be constant on \( I \), i.e. \( H(x) = H(a) = F(a) - G(a) = F(a) \) for each \( x \in I \). Evaluating at \( x = b \) yields
\[
F(a) = H(b) = F(b) - \int_a^b F'(t) \, dt ,
\]
thereby establishing the case.

Now we consider the remaining case of a differentiable $F$ with integrable derivative $F' \in \mathcal{R}(I)$. Consider a partition $\Delta = (x_0, \ldots, x_N)$ of $I$ as in Def. 10.3. Then, for each $j \in \{1, \ldots, N\}$, the mean value theorem provides $\xi_j \in ]x_{j-1}, x_j[\text{ such that } F(x_j) - F(x_{j-1}) = (x_j - x_{j-1}) F'(\xi_j)$. Thus,

$$F(b) - F(a) = \sum_{j=1}^{N} (F(x_j) - F(x_{j-1})) = \sum_{j=1}^{N} (x_j - x_{j-1}) F'(\xi_j) = \rho(\Delta, F'). \tag{10.62}$$

If we choose a sequence of partitions $\Delta$ of $I$ such that $|\Delta| \to 0$, then the integrability of $f$ implies that the right-hand side of (10.62) converges to $\int_{a}^{b} F'$, once again establishing the case. ■

Definition 10.20. If $I \subseteq \mathbb{R}$, $f : I \rightarrow \mathbb{K}$, and $F : I \rightarrow \mathbb{K}$ is a differentiable function with $F' = f$, then $F$ is called a primitive or antiderivative of $f$.

Example 10.21. Due to the fundamental theorem, if we know a function’s antiderivative, we can easily compute its integral over a given interval. Here are three simple examples:

$$\int_{0}^{1} (x^5 - 3x) \, dx = \left[ \frac{x^6}{6} - \frac{3x^2}{2} \right]_{0}^{1} = \frac{1}{6} - \frac{3}{2} = -\frac{4}{3}, \tag{10.63a}$$

$$\int_{1}^{e} \frac{1}{x} \, dx = [\ln x]_{1}^{e} = \ln e - \ln 1 = 1, \quad \tag{10.63b}$$

$$\int_{0}^{\pi} \sin x \, dx = [-\cos x]_{0}^{\pi} = 2. \tag{10.63c}$$

10.2.2 Integration by Parts Formula

Theorem 10.22. Let $a, b \in \mathbb{R}$, $a < b$, $I := [a, b]$. If $f, g \in C^1(I, \mathbb{C})$, then the following integration by parts formula holds:

$$\int_{a}^{b} f g' = [fg]_{a}^{b} - \int_{a}^{b} f' g. \tag{10.64}$$

Proof. If $f, g \in C^1(I, \mathbb{C})$, then, according to the product rule, $fg \in C^1(I, \mathbb{C})$ with $(fg)' = f'g + fg'$. Applying (10.56a), we obtain

$$[fg]_{a}^{b} = \int_{a}^{b} (fg)' = \int_{a}^{b} f' g + \int_{a}^{b} fg', \tag{10.65}$$

which is precisely (10.64). ■

Example 10.23. We compute the integral $\int_{0}^{2\pi} \sin^2 t \, dt$:

$$\int_{0}^{2\pi} \sin^2 t \, dt = [-\sin t \cos t]_{0}^{2\pi} + \int_{0}^{2\pi} \cos^2 t \, dt = \int_{0}^{2\pi} \cos^2 t \, dt. \tag{10.66}$$
Adding $\int_0^{2\pi} \sin^2 t \, dt$ on both sides of (10.66) and using $\sin^2 + \cos^2 \equiv 1$ yields
\begin{equation}
2 \int_0^{2\pi} \sin^2 t \, dt = \int_0^{2\pi} 1 \, dt = 2\pi,
\end{equation}
i.e. $\int_0^{2\pi} \sin^2 t \, dt = \pi$.

### 10.2.3 Change of Variables

**Theorem 10.24.** Let $I, J \subseteq \mathbb{R}$ be intervals, $\phi \in C^1(I)$ and $f \in C(J, \mathbb{C})$. If $\phi(I) \subseteq J$, then the following change of variables formula holds for each $a, b \in I$:
\begin{equation}
\int_{\phi(a)}^{\phi(b)} f = \int_a^b f(\phi(t)) \phi'(t) \, dt = \int_a^b (f \circ \phi) \phi'.
\end{equation}

**Proof.** Let
\begin{equation}
F : J \longrightarrow \mathbb{C}, \quad F(x) := \int_{\phi(a)}^{x} f(t) \, dt.
\end{equation}
According to Th. 10.19(a) and the chain rule of Th. 9.10, we obtain
\begin{equation}
(F \circ \phi)' : I \longrightarrow \mathbb{C}, \quad (F \circ \phi)'(x) = \phi'(x)f(\phi(x)).
\end{equation}
Thus, we can apply (10.56a), which yields
\begin{equation}
\int_{\phi(a)}^{\phi(b)} f = F(\phi(b)) - F(\phi(a)) = \int_a^b (f \circ \phi) \phi',
\end{equation}
proving (10.68).

**Example 10.25.** We compute the integral $\int_0^1 t^2 \sqrt{1 - t} \, dt$ using the change of variables $x := \phi(t) := 1 - t$, $\phi'(t) = -1$:
\begin{align*}
\int_0^1 t^2 \sqrt{1 - t} \, dt &= -\int_0^1 (1 - x)^2 \sqrt{x} \, dx = \int_0^1 (\sqrt{x} - 2x\sqrt{x} + x^2\sqrt{x}) \, dx \\
&= \left[ \frac{2x^{\frac{7}{2}}}{3} - \frac{4x^{\frac{5}{2}}}{5} + \frac{2x^2}{7} \right]_0^1 = \frac{16}{105}.
\end{align*}

### 10.3 Improper Integrals

For our definition of the Riemann integral in Def. 10.5, it was important that we considered bounded functions on compact intervals (where the boundedness of the intervals was more important than the closedness) – for unbounded functions and/or unbounded intervals, even Def. 10.4 of lower and upper Riemann sums no longer makes sense.

Still, for sufficiently benign functions, it is possible to extend the notion of a definite Riemann integral to both unbounded intervals and unbounded functions, and in such situations we will speak of improper integrals (cf. Def. 10.29 below).
Definition 10.26. Let \( \emptyset \neq I \subseteq \mathbb{R} \) be an interval. We call \( f : I \rightarrow \mathbb{R} \) to be locally Riemann integrable if, and only if, \( f \in \mathcal{R}(J) \) for each compact interval \( J \subseteq I \). Let \( \mathcal{R}_{\text{loc}}(I) \) denote the set of all locally Riemann integrable functions on \( I \).

Remark 10.27. In particular, locally Riemann integrable functions are bounded on every compact interval. Moreover, \( \mathcal{R}_{\text{loc}}(I) = \mathcal{R}(I) \) if, and only if, \( I \) is a compact interval. For example, for each \( a, b \in \mathbb{R} \) with \( a < b \), the function given by the assignment rule

\[
f(x) := \frac{1}{(x-a)(x-b)}
\]

is clearly locally Riemann integrable, but not bounded on each of the intervals \([ -\infty, a [ , ] a, b [, \text{ and } ] b, \infty [ .

Before we can define improper Riemann integrals, we define, in partial extension of Def. 8.17:

Definition 10.28. Let \( M \subseteq \mathbb{R} \). If \( M \) is unbounded from above (resp. below, then \( f : M \rightarrow \mathbb{K} \) is said to tend to \( \eta \in \mathbb{K} \) (or to have the limit \( \eta \in \mathbb{K} \)) for \( x \rightarrow \infty \) (resp., for \( x \rightarrow -\infty \)) (denoted by \( \lim_{x \rightarrow \pm \infty} f(x) = \eta \)) if, and only if, for each sequence \( (\xi_k)_{k \in \mathbb{N}} \) in \( M \) with \( \lim_{k \rightarrow \infty} \xi_k = \pm \infty \) (resp. with \( \lim_{k \rightarrow \infty} \xi_k = -\infty \)), the sequence \( (f(\xi_k))_{k \in \mathbb{N}} \) converges to \( \eta \in \mathbb{K} \), i.e.

\[
\lim_{x \rightarrow \pm \infty} f(x) = \eta \iff \forall (\xi_k)_{k \in \mathbb{N}} \text{ in } M \left( \lim_{k \rightarrow \infty} \xi_k = \pm \infty \Rightarrow \lim_{k \rightarrow \infty} f(\xi_k) = \eta \right). \tag{10.73}
\]

Definition 10.29. Let \( a < c < b \) (\( a = -\infty, b = \infty \) is admissible).

(a) Let \( I := [c,b[, f \in \mathcal{R}_{\text{loc}}(I) \), and assume \( b = \infty \) or \( f \) is unbounded. Consider the function

\[
F : I \rightarrow \mathbb{R}, \quad F(x) := \int_x^c f.
\]

If the limit

\[
\lim_{x \rightarrow b} F(x) = \lim_{x \rightarrow b} \int_x^c f
\]

exists in \( \mathbb{R} \), then we define

\[
\int_I f := \int_c^b f(t) \, dt := \int_c^b f := \lim_{x \rightarrow b} \int_x^c f.
\]

(b) Let \( I := [a,c] \), \( f \in \mathcal{R}_{\text{loc}}(I) \), and assume \( a = -\infty \) or \( f \) is unbounded. Consider the function

\[
F : I \rightarrow \mathbb{R}, \quad F(x) := \int_x^c f.
\]

If the limit

\[
\lim_{x \rightarrow a} F(x) = \lim_{x \rightarrow a} \int_x^c f
\]

exists in \( \mathbb{R} \), then we define

\[
\int_I f := \int_c^b f := \lim_{x \rightarrow a} \int_x^c f.
\]

(10.74)
exists in \( \mathbb{R} \), then we define
\[
\int_I f := \int_a^c f(t) \, dt := \int_a^c f := \lim_{x \to a} \int_x^c f.
\]

(c) Let \( I = ]a, b[ \), \( f \in \mathcal{R}_{\text{loc}}(I) \). If the conditions of both (a) and (b) hold, i.e. (i) – (iv), where
\begin{enumerate}
\item \( b = \infty \) or \( f \) is unbounded on \([c, b]\),
\item \( \lim_{x \to b} \int_c^x f \) exists in \( \mathbb{R} \),
\item \( a = -\infty \) or \( f \) is unbounded on \([a, c]\),
\item \( \lim_{x \to a} \int_x^c f \) exists in \( \mathbb{R} \),
\end{enumerate}
then we define
\[
\int_I f := \int_a^b f(t) \, dt := \int_a^b f := \int_a^c f + \int_c^b f.
\]

All the above limits of Riemann integrals (if they exist) are called *improper* Riemann integrals. In each case, if the limit exists, we call \( f \) improperly Riemann integrable and write \( f \in \mathcal{R}(I) \).

**Remark 10.30.** (a) The definitions in Def. 10.29 are consistent with what occurs if the limits are *proper* Riemann integrals: Let \( a, c, b \in \mathbb{R} \), \( a < c < b \), and \( f \in \mathcal{R}[a, b] \). Then
\[
\lim_{x \to b} \int_c^x f = \int_c^b f \quad \text{and} \quad \lim_{x \to a} \int_x^c f = \int_a^c f. \tag{10.76}
\]
Indeed, since \( f \in \mathcal{R}[a, b] \), \( |f| \) is bounded by some \( M \in \mathbb{R}^+ \); and if \((x_k)_{k \in \mathbb{N}}\) is a sequence in \([a, b]\) such that \( \lim_{k \to \infty} x_k = b \), then
\[
\left| \int_{x_k}^b f \right| \leq \int_{x_k}^b |f| \leq M (b - x_k) \to 0 \quad \text{for} \quad k \to \infty,
\]
implying
\[
\lim_{k \to \infty} \int_c^{x_k} f \quad \text{Th. 10.11(b)} \quad \lim_{k \to \infty} \left( \int_c^b f - \int_{x_k}^b f \right) = \int_c^b f - 0 = \int_c^b f.
\]
An analogous argument shows the remaining equality in (10.76).

(b) In Def. 10.29(c), it can occur that \( \int_{-\infty}^x f \) does not exist, even though the limit \( \lim_{x \to \infty} \int_{-x}^x f \) exists: For example, if \( f : \mathbb{R} \to \mathbb{R}, f(x) = x \), then \( f \in \mathcal{R}_{\text{loc}}(\mathbb{R}) \), and, for each sequence \((x_k)_{k \in \mathbb{N}}\) in \( \mathbb{R} \) such that \( \lim_{k \to \infty} x_k = \infty \) and each \( c \in \mathbb{R} \), one has
\begin{align*}
\lim_{k \to \infty} \int_{-x_k}^{x_k} t \, dt &= \lim_{k \to \infty} \left[ \frac{t^2}{2} \right]_{-x_k}^{x_k} = \lim_{k \to \infty} \frac{x_k^2 - x_k^2}{2} = 0, \\
\lim_{k \to \infty} \int_c^{x_k} t \, dt &= \lim_{k \to \infty} \left[ \frac{t^2}{2} \right]_c^{x_k} = \lim_{k \to \infty} \frac{x_k^2 - c^2}{2} = \infty, \\
\lim_{k \to \infty} \int_{-x_k}^c t \, dt &= \lim_{k \to \infty} \left[ \frac{t^2}{2} \right]_{-x_k}^c = \lim_{k \to \infty} \frac{c^2 - x_k^2}{2} = -\infty.
\end{align*}
i.e. \( \lim_{x \to \infty} \int_{-x}^x t \, dt = 0 \), but neither \( \lim_{x \to \infty} \int_c^x t \, dt \) nor \( \lim_{x \to -\infty} \int_{-x}^c t \, dt \) exists in \( \mathbb{R} \).

(c) Let \( a < c_1 < c_2 < b \) (\( a = -\infty, b = \infty \) is admissible). If \( I := [c_1, b[ \in \mathcal{R}_{\text{loc}}(I) \), and \( b = \infty \) or \( f \) is unbounded, then \( \int_{c_1}^b f \) exists if, and only if, \( \int_{c_2}^b f \) exists. Moreover, if the integrals exist, then

\[
\int_{c_1}^b f = \int_{c_1}^{c_2} f + \int_{c_2}^b f. \tag{10.77a}
\]

Indeed, if \( (x_k)_{k \to \infty} \) is a sequence in \( [c_1, b[ \) such that \( \lim_{k \to \infty} x_k = b \) and if \( \int_{c_1}^b f \) exists, then

\[
\lim_{k \to \infty} \int_{c_2}^{x_k} f = \lim_{k \to \infty} \left( \int_{c_1}^{x_k} f - \int_{c_1}^{c_2} f \right) = \int_{c_1}^b f - \int_{c_1}^{c_2} f,
\]

proving \( \int_{c_2}^b f \) exists and (10.77a) holds. Conversely, if \( \int_{c_2}^b f \) exists, then

\[
\lim_{k \to \infty} \int_{c_1}^{x_k} f = \lim_{k \to \infty} \left( \int_{c_2}^{x_k} f + \int_{c_1}^{c_2} f \right) = \int_{c_2}^b f + \int_{c_1}^{c_2} f,
\]

proving \( \int_{c_1}^b f \) exists and (10.77a) holds. Analogously, one shows that if \( I := ]a, c_2] \), \( f \in \mathcal{R}_{\text{loc}}(I) \), and \( a = -\infty \) or \( f \) is unbounded, then \( \int_{a}^{c_1} f \) exists if, and only if, \( \int_{a}^{c_2} f \) exists, where, if the integrals exist, then

\[
\int_{a}^{c_2} f = \int_{c_1}^{c_2} f + \int_{a}^{c_1} f. \tag{10.77b}
\]

In particular, we see that neither the existence nor the value of the improper integral in Def. 10.29(c) depends on the choice of \( c \).

**Example 10.31.** (a) Let \( 0 < \alpha < 1 \). We claim that

\[
\int_0^1 \frac{1}{t^\alpha} \, dt = \frac{1}{1 - \alpha} \tag{10.78}
\]

Indeed, if \( (x_k)_{k \in \mathbb{N}} \) is a sequence in \( [0, 1] \) such that \( \lim_{k \to \infty} x_k = 0 \), then

\[
\lim_{k \to \infty} \int_0^1 \frac{1}{t^\alpha} \, dt = \lim_{k \to \infty} \left[ \frac{t^{1-\alpha}}{1-\alpha} \right]_0^{1/x_k} = \lim_{k \to \infty} \frac{1 - x_k^{1-\alpha}}{1 - \alpha} = \frac{1}{1 - \alpha}.
\]

(b) If \( (x_k)_{k \in \mathbb{N}} \) is a sequence in \( [0, 1] \) such that \( \lim_{k \to \infty} x_k = 0 \), then

\[
\lim_{k \to \infty} \int_{x_k}^1 \frac{1}{t} \, dt = \lim_{k \to \infty} \left[ \ln t \right]_{x_k}^1 = \lim_{k \to \infty} \left( 0 - \ln x_k \right) = \infty,
\]

showing the limit does not exist in \( \mathbb{R} \), but diverges to \( \infty \). Sometimes, this is stated in the form

\[
\int_0^1 \frac{1}{t} \, dt = \infty. \tag{10.79}
\]
(c) We claim that
\[
\int_{-\infty}^{0} e^t \, dt = 1.
\] (10.80)
Indeed, if \((x_k)_{k \in \mathbb{N}}\) is a sequence in \(\mathbb{R}_0^-\) such that \(\lim_{k \to \infty} x_k = -\infty\), then
\[
\lim_{k \to \infty} \int_{x_k}^{0} e^t \, dt = \lim_{k \to \infty} \left[ e^t \right]_{x_k}^{0} = \lim_{k \to \infty} (1 - e^{x_k}) = 1.
\]

(d) Consider the function
\[
f : \mathbb{R}_0^+ \to \mathbb{R}, \quad f(t) := \begin{cases} n & \text{for } n \leq t \leq n + \frac{1}{n^2}, \, n \in \mathbb{N}, \\ 0 & \text{otherwise}. \end{cases}
\]
Then \(\lim_{t \to \infty} f(t)\) does not exist and \(f\) is not even bounded. However \(f \in \mathcal{R}(\mathbb{R}_0^+)\) and
\[
\int_{0}^{\infty} f(t) \, dt = \sum_{n=1}^{\infty} \int_{n}^{n+1/(n^2)} n \, dt = \sum_{n=1}^{\infty} 2^{-n} = \frac{1}{1 - \frac{1}{2}} - 1 = 1.
\]

**Lemma 10.32.** Let \(a < c < b\) (\(a = -\infty, \, b = \infty\) is admissible). Let \(I \subseteq [a,b]\) be one of the three kinds of intervals occurring in Def. 10.29 (i.e. \(I = [c,b], \, I = ]a,c], \) or \(I = ]a,b]\), and assume \(f,g : I \to \mathbb{R}\) to be improperly Riemann integrable over \(I\).

(a) Linearity: For each \(\lambda, \mu \in \mathbb{R}\), \(\lambda f + \mu g\) is improperly Riemann integrable over \(I\) and
\[
\int_{I} (\lambda f + \mu g) = \lambda \int_{I} f + \mu \int_{I} g.
\]
(b) Monotonicity: If \(f \leq g\), then
\[
\int_{I} f \leq \int_{I} g.
\]

**Proof.** We conduct the proof for the case \(I = [c,b]\) – the case \(I = ]a,c]\) can be shown analogously, and the case \(I = ]a,b]\) then also follows. Let \((x_k)_{k \in \mathbb{N}}\) be a sequence in \(I\) such that \(\lim_{k \to \infty} x_k = b\).

(a): One computes
\[
\lim_{k \to \infty} \int_{c}^{x_k} (\lambda f + \mu g) \stackrel{\text{Th. 10.11(a)}}{=} \lim_{k \to \infty} \left( \lambda \int_{c}^{x_k} f + \mu \int_{c}^{x_k} g \right) = \lambda \int_{c}^{b} f + \mu \int_{c}^{b} g,
\]
showing \((\lambda f + \mu g) \in \mathcal{R}(I)\) and proving (a).

(b): One estimates
\[
\int_{c}^{b} f = \lim_{k \to \infty} \int_{c}^{x_k} f \stackrel{\text{Th. 10.11(c)}}{\leq} \lim_{k \to \infty} \int_{c}^{x_k} g = \int_{c}^{b} g,
\]
proving (b).
Definition 10.33. Let \( a < c < b \) (\( a = -\infty \), \( b = \infty \) is admissible). Let \( I \subseteq \mathbb{R} \) be one of the three kinds of intervals occurring in Def. 10.29 (i.e. \( I = [c,b], I = ]a,c], \) or \( I = ]a,b[ \)), and assume \( f \in R_{\text{loc}}(I) \). Then, by Th. 10.17(c), \( |f| \in R_{\text{loc}}(I) \). If \( \int_I |f| \) exists as an improper integral, then we call the improper integral \( \int_I f \) absolutely convergent.

Before we can proceed to Prop. 10.35 about convergence criteria for improper integrals, we need to prove the analogon of Th. 7.19 for limits of functions.

Proposition 10.34. Let \( \emptyset \neq M \subseteq \mathbb{R} \), \( a \in \mathbb{R} \cup \{-\infty\} \), \( b \in \mathbb{R} \cup \{\infty\} \), and assume

\[
a = \begin{cases} 
\inf(M \setminus \{a\}) & \text{if } M \text{ is bounded from below,} \\
-\infty & \text{if } M \text{ is unbounded from below,}
\end{cases} \tag{10.81a}
\]

\[
b = \begin{cases} 
\sup(M \setminus \{a\}) & \text{if } M \text{ is bounded from above,} \\
\infty & \text{if } M \text{ is unbounded from above.}
\end{cases} \tag{10.81b}
\]

Let \( f : M \rightarrow \mathbb{R} \) be monotone (increasing or decreasing). Defining \( A := f(M) = \{f(x) : x \in M\} \), the following holds:

\[
\lim_{x \to b} f(x) = \begin{cases} 
\sup A & \text{if } f \text{ is increasing and } A \text{ is bounded from above,} \\
\infty & \text{if } f \text{ is increasing and } A \text{ is not bounded from above,} \\
\inf A & \text{if } f \text{ is decreasing and } A \text{ is bounded from below,} \\
-\infty & \text{if } f \text{ is decreasing and } A \text{ is not bounded from below.}
\end{cases} \tag{10.82a}
\]

\[
\lim_{x \to a} f(x) = \begin{cases} 
\sup A & \text{if } f \text{ is decreasing and } A \text{ is bounded from above,} \\
\infty & \text{if } f \text{ is decreasing and } A \text{ is not bounded from above,} \\
\inf A & \text{if } f \text{ is increasing and } A \text{ is bounded from below,} \\
-\infty & \text{if } f \text{ is increasing and } A \text{ is not bounded from below.}
\end{cases} \tag{10.82b}
\]

Proof. We prove (10.82a) for the case, where \( f \) is increasing – the remaining case of (10.82a) as well as (10.82b) can be proved completely analogous. Let \( (x_k)_{k \in \mathbb{N}} \) be a sequence in \( M \setminus \{b\} \) such that \( \lim_{k \to \infty} x_k = b \). We have to show that \( \lim_{k \to \infty} f(x_k) = \eta \), where \( \eta := \sup A \) for \( A \) bounded from above and \( \eta := \infty \) for \( A \) not bounded from above. Seeking a contradiction, assume \( \lim_{k \to \infty} f(x_k) = \eta \) does not hold. Due to the choice of \( b \), there then must be \( \epsilon > 0 \) and a subsequence \( (y_k)_{k \in \mathbb{N}} \) of \( (x_k)_{k \in \mathbb{N}} \) such that \( (y_k)_{k \in \mathbb{N}} \) is strictly increasing and

\[
\forall k \in \mathbb{N} \quad f(y_k) \leq \begin{cases} 
\eta - \epsilon & \text{if } \eta = \sup A, \\
\epsilon & \text{if } \eta = \infty.
\end{cases}
\]

Since \( \lim_{k \to \infty} y_k = b \) and \( f \) is increasing, this means \( \sup A \leq \eta - \epsilon \) or \( \sup A = \epsilon \), which means a contradiction in each case. Thus, \( \lim_{k \to \infty} f(x_k) = \eta \) must hold and the proof is complete. \( \blacksquare \)
Proposition 10.35. Let \( a < c < b \) (\( a = -\infty, b = \infty \) is admissible). Let \( I \subseteq ]a, b[ \) be one of the three kinds of intervals occurring in Def. 10.29 (i.e. \( I = [c, b[, I = ]a, c[, \) or \( I = ]a, b[ \)), and assume \( f \in \mathcal{R}_{\text{loc}}(I) \).

(a) If \( g \in \mathcal{R}_{\text{loc}}(I) \), \( 0 \leq f \leq g \) and \( \int_I g \) exists, then \( \int_I f \) exists as well. Conversely, if \( 0 \leq g \leq f \) and \( \int_I g \) diverges, then \( \int_I f \) diverges as well.

(b) If \( \int_I f \) is an improper integral that is absolutely convergent, then it is also convergent.

Proof. (a): We consider the case \( I = [c, b[ \) – the proof for the case \( I = ]a, c[ \) is completely analogous, and the case \( I = ]a, b[ \) then also follows. First, suppose \( 0 \leq f \leq g \), and \( \int_I g \) exists. Since \( 0 \leq f \), the function \( F: [c, b[ \rightarrow \mathbb{R}_0^+ \), \( F(x) := \int_c^x f \),

is increasing. Due to

\[
\forall x \in [c, b[ \quad F(x) = \int_c^x f \leq \int_c^x g \leq \int_c^b g \in \mathbb{R}_0^+, 
\]

\( F \) is also bounded from above (in the sense that \( \{F(x) : x \in [c, b[\} \) is bounded from above), i.e. Prop. 10.34 yields that \( \lim_{x \to b} F(x) = \lim_{x \to b} \int_c^x f \) exists in \( \mathbb{R} \) as claimed.

Now suppose \( 0 \leq g \leq f \) and \( \int_I g \) diverges. As the function \( F \) above, the function

\( G: [c, b[ \rightarrow \mathbb{R}_0^+ \), \( G(x) := \int_c^x g \),

is increasing. Since we assume that \( \lim_{x \to b} G(x) \) does not exist in \( \mathbb{R} \), Prop. 10.34 implies \( \lim_{x \to b} G(x) = \infty \). As a consequence, if \( (x_k)_{k \in \mathbb{N}} \) is a sequence in \( [c, b[ \) such that \( \lim_{k \to \infty} x_k = b \), then

\[
\lim_{k \to \infty} F(x_k) = \lim_{k \to \infty} \int_c^{x_k} f = \lim_{k \to \infty} \int_c^{x_k} g = \infty,
\]

showing that \( \int_I f \) diverges as well.

(b): We assume \( \int_I f \) to converge absolutely, i.e. \( \int_I |f| \) must exist in \( \mathbb{R} \). Since \( 0 \leq f^+ \leq |f| \) and \( 0 \leq f^- \leq |f| \), (a) then implies the existence of \( \int_I f^+ \) and of \( \int_I f^- \). Thus, according to Lem. 10.32(a), \( \int_I f = \int_I f^+ - \int_I f^- \) must also exist.

Example 10.36. (a) We will use Prop. 10.35(a) to show that the improper integral

\[
\int_{0}^{\infty} e^{-t^2} \, dt
\]

exists. Indeed,

\[
\forall t \in \mathbb{R} \quad (t - 1)^2 = t^2 - 2t + 1 \geq 0 \quad \Rightarrow \quad -t^2 \leq -2t + 1 \quad \Rightarrow \quad 0 \leq e^{-t^2} \leq e^{-2t+1},
\]
and, since
\[
\int_0^\infty e^{-2t+1} \, dt = \lim_{x \to \infty} \int_0^x e^{-2t+1} \, dt = \lim_{x \to \infty} \left[ -\frac{e^{-2t+1}}{2} \right]_0^x = \lim_{x \to \infty} \frac{e^{-2x+1}}{2} = \frac{e}{2}.
\]

Prop. 10.35(a) implies that \( f_0^\infty e^{-t^2} \, dt \) exists in \( \mathbb{R} \).

(b) We will use Prop. 10.35(a) to show that
\[
\int_0^\infty e^{t^2} \, dt
\]
diverges. Indeed,
\[
\forall t \in \mathbb{R}, \quad (t^2 \geq 0 \Rightarrow e^{t^2} \geq 1),
\]
and, since
\[
\lim_{x \to \infty} \int_0^x 1 \, dt = \lim_{x \to \infty} x = \infty,
\]
Prop. 10.35(a) implies that \( f_0^\infty e^{t^2} \, dt = \infty \).

(c) We provide an example that shows an improper integral can converge without converging absolutely: Consider the function
\[
f : [0, \infty[ \to \mathbb{R}, \quad f(t) := \begin{cases} (-1)^{n+1} & \text{for } n \leq t \leq n + \frac{1}{n}, \ n \in \mathbb{N}, \\ 0 & \text{otherwise}. \end{cases}
\]
Then
\[
\int_0^\infty |f| = \lim_{k \to \infty} \sum_{k=1}^n \int_k^{k+\frac{1}{n}} 1 \, dt = \lim_{k \to \infty} \sum_{k=1}^n \frac{1}{k} = \infty,
\]
showing \( f_0^\infty f \) does not converge absolutely. However, we will show
\[
\int_0^\infty f = \sum_{j=1}^\infty \left( \frac{(-1)^{j+1}}{j} \right) =: \alpha > 0.
\]
We know \( \alpha > 0 \) from Ex. 7.86(a) and Th. 7.85. Let \( (x_k)_{k \in \mathbb{N}} \) be a sequence in \( \mathbb{R}_0^+ \) such that \( \lim_{k \to \infty} x_k = \infty \). Given \( \varepsilon > 0 \), choose \( K \in \mathbb{N} \) such that \( \frac{1}{K} < \frac{\varepsilon}{2} \) and \( N \in \mathbb{N} \) such that
\[
\forall k > N, \quad x_k > K.
\]
Then, for each \( k > N \), there exists \( K_1 \in \mathbb{N} \) such that \( K < K_1 \leq x_k < K_1 + 1 \). Thus
\[
\int_0^{x_k} f(t) \, dt = \min \left\{ x_k - K_1, \frac{1}{K_1} \right\} + \sum_{j=1}^{K_1-1} \left( \frac{(-1)^{j+1}}{j} \right)
\]
and

$$\left| \alpha - \int_0^{x_k} f(t) \, dt \right| = \left| \sum_{j=K_1}^{\infty} \frac{(-1)^{j+1}}{j} - \min \left\{ x_k - K_1, \frac{1}{K_1} \right\} \right| \leq \frac{1}{K_1} + \frac{1}{K_1}$$

$$< \frac{2}{K} < 2 \cdot \frac{\epsilon}{2} = \epsilon,$$

thereby proving (10.85).

References


