## V. Functions of several variables

## V.1. $\mathbb{R}^{n}$ as a linear and metric space

Definition. The set $\mathbb{R}^{n}, n \in \mathbb{N}$, is the set of all ordered $n$-tuples of real numbers, i.e.

$$
\mathbb{R}^{n}=\left\{\left[x_{1}, \ldots, x_{n}\right]: x_{1}, \ldots, x_{n} \in \mathbb{R}\right\}
$$

For $\boldsymbol{x}=\left[x_{1}, \ldots, x_{n}\right] \in \mathbb{R}^{n}, \boldsymbol{y}=\left[y_{1}, \ldots, y_{n}\right] \in \mathbb{R}^{n}$ and $\alpha \in \mathbb{R}$ we set

$$
\boldsymbol{x}+\boldsymbol{y}=\left[x_{1}+y_{1}, \ldots, x_{n}+y_{n}\right], \quad \alpha \boldsymbol{x}=\left[\alpha x_{1}, \ldots, \alpha x_{n}\right] .
$$

Further, we denote $\boldsymbol{o}=[0, \ldots, 0]-$ the origin.
Definition. The Euclidean metric (distance) on $\mathbb{R}^{n}$ is the function $\rho: \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow[0,+\infty$ ) defined by

$$
\rho(\boldsymbol{x}, \boldsymbol{y})=\sqrt{\sum_{i=1}^{n}\left(x_{i}-y_{i}\right)^{2}}
$$

The number $\rho(\boldsymbol{x}, \boldsymbol{y})$ is called the distance of the point $\boldsymbol{x}$ from the point $\boldsymbol{y}$.
Theorem 1 (properties of the Euclidean metric). The Euclidean metric $\rho$ has the following properties:
(i) $\forall \boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}: \rho(\boldsymbol{x}, \boldsymbol{y})=0 \Leftrightarrow \boldsymbol{x}=\boldsymbol{y}$,
(ii) $\forall \boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}: \rho(\boldsymbol{x}, \boldsymbol{y})=\rho(\boldsymbol{y}, \boldsymbol{x})$,
(symmetry)
(iii) $\forall \boldsymbol{x}, \boldsymbol{y}, \boldsymbol{z} \in \mathbb{R}^{n}: \rho(\boldsymbol{x}, \boldsymbol{y}) \leq \rho(\boldsymbol{x}, \boldsymbol{z})+\rho(\boldsymbol{z}, \boldsymbol{y})$,
(triangle inequality)
(iv) $\forall \boldsymbol{x}, \boldsymbol{y} \in \mathbb{R}^{n}, \forall \lambda \in \mathbb{R}: \rho(\lambda \boldsymbol{x}, \lambda \boldsymbol{y})=|\lambda| \rho(\boldsymbol{x}, \boldsymbol{y})$,
(homogeneity)
(v) $\forall \boldsymbol{x}, \boldsymbol{y}, \boldsymbol{z} \in \mathbb{R}^{n}: \rho(\boldsymbol{x}+\boldsymbol{z}, \boldsymbol{y}+\boldsymbol{z})=\rho(\boldsymbol{x}, \boldsymbol{y})$.

Definition. Let $\boldsymbol{x} \in \mathbb{R}^{n}, r \in \mathbb{R}, r>0$. The set $B(\boldsymbol{x}, r)$ defined by

$$
B(\boldsymbol{x}, r)=\left\{\boldsymbol{y} \in \mathbb{R}^{n} ; \rho(\boldsymbol{x}, \boldsymbol{y})<r\right\}
$$

is called an open ball with radius $r$ centred at $\boldsymbol{x}$ or the neighbourhood of $\boldsymbol{x}$.
Definition. Let $M \subset \mathbb{R}^{n}$. We say that $\boldsymbol{x} \in \mathbb{R}^{n}$ is an interior point of $M$, if there exists $r>0$ such that $B(\boldsymbol{x}, r) \subset M$.
The set of all interior points of $M$ is called the interior of $M$ and is denoted by $\operatorname{Int} M$.
The set $M \subset \mathbb{R}^{n}$ is open in $\mathbb{R}^{n}$, if each point of $M$ is an interior point of $M$, i.e. if $M=\operatorname{Int} M$.
Theorem 2 (properties of open sets).
(i) The empty set and $\mathbb{R}^{n}$ are open in $\mathbb{R}^{n}$.
(ii) Let $G_{\alpha} \subset \mathbb{R}^{n}, \alpha \in A \neq \emptyset$, be open in $\mathbb{R}^{n}$. Then $\bigcup_{\alpha \in A} G_{\alpha}$ is open in $\mathbb{R}^{n}$.
(iii) Let $G_{i} \subset \mathbb{R}^{n}, i=1, \ldots, m$, be open in $\mathbb{R}^{n}$. Then $\bigcap_{i=1}^{m} G_{i}$ is open in $\mathbb{R}^{n}$.

## Remark.

(ii) A union of an arbitrary system of open sets is an open set.
(iii) An intersection of a finitely many open sets is an open set.

Definition. Let $M \subset \mathbb{R}^{n}$ and $\boldsymbol{x} \in \mathbb{R}^{n}$. We say that $\boldsymbol{x}$ is a boundary point of $M$ if for each $r>0$

$$
B(\boldsymbol{x}, r) \cap M \neq \emptyset \quad \text { and } \quad B(\boldsymbol{x}, r) \cap\left(\mathbb{R}^{n} \backslash M\right) \neq \emptyset
$$

The boundary of $M$ is the set of all boundary points of $M$ (notation bd $M$ ).
The closure of $M$ is the set $M \cup \mathrm{bd} M$ (notation $\bar{M}$ ).
A set $M \subset \mathbb{R}^{n}$ is said to be closed in $\mathbb{R}^{n}$ if it contains all its boundary points, i.e. if bd $M \subset M$, or in other words if $\bar{M}=M$.
Definition. Let $\boldsymbol{x}^{j} \in \mathbb{R}^{n}$ for each $j \in \mathbb{N}$ and $\boldsymbol{x} \in \mathbb{R}^{n}$. We say that a sequence $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$ converges to $\boldsymbol{x}$, if

$$
\lim _{j \rightarrow \infty} \rho\left(\boldsymbol{x}, \boldsymbol{x}^{j}\right)=0
$$

The vector $\boldsymbol{x}$ is called the limit of the sequence $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$.
The sequence $\left\{\boldsymbol{y}^{j}\right\}_{j=1}^{\infty}$ of points in $\mathbb{R}^{n}$ is called convergent if there exists $\boldsymbol{y} \in \mathbb{R}^{n}$ such that $\left\{\boldsymbol{y}^{j}\right\}_{j=1}^{\infty}$ converges to $\boldsymbol{y}$.

Remark. The sequence $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$ converges to $\boldsymbol{x} \in \mathbb{R}^{n}$ if and only if

$$
\forall \varepsilon \in \mathbb{R}, \varepsilon>0 \exists j_{0} \in \mathbb{N} \forall j \in \mathbb{N}, j \geq j_{0}: \boldsymbol{x}^{j} \in B(\boldsymbol{x}, \varepsilon) .
$$

Theorem 3 (convergence is coordinatewise). Let $\boldsymbol{x}^{j} \in \mathbb{R}^{n}$ for each $j \in \mathbb{N}$ and let $\boldsymbol{x} \in \mathbb{R}^{n}$. The sequence $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$ converges to $\boldsymbol{x}$ if and only iffor each $i \in\{1, \ldots, n\}$ the sequence of real numbers $\left\{x_{i}^{j}\right\}_{j=1}^{\infty}$ converges to the real number $x_{i}$.

Remark. Theorem 3 says that the convergence in the space $\mathbb{R}^{n}$ is the same as the "coordinatewise" convergence. It follows that a sequence $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$ has at most one limit. If it exists, then we denote it by $\lim _{j \rightarrow \infty} \boldsymbol{x}^{j}$. Sometimes we also write simply $\boldsymbol{x}^{j} \rightarrow \boldsymbol{x}$ instead of $\lim _{j \rightarrow \infty} \boldsymbol{x}^{j}=\boldsymbol{x}$.

Theorem 4 (characterisation of closed sets). Let $M \subset \mathbb{R}^{n}$. Then the following statements are equivalent:
(i) $M$ is closed in $\mathbb{R}^{n}$.
(ii) $\mathbb{R}^{n} \backslash M$ is open in $\mathbb{R}^{n}$.
(iii) Any $\boldsymbol{x} \in \mathbb{R}^{n}$ which is a limit of a sequence from $M$ belongs to $M$.

Theorem 5 (properties of closed sets).
(i) The empty set and the whole space $\mathbb{R}^{n}$ are closed in $\mathbb{R}^{n}$.
(ii) Let $F_{\alpha} \subset \mathbb{R}^{n}, \alpha \in A \neq \emptyset$, be closed in $\mathbb{R}^{n}$. Then $\bigcap_{\alpha \in A} F_{\alpha}$ is closed in $\mathbb{R}^{n}$.
(iii) Let $F_{i} \subset \mathbb{R}^{n}, i=1, \ldots, m$, be closed in $\mathbb{R}^{n}$. Then $\bigcup_{i=1}^{m} F_{i}$ is closed in $\mathbb{R}^{n}$.

Remark.
(ii) An intersection of an arbitrary system of closed sets is closed.
(iii) A union of finitely many closed sets is closed.

Theorem 6. Let $M \subset \mathbb{R}^{n}$. Then the following holds:
(i) The set $\bar{M}$ is closed in $\mathbb{R}^{n}$.
(ii) The set $\operatorname{Int} M$ is open in $\mathbb{R}^{n}$.
(iii) The set $M$ is open in $\mathbb{R}^{n}$ if and only if $M=\operatorname{Int} M$.

Remark. The set Int $M$ is the largest open set contained in $M$ in the following sense: If $G$ is a set open in $\mathbb{R}^{n}$ and satisfying $G \subset M$, then $G \subset \operatorname{Int} M$. Similarly $\bar{M}$ is the smallest closed set containing $M$.

Definition. We say that the set $M \subset \mathbb{R}^{n}$ is bounded if there exists $r>0$ such that $M \subset B(\boldsymbol{o}, r)$. A sequence of points in $\mathbb{R}^{n}$ is bounded if the set of its members is bounded.

Theorem 7. A set $M \subset \mathbb{R}^{n}$ is bounded if and only if its closure $\bar{M}$ is bounded.

## V.2. Continuous functions of several variables

Definition. Let $M \subset \mathbb{R}^{n}, \boldsymbol{x} \in M$, and $f: M \rightarrow \mathbb{R}$. We say that $f$ is continuous at $\boldsymbol{x}$ with respect to $M$, if we

$$
\forall \varepsilon \in \mathbb{R}, \varepsilon>0 \exists \delta \in \mathbb{R}, \delta>0 \forall \boldsymbol{y} \in B(\boldsymbol{x}, \delta) \cap M: f(\boldsymbol{y}) \in B(f(\boldsymbol{x}), \varepsilon) .
$$

We say that $f$ is continuous at the point $\boldsymbol{x}$ if it is continuous at $\boldsymbol{x}$ with respect to a neighbourhood of $\boldsymbol{x}$, i.e.

$$
\forall \varepsilon \in \mathbb{R}, \varepsilon>0 \exists \delta \in \mathbb{R}, \delta>0 \forall \boldsymbol{y} \in B(\boldsymbol{x}, \delta): f(\boldsymbol{y}) \in B(f(\boldsymbol{x}), \varepsilon) .
$$

Theorem 8. Let $M \subset \mathbb{R}^{n}, \boldsymbol{x} \in M, f: M \rightarrow \mathbb{R}, g: M \rightarrow \mathbb{R}$, and $c \in \mathbb{R}$. If $f$ and $g$ are continuous at the point $\boldsymbol{x}$ with respect to $M$, then the functions $c f, f+g$ a fg are continuous at $\boldsymbol{x}$ with respect to $M$. If the function $g$ is nonzero at $\boldsymbol{x}$, then also the function $f / g$ is continuous at $\boldsymbol{x}$ with respect to $M$.

Theorem 9. Let $r, s \in \mathbb{N}, M \subset \mathbb{R}^{s}, L \subset \mathbb{R}^{r}$, and $\boldsymbol{y} \in M$. Let $\varphi_{1}, \ldots, \varphi_{r}$ be functions defined on $M$, which are continuous at $\boldsymbol{y}$ with respect to $M$ and $\left[\varphi_{1}(\boldsymbol{x}), \ldots, \varphi_{r}(\boldsymbol{x})\right] \in L$ for each $\boldsymbol{x} \in M$. Let $f: L \rightarrow \mathbb{R}$ be continuous at the point $\left[\varphi_{1}(\boldsymbol{y}), \ldots, \varphi_{r}(\boldsymbol{y})\right]$ with respect to $L$. Then the compound function $F: M \rightarrow \mathbb{R}$ defined by

$$
F(\boldsymbol{x})=f\left(\varphi_{1}(\boldsymbol{x}), \ldots, \varphi_{r}(\boldsymbol{x})\right), \quad \boldsymbol{x} \in M,
$$

is continuous at $\boldsymbol{y}$ with respect to $M$.

Theorem 10 (Heine). Let $M \subset \mathbb{R}^{n}, \boldsymbol{x} \in M$, and $f: M \rightarrow \mathbb{R}$. Then the following are equivalent.
(i) The function $f$ is continuous at $\boldsymbol{x}$ with respect to $M$.
(ii) $\lim _{j \rightarrow \infty} f\left(\boldsymbol{x}^{j}\right)=f(\boldsymbol{x})$ for each sequence $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$ such that $\boldsymbol{x}^{j} \in M$ for $j \in \mathbb{N}$ and $\lim _{j \rightarrow \infty} \boldsymbol{x}^{j}=\boldsymbol{x}$.

Definition. Let $M \subset \mathbb{R}^{n}$ and $f: M \rightarrow \mathbb{R}$. We say that $f$ is continuous on $M$ if it is continuous at each point $\boldsymbol{x} \in M$ with respect to $M$.

Remark. The functions $\pi_{j}: \mathbb{R}^{n} \rightarrow \mathbb{R}, \pi_{j}(\boldsymbol{x})=x_{j}, 1 \leq j \leq n$, are continuous on $\mathbb{R}^{n}$. They are called coordinate projections.
Theorem 11. Let $f$ be a continuous function on $\mathbb{R}^{n}$ and $c \in \mathbb{R}$. Then the following holds:
(i) The set $\left\{\boldsymbol{x} \in \mathbb{R}^{n} ; f(\boldsymbol{x})<c\right\}$ is open in $\mathbb{R}^{n}$.
(ii) The set $\left\{\boldsymbol{x} \in \mathbb{R}^{n} ; f(\boldsymbol{x})>c\right\}$ is open in $\mathbb{R}^{n}$.
(iii) The set $\left\{\boldsymbol{x} \in \mathbb{R}^{n} ; f(\boldsymbol{x}) \leq c\right\}$ is closed in $\mathbb{R}^{n}$.
(iv) The set $\left\{\boldsymbol{x} \in \mathbb{R}^{n} ; f(\boldsymbol{x}) \geq c\right\}$ is closed in $\mathbb{R}^{n}$.
(v) The set $\left\{\boldsymbol{x} \in \mathbb{R}^{n} ; f(\boldsymbol{x})=c\right\}$ is closed in $\mathbb{R}^{n}$.

Definition. We say that a set $M \subset \mathbb{R}^{n}$ is compact if for each sequence of elements of $M$ there exists a convergent subsequence with a limit in $M$.

Theorem 12 (characterisation of compact subsets of $\mathbb{R}^{n}$ ). The set $M \subset \mathbb{R}^{n}$ is compact if and only if $M$ is bounded and closed.
Lemma 13. Let $\left\{\boldsymbol{x}^{j}\right\}_{j=1}^{\infty}$ be a bounded sequence in $\mathbb{R}^{n}$. Then it has a convergent subsequence.
Definition. Let $M \subset \mathbb{R}^{n}, \boldsymbol{x} \in M$, and let $f$ be a function defined at least on $M$ (i.e. $M \subset D_{f}$ ). We say that $f$ attains at the point $\boldsymbol{x}$ its

- maximum on $M$ if $f(\boldsymbol{y}) \leq f(\boldsymbol{x})$ for every $\boldsymbol{y} \in M$,
- local maximum with respect to $M$ if there exists $\delta>0$ such that $f(\boldsymbol{y}) \leq f(\boldsymbol{x})$ for every $\boldsymbol{y} \in B(\boldsymbol{x}, \delta) \cap M$,
- strict local maximum with respect to $M$ if there exists $\delta>0$ such that $f(\boldsymbol{y})<f(\boldsymbol{x})$ for every $\boldsymbol{y} \in(B(\boldsymbol{x}, \delta) \backslash\{\boldsymbol{x}\}) \cap M$. The notions of a minimum, a local minimum, and a strict local minimum with respect to $M$ are defined in analogous way.

Definition. We say that a function $f$ attains a local maximum at a point $\boldsymbol{x} \in \mathbb{R}^{n}$ if $\boldsymbol{x}$ is a local maximum with respect to some neighbourhood of $\boldsymbol{x}$.

Similarly we define local minimum, strict local maximum and strict local minimum.
Theorem 14 (attaining extrema). Let $M \subset \mathbb{R}^{n}$ be a non-empty compact set and $f: M \rightarrow \mathbb{R}$ a function continuous on $M$. Then $f$ attains its maximum and minimum on $M$.

Corollary. Let $M \subset \mathbb{R}^{n}$ be a non-empty compact set and $f: M \rightarrow \mathbb{R}$ a continuous function on $M$. Then $f$ is bounded on $M$.
Definition. We say that a function $f$ of $n$ variables has a limit at a point $\boldsymbol{a} \in \mathbb{R}^{n}$ equal to $A \in \mathbb{R}^{*}$ if

$$
\forall \varepsilon \in \mathbb{R}, \varepsilon>0 \exists \delta \in \mathbb{R}, \delta>0 \forall x \in B(\boldsymbol{a}, \delta) \backslash\{a\}: f(\boldsymbol{x}) \in B(A, \varepsilon) .
$$

Remark.

- Each function has at a given point at most one limit. We write $\lim _{\boldsymbol{x} \rightarrow \boldsymbol{a}} f(\boldsymbol{x})=A$.
- The function $f$ is continuous at $\boldsymbol{a}$ if and only if $\lim _{\boldsymbol{x} \rightarrow \boldsymbol{a}} f(\boldsymbol{x})=f(\boldsymbol{a})$.
- For limits of functions of several variables one can prove similar theorems as for limits of functions of one variable (arithmetics, the sandwich theorem, ...).

Theorem 15. Let $r, s \in \mathbb{N}, \boldsymbol{a} \in \mathbb{R}^{s}$, and let $\varphi_{1}, \ldots, \varphi_{r}$ be functions of $s$ variables such that $\lim _{\boldsymbol{x} \rightarrow \boldsymbol{a}} \varphi_{j}(\boldsymbol{x})=b_{j}, j=1, \ldots, r$. Set $\boldsymbol{b}=\left[b_{1}, \ldots, b_{r}\right]$. Let $f$ be a function of $r$ variables which is continuous at the point $\boldsymbol{b}$. If we define a compound function $F$ of $s$ variables by

$$
F(\boldsymbol{x})=f\left(\varphi_{1}(\boldsymbol{x}), \varphi_{2}(\boldsymbol{x}), \ldots, \varphi_{r}(\boldsymbol{x})\right),
$$

then $\lim _{\boldsymbol{x} \rightarrow \boldsymbol{a}} F(\boldsymbol{x})=f(\boldsymbol{b})$.

## V.3. Partial derivatives and tangent hyperplane

$$
\text { Set } \boldsymbol{e}^{j}=[0, \ldots, 0, \underset{j \text { th coordinate }}{1}, 0, \ldots, 0] \text {. }
$$

Definition. Let $f$ be a function of $n$ variables, $j \in\{1, \ldots, n\}, \boldsymbol{a} \in \mathbb{R}^{n}$. Then the number

$$
\begin{aligned}
\frac{\partial f}{\partial x_{j}}(\boldsymbol{a}) & =\lim _{t \rightarrow 0} \frac{f\left(\boldsymbol{a}+t \boldsymbol{e}^{j}\right)-f(\boldsymbol{a})}{t} \\
& =\lim _{t \rightarrow 0} \frac{f\left(a_{1}, \ldots, a_{j-1}, a_{j}+t, a_{j+1}, \ldots, a_{n}\right)-f\left(a_{1}, \ldots, a_{n}\right)}{t}
\end{aligned}
$$

is called the partial derivative (of first order) of function $f$ according to $j$ th variable at the point $\boldsymbol{a}$ (if the limit exists).
Theorem 16 (necessary condition of the existence of local extremum). Let $G \subset \mathbb{R}^{n}$ be an open set, $\boldsymbol{a} \in G$, and suppose that a function $f: G \rightarrow \mathbb{R}$ has a local extremum (i.e. a local maximum or a local minimum) at the point a. Then for each $j \in\{1, \ldots, n\}$ the following holds:

The partial derivative $\frac{\partial f}{\partial x_{j}}(\boldsymbol{a})$ either does not exist or it is equal to zero.
Definition. Let $G \subset \mathbb{R}^{n}$ be a non-empty open set. If a function $f: G \rightarrow \mathbb{R}$ has all partial derivatives continuous at each point of the set $G$ (i.e. the function $\boldsymbol{x} \mapsto \frac{\partial f}{\partial x_{j}}(\boldsymbol{x})$ is continuous on $G$ for each $j \in\{1, \ldots, n\}$ ), then we say that $f$ is of the class $\mathcal{C}^{1}$ on $G$. The set of all of these functions is denoted by $C^{1}(G)$.

Remark. If $G \subset \mathbb{R}^{n}$ is a non-empty open set and and $f, g \in C^{1}(G)$, then $f+g \in C^{1}(G), f-g \in C^{1}(G)$, and $f g \in C^{1}(G)$. If moreover $g(\boldsymbol{x}) \neq 0$ for each $\boldsymbol{x} \in G$, then $f / g \in C^{1}(G)$.

Proposition 17 (weak Lagrange theorem). Let $n \in \mathbb{N}, I_{1}, \ldots, I_{n} \subset \mathbb{R}$ be open intervals, $I=I_{1} \times I_{2} \times \cdots \times I_{n}, f \in C^{1}(I)$, and $\boldsymbol{a}, \boldsymbol{b} \in I$. Then there exist points $\boldsymbol{\xi}^{1}, \ldots, \boldsymbol{\xi}^{n} \in I$ with $\xi_{j}^{i} \in\left[a_{j}, b_{j}\right]$ for each $i, j \in\{1, \ldots, n\}$, such that

$$
f(\boldsymbol{b})-f(\boldsymbol{a})=\sum_{i=1}^{n} \frac{\partial f}{\partial x_{i}}\left(\xi^{i}\right)\left(b_{i}-a_{i}\right)
$$

Definition. Let $G \subset \mathbb{R}^{n}$ be an open set, $\boldsymbol{a} \in G$, and $f \in C^{1}(G)$. Then the graph of the function

$$
T: \boldsymbol{x} \mapsto f(\boldsymbol{a})+\frac{\partial f}{\partial x_{1}}(\boldsymbol{a})\left(x_{1}-a_{1}\right)+\frac{\partial f}{\partial x_{2}}(\boldsymbol{a})\left(x_{2}-a_{2}\right)+\cdots+\frac{\partial f}{\partial x_{n}}(\boldsymbol{a})\left(x_{n}-a_{n}\right), \quad \boldsymbol{x} \in \mathbb{R}^{n}
$$

is called the tangent hyperplane to the graph of the function $f$ at the point $[\boldsymbol{a}, f(\boldsymbol{a})]$.
Theorem 18 (tangent hyperplane). Let $G \subset \mathbb{R}^{n}$ be an open set, $a \in G, f \in C^{1}(G)$, and let $T$ be a function whose graph is the tangent hyperplane of the function $f$ at the point $[\boldsymbol{a}, f(\boldsymbol{a})]$. Then

$$
\lim _{\boldsymbol{x} \rightarrow \boldsymbol{a}} \frac{f(\boldsymbol{x})-T(\boldsymbol{x})}{\rho(\boldsymbol{x}, \boldsymbol{a})}=0
$$

Theorem 19. Let $G \subset \mathbb{R}^{n}$ be an open non-empty set and $f \in C^{1}(G)$. Then $f$ is continuous on $G$.
Theorem 20 (derivative of a compound function; chain rule). Let $r, s \in \mathbb{N}$ and let $G \subset \mathbb{R}^{s}, H \subset \mathbb{R}^{r}$ be open sets. Let $\varphi_{1}, \ldots, \varphi_{r} \in C^{1}(G), f \in C^{1}(H)$ and $\left[\varphi_{1}(\boldsymbol{x}), \ldots, \varphi_{r}(\boldsymbol{x})\right] \in H$ for each $\boldsymbol{x} \in G$. Then the compound function $F: G \rightarrow \mathbb{R}$ defined by

$$
F(\boldsymbol{x})=f\left(\varphi_{1}(\boldsymbol{x}), \varphi_{2}(\boldsymbol{x}), \ldots, \varphi_{r}(\boldsymbol{x})\right), \quad \boldsymbol{x} \in G,
$$

is of the class $\mathcal{C}^{1}$ on $G$. Let $\boldsymbol{a} \in G$ and $\boldsymbol{b}=\left[\varphi_{1}(\boldsymbol{a}), \ldots, \varphi_{r}(\boldsymbol{a})\right]$. Then for each $j \in\{1, \ldots, s\}$ we have

$$
\frac{\partial F}{\partial x_{j}}(\boldsymbol{a})=\sum_{i=1}^{r} \frac{\partial f}{\partial y_{i}}(\boldsymbol{b}) \frac{\partial \varphi_{i}}{\partial x_{j}}(\boldsymbol{a}) .
$$

Definition. Let $G \subset \mathbb{R}^{n}$ be an open set, $\boldsymbol{a} \in G$, and $f \in C^{1}(G)$. The gradient of $f$ at the point $\boldsymbol{a}$ is the vector

$$
\nabla f(\boldsymbol{a})=\left[\frac{\partial f}{\partial x_{1}}(\boldsymbol{a}), \frac{\partial f}{\partial x_{2}}(\boldsymbol{a}), \ldots, \frac{\partial f}{\partial x_{n}}(\boldsymbol{a})\right] .
$$

Definition. Let $G \subset \mathbb{R}^{n}$ be an open set, $\boldsymbol{a} \in G, f \in C^{1}(G)$, and $\nabla f(\boldsymbol{a})=\boldsymbol{o}$. Then the point $\boldsymbol{a}$ is called a stationary (or critical) point of the function $f$.

Definition. Let $G \subset \mathbb{R}^{n}$ be an open set, $f: G \rightarrow \mathbb{R}, i, j \in\{1, \ldots, n\}$, and suppose that $\frac{\partial f}{\partial x_{i}}(\boldsymbol{x})$ exists finite for each $\boldsymbol{x} \in G$. Then the partial derivative of the second order of the function $f$ according to $i$ th and $j$ th variable at a point $\boldsymbol{a} \in G$ is defined by

$$
\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}(\boldsymbol{a})=\frac{\partial\left(\frac{\partial f}{\partial x_{i}}\right)}{\partial x_{j}}(\boldsymbol{a})
$$

If $i=j$ then we use the notation $\frac{\partial^{2} f}{\partial x_{i}^{2}}(\boldsymbol{a})$.
Similarly we define higher order partial derivatives.
Remark. In general it is not true that $\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}(\boldsymbol{a})=\frac{\partial^{2} f}{\partial x_{j} \partial x_{i}}(\boldsymbol{a})$.
Theorem 21 (interchanging of partial derivatives). Let $i, j \in\{1, \ldots, n\}$ and suppose that a function $f$ has both partial derivatives $\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}$ and $\frac{\partial^{2} f}{\partial x_{j} \partial x_{i}}$ on a neighbourhood of a point $\boldsymbol{a} \in \mathbb{R}^{n}$ and that these functions are continuous at $\boldsymbol{a}$. Then

$$
\frac{\partial^{2} f}{\partial x_{i} \partial x_{j}}(\boldsymbol{a})=\frac{\partial^{2} f}{\partial x_{j} \partial x_{i}}(\boldsymbol{a}) .
$$

Definition. Let $G \subset \mathbb{R}^{n}$ be an open set and $k \in \mathbb{N}$. We say that a function $f$ is of the class $C^{k}$ on $G$, if all partial derivatives of $f$ of all orders up to $k$ are continuous on $G$. The set of all of these functions is denoted by $C^{k}(G)$.

We say that a function $f$ is of the class $\mathcal{C}^{\infty}$ on $G$, if all partial derivatives of all orders of $f$ are continuous on $G$. The set of all of these functions is denoted by $C^{\infty}(G)$.

## V.4. Implicit function theorem

Theorem 22 (implicit function). Let $G \subset \mathbb{R}^{n+1}$ be an open set, $F: G \rightarrow \mathbb{R}$, and $\tilde{\boldsymbol{x}} \in \mathbb{R}^{n}, \tilde{y} \in \mathbb{R}$ such that $[\tilde{\boldsymbol{x}}, \tilde{y}] \in G$. Suppose that
(i) $F \in C^{1}(G)$,
(ii) $F(\tilde{\boldsymbol{x}}, \tilde{y})=0$,
(iii) $\frac{\partial F}{\partial y}(\tilde{\boldsymbol{x}}, \tilde{y}) \neq 0$.

Then there exist a neighbourhood $U \subset \mathbb{R}^{n}$ of the point $\tilde{\boldsymbol{x}}$ and a neighbourhood $V \subset \mathbb{R}$ of the point $\tilde{y}$ such that for each $\boldsymbol{x} \in U$ there exists a unique $y \in V$ satisfying $F(\boldsymbol{x}, y)=0$. If we denote this $y$ by $\varphi(\boldsymbol{x})$, then the resulting function $\varphi$ is in $C^{1}(U)$ and

$$
\frac{\partial \varphi}{\partial x_{j}}(\boldsymbol{x})=-\frac{\frac{\partial F}{\partial x_{j}}(\boldsymbol{x}, \varphi(\boldsymbol{x}))}{\frac{\partial F}{\partial y}(\boldsymbol{x}, \varphi(\boldsymbol{x}))} \quad \text { for } \boldsymbol{x} \in U, j \in\{1, \ldots, n\} .
$$

Theorem 23 (implicit functions). Let $m, n \in \mathbb{N}, k \in \mathbb{N} \cup\{\infty\}, G \subset \mathbb{R}^{n+m}$ an open set, $F_{j}: G \rightarrow \mathbb{R}$ for $j=1, \ldots, m, \tilde{\boldsymbol{x}} \in \mathbb{R}^{n}$, $\tilde{\boldsymbol{y}} \in \mathbb{R}^{m},[\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{y}}] \in G$. Suppose that
(i) $F_{j} \in C^{k}(G)$ for all $j \in\{1, \ldots, m\}$,
(ii) $F_{j}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{y}})=0$ for all $j \in\{1, \ldots, m\}$,
(iii) $\left|\begin{array}{ccc}\frac{\partial F_{1}}{\partial y_{1}}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{y}}) & \ldots & \frac{\partial F_{1}}{\partial y_{m}}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{y}}) \\ \vdots & \ddots & \vdots \\ \frac{\partial F_{m}}{\partial y_{1}}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{y}}) & \ldots & \frac{\partial F_{m}}{\partial y_{m}}(\tilde{\boldsymbol{x}}, \tilde{\boldsymbol{y}})\end{array}\right| \neq 0$.

Then there are a neighbourhood $U \subset \mathbb{R}^{n}$ of $\tilde{\boldsymbol{x}}$ and a neighbourhood $V \subset \mathbb{R}^{m}$ of $\tilde{\boldsymbol{y}}$ such that for each $\boldsymbol{x} \in U$ there exists a unique $\boldsymbol{y} \in V$ satisfying $F_{j}(\boldsymbol{x}, \boldsymbol{y})=0$ for each $j \in\{1, \ldots, m\}$. If we denote the coordinates of this $\boldsymbol{y}$ by $\varphi_{j}(\boldsymbol{x})$, then the resulting functions $\varphi_{j}$ are in $C^{k}(U)$.

Remark. The symbol in the condition (iii) of Theorem 23 is called a determinant. The general definition will be given later.
For $m=1$ we have $|a|=a, a \in \mathbb{R}$. In particular, in this case the condition (iii) in Theorem 23 is the same as the condition (iii) in Theorem 22.

For $m=2$ we have $\left|\begin{array}{ll}a & b \\ c & d\end{array}\right|=a d-b c, a, b, c, d \in \mathbb{R}$.

## V.5. Lagrange multipliers theorem

Theorem 24 (Lagrange multiplier theorem). Let $G \subset \mathbb{R}^{2}$ be an open set, $f, g \in C^{1}(G), M=\{[x, y] \in G ; g(x, y)=0\}$ and let $[\tilde{x}, \tilde{y}] \in M$ be a point of local extremum of $f$ with respect to $M$. Then at least one of the following conditions holds:
(I) $\nabla g(\tilde{x}, \tilde{y})=\boldsymbol{o}$,
(II) there exists $\lambda \in \mathbb{R}$ satisfying

$$
\begin{aligned}
& \frac{\partial f}{\partial x}(\tilde{x}, \tilde{y})+\lambda \frac{\partial g}{\partial x}(\tilde{x}, \tilde{y})=0 \\
& \frac{\partial f}{\partial y}(\tilde{x}, \tilde{y})+\lambda \frac{\partial g}{\partial y}(\tilde{x}, \tilde{y})=0
\end{aligned}
$$

Theorem 25 (Lagrange multipliers theorem). Let $m, n \in \mathbb{N}, m<n, G \subset \mathbb{R}^{n}$ an open set, $f, g_{1}, \ldots, g_{m} \in C^{1}(G)$,

$$
M=\left\{z \in G ; g_{1}(z)=0, g_{2}(z)=0, \ldots, g_{m}(z)=0\right\}
$$

and let $\tilde{\boldsymbol{z}} \in M$ be a point of local extremum of $f$ with respect to the set $M$. Then at least one of the following conditions holds:
(I) the vectors

$$
\nabla g_{1}(\tilde{\boldsymbol{z}}), \nabla g_{2}(\tilde{\boldsymbol{z}}), \ldots, \nabla g_{m}(\tilde{\boldsymbol{z}})
$$

are linearly dependent,
(II) there exist numbers $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m} \in \mathbb{R}$ satisfying

$$
\nabla f(\tilde{\boldsymbol{z}})+\lambda_{1} \nabla g_{1}(\tilde{\boldsymbol{z}})+\lambda_{2} \nabla g_{2}(\tilde{\boldsymbol{z}})+\cdots+\lambda_{m} \nabla g_{m}(\tilde{\boldsymbol{z}})=\boldsymbol{o}
$$

Remark.

- The notion of linearly dependent vectors will be defined later.

For $m=1$ : One vector is linearly dependent if it is the zero vector.
For $m=2$ : Two vectors are linearly dependent if one of them is a multiple of the other one.

- The numbers $\lambda_{1}, \ldots, \lambda_{m}$ are called the Lagrange multipliers.


## V.6. Concave and quasiconcave functions

Definition. Let $M \subset \mathbb{R}^{n}$. We say that $M$ is convex if

$$
\forall \boldsymbol{x}, \boldsymbol{y} \in M \forall t \in[0,1]: t \boldsymbol{x}+(1-t) \boldsymbol{y} \in M .
$$

Definition. Let $M \subset \mathbb{R}^{n}$ be a convex set and $f$ a function defined on $M$. We say that $f$ is

- concave on $M$ if

$$
\forall \boldsymbol{a}, \boldsymbol{b} \in M \forall t \in[0,1]: f(t \boldsymbol{a}+(1-t) \boldsymbol{b}) \geq t f(\boldsymbol{a})+(1-t) f(\boldsymbol{b}),
$$

- strictly concave on $M$ if

$$
\forall \boldsymbol{a}, \boldsymbol{b} \in M, \boldsymbol{a} \neq \boldsymbol{b} \forall t \in(0,1): f(t \boldsymbol{a}+(1-t) \boldsymbol{b})>t f(\boldsymbol{a})+(1-t) f(\boldsymbol{b})
$$

Remark. By changing the inequalities to the opposite we obtain a definition of a convex and a strictly convex function.
Remark. A function $f$ is convex (strictly convex) if and only if the function $-f$ is concave (strictly concave).
All the theorems in this section are formulated for concave and strictly concave functions. They have obvious analogies that hold for convex and strictly convex functions.
Remark.

- If a function $f$ is strictly concave on $M$, then it is concave on $M$.
- Let $f$ be a concave function on $M$. Then $f$ is strictly concave on $M$ if and only if the graph of $f$ "does not contain a segment", i.e.

$$
\neg(\exists \boldsymbol{a}, \boldsymbol{b} \in M, \boldsymbol{a} \neq \boldsymbol{b}, \forall t \in[0,1]: f(t \boldsymbol{a}+(1-t) \boldsymbol{b})=t f(\boldsymbol{a})+(1-t) f(\boldsymbol{b})) .
$$

Theorem 26. Let $f$ be a function concave on an open convex set $G \subset \mathbb{R}^{n}$. Then $f$ is continuous on $G$.
Theorem 27. Let $f$ be a function concave on a convex set $M \subset \mathbb{R}^{n}$. Then for each $\alpha \in \mathbb{R}$ the set $Q_{\alpha}=\{\boldsymbol{x} \in M ; f(\boldsymbol{x}) \geq \alpha\}$ is convex.

Theorem 28 (characterisation of concave functions of the class $\varphi^{1}$ ). Let $G \subset \mathbb{R}^{n}$ be a convex open set and $f \in C^{1}(G)$. Then the function $f$ is concave on $G$ if and only if

$$
\forall \boldsymbol{x}, \boldsymbol{y} \in G: f(\boldsymbol{y}) \leq f(\boldsymbol{x})+\sum_{i=1}^{n} \frac{\partial f}{\partial x_{i}}(\boldsymbol{x})\left(y_{i}-x_{i}\right)
$$

Corollary 29. Let $G \subset \mathbb{R}^{n}$ be a convex open set and let $f \in C^{1}(G)$ be concave on $G$. If a point $\boldsymbol{a} \in G$ is a critical point of $f$ (i.e. $\nabla f(\boldsymbol{a})=\boldsymbol{o}$ ), then $\boldsymbol{a}$ is a point of maximum of $f$ on $G$.

Theorem 30 (characterisation of strictly concave functions of the class $\mathcal{C}^{1}$ ). Let $G \subset \mathbb{R}^{n}$ be a convex open set and $f \in C^{1}(G)$. Then the function $f$ is strictly concave on $G$ if and only if

$$
\forall \boldsymbol{x}, \boldsymbol{y} \in G, \boldsymbol{x} \neq \boldsymbol{y}: f(\boldsymbol{y})<f(\boldsymbol{x})+\sum_{i=1}^{n} \frac{\partial f}{\partial x_{i}}(\boldsymbol{x})\left(y_{i}-x_{i}\right) .
$$

Definition. Let $M \subset \mathbb{R}^{n}$ be a convex set and let $f$ be a function defined on $M$. We say that $f$ is

- quasiconcave na $M$ if

$$
\forall \boldsymbol{a}, \boldsymbol{b} \in M \forall t \in[0,1]: f(t \boldsymbol{a}+(1-t) \boldsymbol{b}) \geq \min \{f(\boldsymbol{a}), f(\boldsymbol{b})\},
$$

- strictly quasiconcave on $M$ if

$$
\forall \boldsymbol{a}, \boldsymbol{b} \in M, \boldsymbol{a} \neq \boldsymbol{b}, \forall t \in(0,1): f(t \boldsymbol{a}+(1-t) \boldsymbol{b})>\min \{f(\boldsymbol{a}), f(\boldsymbol{b})\}
$$

Remark. By changing the inequalities to the opposite and changing the minimum to a maximum we obtain a definition of a quasiconvex and a strictly quasiconvex function.
Remark. A function $f$ is quasiconvex (strictly quasiconvex) if and only if the function $-f$ is quasiconcave (strictly quasiconcave).

All the theorems in this section are formulated for quasiconcave and strictly quasiconcave functions. They have obvious analogies that hold for quasiconvex and strictly quasiconvex functions.
Remark.

- If a function $f$ is strictly quasiconcave on $M$, then it is quasiconcave on $M$.
- Let $f$ be a quasiconcave function on $M$. Then $f$ is strictly quasiconcave on $M$ if and only if the graph of $f$ "does not contain a horizontal segment", i.e.

$$
\neg(\exists \boldsymbol{a}, \boldsymbol{b} \in M, \boldsymbol{a} \neq \boldsymbol{b}, \forall t \in[0,1]: f(t \boldsymbol{a}+(1-t) \boldsymbol{b})=f(\boldsymbol{a})) .
$$

Remark. Let $M \subset \mathbb{R}^{n}$ be a convex set and $f$ a function defined on $M$.

- If $f$ is concave on $M$, then $f$ is quasiconcave on $M$.
- If $f$ is strictly concave on $M$, then $f$ is strictly quasiconcave on $M$.

Theorem 31 (a uniqueness of an extremum). Let $f$ be a strictly quasiconcave function on a convex set $M \subset \mathbb{R}^{n}$. Then there exists at most one point of maximum of $f$.

Corollary. Let $M \subset \mathbb{R}^{n}$ be a convex, closed, bounded and nonempty set and $f$ a continuous and strictly quasiconcave function on $M$. Then $f$ attains its maximum at exactly one point.

Theorem 32 (characterization of quasiconcave functions using level sets). Let $M \subset \mathbb{R}^{n}$ be a convex set and $f$ a function defined on $M$. Then $f$ is quasiconcave on $M$ if and only iffor each $\alpha \in \mathbb{R}$ the set $Q_{\alpha}=\{\boldsymbol{x} \in M ; f(\boldsymbol{x}) \geq \alpha\}$ is convex.

## VI. Matrix calculus

## VI.1. Basic operations with matrices

Definition. A table of numbers

$$
\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
a_{21} & a_{22} & \ldots & a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & a_{m 2} & \ldots & a_{m n}
\end{array}\right)
$$

where $a_{i j} \in \mathbb{R}, i=1, \ldots, m, j=1, \ldots, n$, is called a matrix of type $m \times n$ (shortly, an $m$-by-n matrix). We also write $\left(a_{i j}\right)_{\substack{i=1 . m \\ j=1 . . n}}$ for short.

An $n$-by- $n$ matrix is called a square matrix of order $n$.
The set of all $m$-by- $n$ matrices is denoted by $M(m \times n)$.
Definition. Let

$$
A=\left(\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
a_{21} & a_{22} & \ldots & a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1} & a_{m 2} & \ldots & a_{m n}
\end{array}\right)
$$

The $n$-tuple $\left(a_{i 1}, a_{i 2}, \ldots, a_{i n}\right)$, where $i \in\{1,2, \ldots, m\}$, is called the $i$ th row of the matrix $A$.
The $m$-tuple $\left(\begin{array}{c}a_{1 j} \\ a_{2 j} \\ \vdots \\ a_{m j}\end{array}\right)$, where $j \in\{1,2, \ldots, n\}$, is called the $j$ th column of the matrix $A$.
Definition. We say that two matrices are equal, if they are of the same type and the corresponding elements are equal, i.e. if $A=$ $\left(a_{i j}\right)_{\substack{i=1 . . m \\ j=1 \ldots n}}$ and $\mathbb{B}=\left(b_{u v}\right)_{\substack{u=1 . . r \\ v=1 . . s}}$, then $\mathbb{A}=\mathbb{B}$ if and only if $m=r, n=s$ and $a_{i j}=b_{i j} \forall i \in\{1, \ldots, m\}, \forall j \in\{1, \ldots, n\}$.
Definition. Let $A, \mathbb{B} \in M(m \times n), A=\left(a_{i j}\right)_{\substack{i=1 . . m \\ j=1 . . n}}, \mathbb{B}=\left(b_{i j}\right)_{\substack{i=1 . . m \\ j=1 . . n}}, \lambda \in \mathbb{R}$. The sum of the matrices $\mathbb{A}$ and $\mathbb{B}$ is the matrix defined by

$$
A+\mathbb{B}=\left(\begin{array}{cccc}
a_{11}+b_{11} & a_{12}+b_{12} & \ldots & a_{1 n}+b_{1 n} \\
a_{21}+b_{21} & a_{22}+b_{22} & \ldots & a_{2 n}+b_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
a_{m 1}+b_{m 1} & a_{m 2}+b_{m 1} & \ldots & a_{m n}+b_{m n}
\end{array}\right)
$$

The product of the real number $\lambda$ and the matrix $\mathbb{A}$ (or the $\lambda$-multiple of the matrix $A$ ) is the matrix defined by

$$
\lambda A=\left(\begin{array}{cccc}
\lambda a_{11} & \lambda a_{12} & \ldots & \lambda a_{1 n} \\
\lambda a_{21} & \lambda a_{22} & \ldots & \lambda a_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
\lambda a_{m 1} & \lambda a_{m 2} & \ldots & \lambda a_{m n}
\end{array}\right)
$$

Proposition 33 (basic properties of the sum of matrices and of a multiplication by a scalar). The following holds:

- $\forall A, \mathbb{B}, \mathbb{C} \in M(m \times n): A+(\mathbb{B}+\mathbb{C})=(\mathbb{A}+\mathbb{B})+\mathbb{C}$,
- $\forall A, \mathbb{B} \in M(m \times n): A+\mathbb{B}=\mathbb{B}+A$,
- $\exists!\mathbb{O} \in M(m \times n) \forall A \in M(m \times n): A+\mathbb{O}=A$, (existence of a zero element)
- $\forall \mathbb{A} \in M(m \times n) \exists \mathbb{C}_{A} \in M(m \times n): A+\mathbb{C}_{A}=\mathbb{O}$, (existence of an opposite element)
- $\forall \mathbb{A} \in M(m \times n) \forall \lambda, \mu \in \mathbb{R}:(\lambda \mu) \mathbb{A}=\lambda(\mu A)$,
- $\forall \mathbb{A} \in M(m \times n): 1 \cdot \mathbb{A}=\mathbb{A}$,
- $\forall A \in M(m \times n) \forall \lambda, \mu \in \mathbb{R}:(\lambda+\mu) A=\lambda A+\mu A$,
- $\forall A, \mathbb{B} \in M(m \times n) \forall \lambda \in \mathbb{R}: \lambda(A+\mathbb{B})=\lambda A+\lambda \mathbb{B}$.

Remark.

- The matrix $\mathbb{O}$ from the previous proposition is called a zero matrix and all its elements are all zeros.
- The matrix $\mathbb{C}_{A}$ from the previous proposition is called a matrix opposite to $A$. It is determined uniquely, it is denoted by $-A$, and it satisfies $-A=\left(-a_{i j}\right)_{\substack{i=1 . . m \\ j=1 . . n}}$ and $-A=-1 \cdot A$.

Definition. Let $\mathbb{A} \in M(m \times n), \mathbb{A}=\left(a_{i s}\right)_{\substack{i=1 \ldots m \\ s=1 . . n}}, \mathbb{B} \in M(n \times k), \mathbb{B}=\left(b_{s j}\right)_{\substack{s=1 \ldots . n \\ j=1 . . k}}$. Then the product of matrices $\mathbb{A}$ and $\mathbb{B}$ is defined as a matrix $\mathbb{A} \mathbb{B} \in M(m \times k), A \mathbb{B}=\left(c_{i j}\right)_{\substack{i=1 . . m \\ j=1 . . k}}$, where

$$
c_{i j}=\sum_{s=1}^{n} a_{i s} b_{s j}
$$

Theorem 34 (properties of the matrix multiplication). Let $m, n, k, l \in \mathbb{N}$. Then:
(i) $\forall \mathbb{A} \in M(m \times n) \forall \mathbb{B} \in M(n \times k) \forall \mathbb{C} \in M(k \times l): A(\mathbb{B} \mathbb{C})=(A \mathbb{B}) \mathbb{C}$,
(associativity of multiplication)
(ii) $\forall \mathbb{A} \in M(m \times n) \forall \mathbb{B}, \mathbb{C} \in M(n \times k): A(\mathbb{B}+\mathbb{C})=A \mathbb{B}+A \mathbb{C}$,
(distributivity from the left)
(iii) $\forall \mathbb{A}, \mathbb{B} \in M(m \times n) \forall \mathbb{C} \in M(n \times k):(A+\mathbb{B}) \mathbb{C}=A \mathbb{C}+\mathbb{B} \mathbb{C}$,
(distributivity from the right)
(iv) $\exists$ ! $\mathbb{I} \in M(n \times n) \forall A \in M(n \times n): \mathbb{I} A=A \mathbb{I}=A$. (existence and uniqueness of an identity matrix I)

Remark. Warning! The matrix multiplication is not commutative.
Definition. A transpose of a matrix

$$
A=\left(\begin{array}{ccccc}
a_{11} & a_{12} & a_{13} & \ldots & a_{1 n} \\
a_{21} & a_{22} & a_{23} & \ldots & a_{2 n} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
a_{m 1} & a_{m 2} & a_{m 3} & \ldots & a_{m n}
\end{array}\right)
$$

is the matrix

$$
A^{T}=\left(\begin{array}{cccc}
a_{11} & a_{21} & \ldots & a_{m 1} \\
a_{12} & a_{22} & \ldots & a_{m 2} \\
a_{13} & a_{23} & \ldots & a_{m 3} \\
\vdots & \vdots & \ddots & \vdots \\
a_{1 n} & a_{2 n} & \ldots & a_{m n}
\end{array}\right)
$$

i.e. if $A=\left(a_{i j}\right)_{\substack{i=1 . . m \\ j=1 . . n}}$, then $A^{T}=\left(b_{u v}\right)_{\substack{u=1 . . n \\ v=1 . . m}}$, where $b_{u v}=a_{v u}$ for each $u \in\{1, \ldots, n\}, v \in\{1,2, \ldots, m\}$.

Theorem 35 (properties of the transpose of a matrix). Platí:
(i) $\forall \mathbb{A} \in M(m \times n):\left(\mathbb{A}^{T}\right)^{T}=A$,
(ii) $\forall A, \mathbb{B} \in M(m \times n):(\mathbb{A}+\mathbb{B})^{T}=A^{T}+\mathbb{B}^{T}$,
(iii) $\forall \mathbb{A} \in M(m \times n) \forall \mathbb{B} \in M(n \times k):(\mathbb{A} \mathbb{B})^{T}=\mathbb{B}^{T} A^{T}$.

## VI.2. Invertible matrices

Definition. Let $\mathbb{A} \in M(n \times n)$. We say that $A$ is an invertible matrix if there exist $\mathbb{B} \in M(n \times n)$ such that

$$
A \mathbb{B}=\mathbb{B} A=\mathbb{I}
$$

Definition. We say that the matrix $\mathbb{B} \in M(n \times n)$ is an inverse of a matrix $A \in M(n \times n)$ if $A \mathbb{B}=\mathbb{B} \mathbb{A}=\mathbb{I}$.
Remark. A matrix $A \in M(n \times n)$ is invertible if and only if it has an inverse.
Remark.

- If $A \in M(n \times n)$ is invertible, then it has exactly one inverse, which is denoted by $A^{-1}$.
- If some matrices $A, \mathbb{B} \in M(n \times n)$ satisfy $A \mathbb{B}=\mathbb{I}$, then also $\mathbb{B} A=\mathbb{I}$.

Theorem 36 (operations with invertible matrices). Let $A, \mathbb{B} \in M(n \times n)$ be invertible matrices. Then
(i) $\mathbb{A}^{-1}$ is invertible and $\left(\mathbb{A}^{-1}\right)^{-1}=\mathbb{A}$,
(ii) $A^{T}$ is invertible and $\left(A^{T}\right)^{-1}=\left(A^{-1}\right)^{T}$,
(iii) $A \mathbb{B}$ is invertible and $(\mathbb{A} \mathbb{B})^{-1}=\mathbb{B}^{-1} \mathbb{A}^{-1}$.

Definition. Let $k, n \in \mathbb{N}$ and $\boldsymbol{v}^{1}, \ldots, \boldsymbol{v}^{k} \in \mathbb{R}^{n}$. We say that a vector $\boldsymbol{u} \in \mathbb{R}^{n}$ is a linear combination of the vectors $\boldsymbol{v}^{1}, \ldots, \boldsymbol{v}^{k}$ with coefficients $\lambda_{1}, \ldots, \lambda_{k} \in \mathbb{R}$ if

$$
\boldsymbol{u}=\lambda_{1} \boldsymbol{v}^{1}+\cdots+\lambda_{k} \boldsymbol{v}^{k} .
$$

By a trivial linear combination of vectors $\boldsymbol{v}^{1}, \ldots, \boldsymbol{v}^{k}$ we mean the linear combination $0 \cdot \boldsymbol{v}^{1}+\cdots+0 \cdot \boldsymbol{v}^{k}$. Linear combination which is not trivial is called non-trivial.

Definition. We say that vectors $\boldsymbol{v}^{1}, \ldots, \boldsymbol{v}^{k} \in \mathbb{R}^{n}$ are linearly dependent if there exists their non-trivial linear combination which is equal to the zero vector. We say that vectors $\boldsymbol{v}^{1}, \ldots, \boldsymbol{v}^{k} \in \mathbb{R}^{n}$ are linearly independent if they are not linearly dependent, i.e. if whenever $\lambda_{1}, \ldots, \lambda_{k} \in \mathbb{R}$ satisfy $\lambda_{1} \boldsymbol{v}^{1}+\cdots+\lambda_{k} \boldsymbol{v}^{k}=\boldsymbol{o}$, then $\lambda_{1}=\lambda_{2}=\cdots=\lambda_{k}=0$.

Remark. Vectors $\boldsymbol{v}^{1}, \ldots, \boldsymbol{v}^{k}$ are linearly dependent if and only if one of them can be expressed as a linear combination of the others.

Definition. Let $\mathbb{A} \in M(m \times n)$. The $\operatorname{rank}$ of the matrix $\mathbb{A}$ is the maximal number of linearly independent row vectors of $A$, i.e. the rank is equal to $k \in \mathbb{N}$ if
(i) there is $k$ linearly independent row vectors of $\mathbb{A}$ and
(ii) each $l$-tuple of row vectors of $A$, where $l>k$, is linearly dependent.

The rank of the zero matrix is zero. $\operatorname{Rank}$ of $\mathbb{A}$ is denoted by $\operatorname{rank}(\mathbb{A})$.
Definition. We say that a matrix $A \in M(m \times n)$ is in a row echelon form if for each $i \in\{2, \ldots, m\}$ the $i$ th row of $A$ is either a zero vector or it has more zeros at the beginning than the $(i-1)$ th row.

Remark. The rank of a row echelon matrix is equal to the number of its non-zero rows.
Definition. The elementary row operations on the matrix $A$ are:
(i) interchange of two rows,
(ii) multiplication of a row by a non-zero real number,
(iii) addition of a multiple of a row to another row.

Definition. A matrix transformation is a finite sequence of elementary row operations. If a matrix $\mathbb{B} \in M(m \times n)$ results from


Theorem 37 (properties of matrix transformations).
(i) Let $\mathbb{A} \in M(m \times n)$. Then there exists a transformation transforming $A$ to a row echelon matrix.
(ii) Let $T_{1}$ be a transformation applicable to m-by-n matrices. Then there exists a transformation $T_{2}$ applicable to m-by-n matrices such that for any two matrices $A, \mathbb{B} \in M(m \times n)$ we have $A \underset{\sim}{\sim} \mathbb{T}_{1}$ if and only if $\mathbb{B} \xrightarrow{T_{2}} A$.
(iii) Let $\mathbb{A}, \mathbb{B} \in M(m \times n)$ and there exist a transformation $T$ such that $\underset{A}{\sim} \xrightarrow{T} \mathbb{B}$. Then $\operatorname{rank}(A)=\operatorname{rank}(\mathbb{B})$.

Remark. Similarly as the elementary row operations one can define also elementary column operations. It can be shown that the elementary column operations do not change the rank of the matrix.
Remark. It can be shown that $\operatorname{rank}(\mathbb{A})=\operatorname{rank}\left(A^{T}\right)$ for any $\mathbb{A} \in M(m \times n)$.
Theorem 38 (multiplication and transformation). Let $\mathbb{A} \in M(m \times k), \mathbb{B} \in M(k \times n), \mathbb{C} \in M(m \times n)$ and $A \mathbb{B}=\mathbb{C}$. Let $T$ be a transformation and $\mathbb{A} \stackrel{T}{\sim} \mathbb{A}^{\prime}$ and $\mathbb{C} \stackrel{T}{\sim} \mathbb{C}^{\prime}$. Then $\mathbb{A}^{\prime} \mathbb{B}=\mathbb{C}^{\prime}$.

Lemma 39. Let $\mathbb{A} \in M(n \times n)$ and $\operatorname{rank}(\mathbb{A})=n$. Then there exists a transformation transforming $\mathbb{A}$ to $\mathbb{I}$.
Theorem 40. Let $\mathbb{A} \in M(n \times n)$. Then $\mathbb{A}$ is invertible if and only if $\operatorname{rank}(\mathbb{A})=n$.

## VI.3. Determinants

Definition. Let $A \in M(n \times n)$. The symbol $A_{i j}$ denotes the $(n-1)$-by- $(n-1)$ matrix which is created from $A$ by omitting the $i$ th row and the $j$ th column.

Definition. Let $\mathbb{A}=\left(a_{i j}\right)_{i, j=1 . . n}$. The determinant of the matrix $A$ is defined by

$$
\operatorname{det} A= \begin{cases}a_{11} & \text { if } n=1 \\ \sum_{i=1}^{n}(-1)^{i+1} a_{i 1} \operatorname{det} A_{i 1} & \text { if } n>1\end{cases}
$$

For $\operatorname{det} A$ we will also use the symbol

$$
\left|\begin{array}{cccc}
a_{11} & a_{12} & \ldots & a_{1 n} \\
a_{21} & a_{22} & \ldots & a_{2 n} \\
\vdots & \ddots & \vdots & \\
a_{n 1} & a_{n 2} & \ldots & a_{n n}
\end{array}\right| .
$$

Theorem 41. Let $j, n \in \mathbb{N}, j \leq n$, and the matrices $A, \mathbb{B}, \mathbb{C} \in M(n \times n)$ coincide at each row except for the $j$ th row. Let the $j$ th row of $A$ be equal to the sum of the $j$ th rows of $\mathbb{B}$ and $\mathbb{C}$. Then $\operatorname{det} A=\operatorname{det} \mathbb{B}+\operatorname{det} \mathbb{C}$.

$$
\left|\begin{array}{ccc}
a_{11} & \ldots & a_{1 n} \\
\vdots & \ddots & \vdots \\
a_{j-1,1} & \ldots & a_{j-1, n} \\
u_{1}+v_{1} & \ldots & u_{n}+v_{n} \\
a_{j+1,1} & \ldots & a_{j+1, n} \\
\vdots & \ddots & \vdots \\
a_{n 1} & \ldots & a_{n n}
\end{array}\right|=\left|\begin{array}{ccc}
a_{11} & \ldots & a_{1 n} \\
\vdots & \ddots & \vdots \\
a_{j-1,1} & \ldots & a_{j-1, n} \\
u_{1} & \ldots & u_{n} \\
a_{j+1,1} & \ldots & a_{j+1, n} \\
\vdots & \ddots & \vdots \\
a_{n 1} & \ldots & a_{n n}
\end{array}\right|+\left|\begin{array}{ccc}
a_{11} & \ldots & a_{1 n} \\
\vdots & \ddots & \vdots \\
a_{j-1,1} & \ldots & a_{j-1, n} \\
v_{j+1,1} & \ldots & v_{n} \\
\vdots & \ddots & \vdots \\
a_{n 1} & \ldots & a_{n n}
\end{array}\right|
$$

Theorem 42 (determinant and transformations). Let $A, A^{\prime} \in M(n \times n)$.
(i) If the matrix $A^{\prime}$ is created from the matrix $\mathbb{A}$ by multiplying one row in $\mathbb{A}$ by a real number $\mu$, then $\operatorname{det} A^{\prime}=\mu \operatorname{det} A$.
(ii) If the matrix $A^{\prime}$ is created from $\mathbb{A}$ by interchanging two rows in $\mathbb{A}$ (i.e. by applying the elementary row operation of the first type), then $\operatorname{det} A^{\prime}=-\operatorname{det} A$.
(iii) If the matrix $A^{\prime}$ is created from $\mathbb{A}$ by adding a $\mu$-multiple of a row in $\mathbb{A}$ to another row in $\mathbb{A}$ (i.e. by applying the elementary row operation of the third type), then $\operatorname{det} A^{\prime}=\operatorname{det} A$.
(iv) If $A^{\prime}$ is created from $\mathbb{A}$ by applying a transformation, then $\operatorname{det} A \neq 0$ if and only if $\operatorname{det} A^{\prime} \neq 0$.

Remark. The determinant of a matrix with a zero row is equal to zero. The determinant of a matrix with two identical rows is also equal to zero.

Definition. Let $A=\left(a_{i j}\right)_{i, j=1 . . n}$. We say that $\mathbb{A}$ is an upper triangular matrix if $a_{i j}=0$ for $i>j, i, j \in\{1, \ldots, n\}$. We say that $A$ is a lower triangular matrix if $a_{i j}=0$ for $i<j, i, j \in\{1, \ldots, n\}$.

Theorem 43 (determinant of a triangular matrix). Let $A=\left(a_{i j}\right)_{i, j=1 . . n}$ be an upper or lower triangular matrix. Then

$$
\operatorname{det} \mathbb{A}=a_{11} \cdot a_{22} \cdots \cdots a_{n n}
$$

Theorem 44 (determinant and invertibility). Let $\mathbb{A} \in M(n \times n)$. Then $A$ is invertible if and only if $\operatorname{det} \mathbb{A} \neq 0$.
Theorem 45 (determinant of a product). Let $A, \mathbb{B} \in M(n \times n)$. Then $\operatorname{det} A \mathbb{B}=\operatorname{det} A \cdot \operatorname{det} \mathbb{B}$.
Theorem 46 (determinant of a transpose). Let $\mathbb{A} \in M(n \times n)$. Then $\operatorname{det} A^{T}=\operatorname{det} A$.
Theorem 47 (cofactor expansion). Let $\mathbb{A}=\left(a_{i j}\right)_{i, j=1 . . n}, k \in\{1, \ldots, n\}$. Then

$$
\begin{array}{ll}
\operatorname{det} A & =\sum_{i=1}^{n}(-1)^{i+k} a_{i k} \operatorname{det} A_{i k} \\
\text { (expansion along kth column), } \\
\operatorname{det} A & =\sum_{j=1}^{n}(-1)^{k+j} a_{k j} \operatorname{det} A_{k j} \\
\text { (expansion along kth row). }
\end{array}
$$

## VI.4. Systems of linear equations

A system of $m$ equations in $n$ unknowns $x_{1}, \ldots, x_{n}$ :

$$
\begin{gather*}
a_{11} x_{1}+a_{12} x_{2}+\cdots+a_{1 n} x_{n}=b_{1} \\
a_{21} x_{1}+a_{22} x_{2}+\cdots+a_{2 n} x_{n}=b_{2}  \tag{S}\\
\vdots \\
a_{m 1} x_{1}+a_{m 2} x_{2}+\cdots+a_{m n} x_{n}=b_{m}
\end{gather*}
$$

where $a_{i j} \in \mathbb{R}, b_{i} \in \mathbb{R}, i=1, \ldots, m, j=1, \ldots, n$. The matrix form is

$$
A \boldsymbol{x}=\boldsymbol{b},
$$

where $A=\left(\begin{array}{ccc}a_{11} & \ldots & a_{1 n} \\ \vdots & \ddots & \vdots \\ a_{m 1} & \ldots & a_{m n}\end{array}\right) \in M(m \times n)$, is called the coefficient matrix, $\boldsymbol{b}=\left(\begin{array}{c}b_{1} \\ \vdots \\ b_{m}\end{array}\right) \in M(m \times 1)$ is called the vector of the right-hand side and $\boldsymbol{x}=\left(\begin{array}{c}x_{1} \\ \vdots \\ x_{n}\end{array}\right) \in M(n \times 1)$ is the vector of unknowns.
Definition. The matrix

$$
(A \mid \boldsymbol{b})=\left(\begin{array}{ccc|c}
a_{11} & \ldots & a_{1 n} & b_{1} \\
\vdots & \ddots & \vdots & \vdots \\
a_{m 1} & \ldots & a_{m n} & b_{m}
\end{array}\right)
$$

is called the augmented matrix of the system (S).
Proposition 48. Let $\mathbb{A} \in M(m \times n), \boldsymbol{b} \in M(m \times 1)$ and let $T$ be a transformation of matrices with $m$ rows. Denote $\mathbb{A} \xrightarrow[\sim]{T} \mathbb{A}^{\prime}$, $\boldsymbol{b} \xrightarrow{T} \boldsymbol{b}^{\prime}$. Then for any $\boldsymbol{y} \in M(n \times 1)$ we have $A \boldsymbol{y}=\boldsymbol{b}$ if and only if $\mathbb{A}^{\prime} \boldsymbol{y}=\boldsymbol{b}^{\prime}$, i.e. the systems $A \boldsymbol{x}=\boldsymbol{b}$ and $A^{\prime} \boldsymbol{x}=\boldsymbol{b}^{\prime}$ have the same set of solutions.

Theorem 49 (Rouché-Fontené). The system (S) has a solution if and only if its coefficient matrix has the same rank as its augmented matrix.

Systems of $n$ equations in $n$ variables
Theorem 50. Let $\mathbb{A} \in M(n \times n)$. Then the following statements are equivalent:
(i) the matrix $A$ is invertible,
(ii) for each $\boldsymbol{b} \in M(n \times 1)$ the system ( S ) has a unique solution,
(iii) for each $\boldsymbol{b} \in M(n \times 1)$ the system ( S ) has at least one solution.

Theorem 51 (Cramer's rule). Let $\mathbb{A} \in M(n \times n)$ be an invertible matrix, $\boldsymbol{b} \in M(n \times 1)$, $\boldsymbol{x} \in M(n \times 1)$, and $\mathbb{A} \boldsymbol{x}=\boldsymbol{b}$. Then

$$
x_{j}=\frac{\left|\begin{array}{ccccccc}
a_{11} & \ldots & a_{1, j-1} & b_{1} & a_{1, j+1} & \ldots & a_{1 n} \\
\vdots & & & \vdots & & & \vdots \\
a_{n 1} & \ldots & a_{n, j-1} & b_{n} & a_{n, j+1} & \ldots & a_{n n}
\end{array}\right|}{\operatorname{det} A}
$$

for $j=1, \ldots, n$.

## VI.5. Matrices and linear mappings

Definition. We say that a mapping $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is linear if
(i) $\forall \boldsymbol{u}, \boldsymbol{v} \in \mathbb{R}^{n}: f(\boldsymbol{u}+\boldsymbol{v})=f(\boldsymbol{u})+f(\boldsymbol{v})$,
(ii) $\forall \lambda \in \mathbb{R} \forall \boldsymbol{u} \in \mathbb{R}^{n}: f(\lambda \boldsymbol{u})=\lambda f(\boldsymbol{u})$.

Definition. Let $i \in\{1, \ldots, n\}$. The vector with $n$ coordinates

$$
\boldsymbol{e}^{i}=\left(\begin{array}{c}
0 \\
\vdots \\
0 \\
1 \\
0 \\
\vdots \\
0
\end{array}\right) \ldots i \text { th coordinate }
$$

is called the $i$ th canonical basis vector of the space $\mathbb{R}^{n}$. The set $\left\{\boldsymbol{e}^{1}, \ldots, \boldsymbol{e}^{n}\right\}$ of all canonical basis vectors in $\mathbb{R}^{n}$ is called the canonical basis of the space $\mathbb{R}^{n}$.

Properties of the canonical basis:
(i) $\forall \boldsymbol{x} \in \mathbb{R}^{n}: \boldsymbol{x}=x_{1} \cdot \boldsymbol{e}^{1}+\cdots+x_{n} \cdot \boldsymbol{e}^{n}$,
(ii) the vectors $\boldsymbol{e}^{1}, \ldots, \boldsymbol{e}^{n}$ are linearly independent.

Theorem 52 (representation of linear mappings). The mapping $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ is linear if and only if there exists a matrix $A \in M(m \times n)$ such that

$$
\forall \boldsymbol{u} \in \mathbb{R}^{n}: f(\boldsymbol{u})=\mathbb{A} \boldsymbol{u}
$$

Remark. The matrix $\mathbb{A}$ from the previous theorem is uniquely determined and is called the representing matrix of the linear mapping $f$.

Theorem 53. Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be a linear mapping. Then the following statements are equivalent:
(i) $f$ is a bijection (i.e. $f$ is a one-to-one mapping of $\mathbb{R}^{n}$ onto $\mathbb{R}^{n}$ ),
(ii) $f$ is a one-to-one mapping,
(iii) $f$ is a mapping of $\mathbb{R}^{n}$ onto $\mathbb{R}^{n}$.

Theorem 54 (composition of linear mappings). Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be a linear mapping represented by a matrix $\mathbb{A} \in M(m \times n)$ and $g: \mathbb{R}^{m} \rightarrow \mathbb{R}^{k}$ a linear mapping represented by a matrix $\mathbb{B} \in M(k \times m)$. Then the composed mapping $g \circ f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{k}$ is linear and is represented by the matrix $\mathbb{B} A$.

## VII. Antiderivatives and Riemann integral

## VII.1. Antiderivatives

Definition. Let $f$ be a function defined on an open interval $I$. We say that a function $F: I \rightarrow \mathbb{R}$ is an antiderivative of $f$ on $I$ if for each $x \in I$ the derivative $F^{\prime}(x)$ exists and $F^{\prime}(x)=f(x)$.

Remark. An antiderivative of $f$ is sometimes called a function primitive to $f$.
If $F$ is an antiderivative of $f$ on $I$, then $F$ is continuous on $I$.
Theorem 55. Let $F$ and $G$ be antiderivatives of $f$ on an open interval $I$. Then there exists $c \in \mathbb{R}$ such that $F(x)=G(x)+c$ for each $x \in I$.

Remark. The set of all antiderivatives of $f$ on an open interval $I$ is denoted by

$$
\int f(x) \mathrm{d} x .
$$

The fact that $F$ is an antiderivative of $f$ on $I$ is expressed by

$$
\int f(x) \mathrm{d} x \stackrel{c}{=} F(x), \quad x \in I
$$

Table of basic antiderivatives

- $\int x^{n} \mathrm{~d} x \stackrel{c}{=} \frac{x^{n+1}}{n+1}$ on $\mathbb{R}$ for $n \in \mathbb{N} \cup\{0\}$; on $(-\infty, 0)$ and on $(0, \infty)$ for $n \in \mathbb{Z}, n<-1$,
- $\int x^{\alpha} \mathrm{d} x \stackrel{c}{=} \frac{x^{\alpha+1}}{\alpha+1}$ on $(0,+\infty)$ for $\alpha \in \mathbb{R} \backslash\{-1\}$,
- $\int \frac{1}{x} \mathrm{~d} x \stackrel{c}{=} \log |x|$ on $(0,+\infty)$ and on $(-\infty, 0)$,
- $\int e^{x} \mathrm{~d} x \stackrel{c}{=} e^{x}$ on $\mathbb{R}$,
- $\int \sin x \mathrm{~d} x \stackrel{c}{=}-\cos x$ on $\mathbb{R}$,
- $\int \cos x \mathrm{~d} x \stackrel{c}{=} \sin x$ on $\mathbb{R}$,
- $\int \frac{1}{\cos ^{2} x} \mathrm{~d} x \stackrel{c}{=} \operatorname{tg} x$ on each of the intervals $\left(-\frac{\pi}{2}+k \pi, \frac{\pi}{2}+k \pi\right), k \in \mathbb{Z}$,
- $\int \frac{1}{\sin ^{2} x} \mathrm{~d} x \stackrel{c}{=}-\operatorname{cotg} x$ on each of the intervals $(k \pi, \pi+k \pi), k \in \mathbb{Z}$,
- $\int \frac{1}{1+x^{2}} \mathrm{~d} x \stackrel{c}{=} \operatorname{arctg} x$ on $\mathbb{R}$,
- $\int \frac{1}{\sqrt{1-x^{2}}} \mathrm{~d} x \stackrel{c}{=} \arcsin x$ on $(-1,1)$,
- $\int-\frac{1}{\sqrt{1-x^{2}}} \mathrm{~d} x \stackrel{c}{=} \arccos x$ on $(-1,1)$.

Theorem 56. Let $f$ be a continuous function on an open interval $I$. Then $f$ has an antiderivative on $I$.
Theorem 57. Suppose that $f$ has an antiderivative $F$ on an open interval $I, g$ has an antiderivative $G$ on $I$, and let $\alpha, \beta \in \mathbb{R}$. Then the function $\alpha F+\beta G$ is an antiderivative of $\alpha f+\beta g$ on $I$.

Theorem 58 (substitution).
(i) Let $F$ be an antiderivative of $f$ on $(a, b)$. Let $\varphi:(\alpha, \beta) \rightarrow(a, b)$ have a finite derivative at each point of $(\alpha, \beta)$. Then

$$
\int f(\varphi(x)) \varphi^{\prime}(x) \mathrm{d} x \stackrel{c}{=} F(\varphi(x)) \quad \text { on }(\alpha, \beta) .
$$

(ii) Let $\varphi$ be a function with a finite derivative in each point of $(\alpha, \beta)$ such that the derivative is either everywhere positive or everywhere negative, and such that $\varphi((\alpha, \beta))=(a, b)$. Let $f$ be a function defined on $(a, b)$ and suppose that

$$
\int f(\varphi(t)) \varphi^{\prime}(t) \mathrm{d} t \stackrel{c}{=} G(t) \quad \text { on }(\alpha, \beta) .
$$

Then

$$
\int f(x) \mathrm{d} x \stackrel{c}{=} G\left(\varphi^{-1}(x)\right) \quad \text { on }(a, b) .
$$

Theorem 59 (integration by parts). Let I be an open interval and let the functions $f$ and $g$ be continuous on $I$. Let $F$ be an antiderivative of $f$ on $I$ and $G$ an antiderivative of $g$ on $I$. Then

$$
\int f(x) G(x) \mathrm{d} x=F(x) G(x)-\int F(x) g(x) \mathrm{d} x \quad \text { on } I .
$$

Example. Denote $I_{n}=\int \frac{1}{\left(1+x^{2}\right)^{n}} \mathrm{~d} x, n \in \mathbb{N}$. Then

$$
\begin{aligned}
I_{n+1} & =\frac{x}{2 n\left(1+x^{2}\right)^{n}}+\frac{2 n-1}{2 n} I_{n}, x \in \mathbb{R}, \quad n \in \mathbb{N}, \\
I_{1} & \stackrel{c}{=} \operatorname{arctg} x, x \in \mathbb{R}
\end{aligned}
$$

Definition. A rational function is a ratio of two polynomials, where the polynomial in the denominator is not a zero polynomial.
Theorem ("fundamental theorem of algebra"). Let $n \in \mathbb{N}, a_{0}, \ldots, a_{n} \in \mathbb{C}$, $a_{n} \neq 0$. Then the equation

$$
a_{n} z^{n}+a_{n-1} z^{n-1}+\cdots+a_{1} z+a_{0}=0
$$

has at least one solution $z \in \mathbb{C}$.
Lemma 60 (polynomial division). Let $P$ and $Q$ be polynomials (with complex coefficients) such that $Q$ is not a zero polynomial. Then there are uniquely determined polynomials $R$ and $Z$ satisfying:

- $\operatorname{deg} Z<\operatorname{deg} Q$,
- $P(x)=R(x) Q(x)+Z(x)$ for all $x \in \mathbb{C}$.

If $P$ and $Q$ have real coefficients then so have $R$ and $Z$.
Corollary. If $P$ is a polynomials and $\lambda \in \mathbb{C}$ its root (i.e. $P(\lambda)=0)$, then there is a polynomial $R$ satisfying $P(x)=(x-\lambda) R(x)$ for all $x \in \mathbb{C}$.

Theorem 61 (factorisation into monomials). Let $P(x)=a_{n} x^{n}+\cdots+a_{1} x+a_{0}$ be a polynomial of degree $n \in \mathbb{N}$. Then there are numbers $x_{1}, \ldots, x_{n} \in \mathbb{C}$ such that

$$
P(x)=a_{n}\left(x-x_{1}\right) \cdots\left(x-x_{n}\right), \quad x \in \mathbb{C} .
$$

Definition. Let $P$ be a polynomial that is not zero, $\lambda \in \mathbb{C}$, and $k \in \mathbb{N}$. We say that $\lambda$ is a root of multiplicity $k$ of the polynomial $P$ if there is a polynomial $R$ satisfying $R(\lambda) \neq 0$ and $P(x)=(x-\lambda)^{k} R(x)$ for all $x \in \mathbb{C}$.

Theorem 62 (roots of a polynomial with real coefficients). Let $P$ be a polynomial with real coefficients and $\lambda \in \mathbb{C}$ a root of $P$ of multiplicity $k \in \mathbb{N}$. Then the also the conjugate number $\bar{\lambda}$ is a root of $P$ of multiplicity $k$.

Theorem 63 (factorisation of a polynomial with real coefficients). Let $P(x)=a_{n} x^{n}+\cdots+a_{1} x+a_{0}$ be a polynomial of degree $n$ with real coefficients. Then there exist real numbers $x_{1}, \ldots, x_{k}, \alpha_{1}, \ldots, \alpha_{l}, \beta_{1}, \ldots, \beta_{l}$ and natural numbers $p_{1}, \ldots, p_{k}$, $q_{1}, \ldots, q_{l}$ such that

- $P(x)=a_{n}\left(x-x_{1}\right)^{p_{1}} \cdots\left(x-x_{k}\right)^{p_{k}}\left(x^{2}+\alpha_{1} x+\beta_{1}\right)^{q_{1}} \cdots\left(x^{2}+\alpha_{l} x+\beta_{l}\right)^{q_{l}}$,
- no two polynomials from $x-x_{1}, x-x_{2}, \ldots, x-x_{k}, x^{2}+\alpha_{1} x+\beta_{1}, \ldots, x^{2}+\alpha_{l} x+\beta_{l}$ have a common root,
- the polynomials $x^{2}+\alpha_{1} x+\beta_{1}, \ldots, x^{2}+\alpha_{l} x+\beta_{l}$ have no real root.

Theorem 64 (decomposition to partial fractions). Let $P, Q$ be polynomials with real coefficients such that $\operatorname{deg} P<\operatorname{deg} Q$ and let

$$
Q(x)=a_{n}\left(x-x_{1}\right)^{p_{1}} \cdots\left(x-x_{k}\right)^{p_{k}}\left(x^{2}+\alpha_{1} x+\beta_{1}\right)^{q_{1}} \cdots\left(x^{2}+\alpha_{l} x+\beta_{l}\right)^{q_{l}}
$$

be a factorisation of from Theorem 63. Then there exist unique real numbers $A_{1}^{1}, \ldots, A_{p_{1}}^{1}, \ldots, A_{1}^{k}, \ldots, A_{p_{k}}^{k}$, $B_{1}^{1}, C_{1}^{1}, \ldots, B_{q_{1}}^{1}, C_{q_{1}}^{1}, \ldots, B_{1}^{l}, C_{1}^{l}, \ldots, B_{q_{l}}^{l}, C_{q_{l}}^{l}$ such that

$$
\begin{aligned}
\frac{P(x)}{Q(x)}= & \frac{A_{1}^{1}}{\left(x-x_{1}\right)}+\cdots+\frac{A_{p_{1}}^{1}}{\left(x-x_{1}\right)^{p_{1}}}+\cdots+\frac{A_{1}^{k}}{\left(x-x_{k}\right)}+\cdots+\frac{A_{p_{k}}^{k}}{\left(x-x_{k}\right)^{p_{k}}}+ \\
& \quad+\frac{B_{1}^{1} x+C_{1}^{1}}{\left(x^{2}+\alpha_{1} x+\beta_{1}\right)}+\cdots+\frac{B_{q_{1}}^{1} x+C_{q_{1}}^{1}}{\left(x^{2}+\alpha_{1} x+\beta_{1}\right)^{q_{1}}}+\cdots+ \\
& +\frac{B_{1}^{l} x+C_{1}^{l}}{\left(x^{2}+\alpha_{l} x+\beta_{l}\right)}+\cdots+\frac{B_{q_{l}}^{l} x+C_{q_{l}}^{l}}{\left(x^{2}+\alpha_{l} x+\beta_{l}\right)^{q_{l}}}, x \in \mathbb{R} \backslash\left\{x_{1}, \ldots, x_{k}\right\} .
\end{aligned}
$$

## VII.2. Riemann integral

Definition. A finite sequence $\left\{x_{j}\right\}_{j=0}^{n}$ is called a partition of the interval $[a, b]$ if

$$
a=x_{0}<x_{1}<\cdots<x_{n}=b
$$

The points $x_{0}, \ldots, x_{n}$ are called the partition points.
We say that a partition $D^{\prime}$ of an interval $[a, b]$ is a refinement of the partition $D$ of $[a, b]$ if each partition point of $D$ is also a partition point of $D^{\prime}$.

Definition. Suppose that $a, b \in \mathbb{R}, a<b$, the function $f$ is bounded on $[a, b]$, and $D=\left\{x_{j}\right\}_{j=0}^{n}$ is a partition of $[a, b]$. Denote

$$
\begin{aligned}
\bar{S}(f, D)= & \sum_{j=1}^{n} M_{j}\left(x_{j}-x_{j-1}\right), \text { where } M_{j}=\sup \left\{f(x) ; x \in\left[x_{j-1}, x_{j}\right]\right\}, \\
\underline{S}(f, D)= & \sum_{j=1}^{n} m_{j}\left(x_{j}-x_{j-1}\right), \text { where } m_{j}=\inf \left\{f(x) ; x \in\left[x_{j-1}, x_{j}\right]\right\}, \\
& \int_{a}^{b} f=\inf \{\bar{S}(f, D) ; D \text { is a partition of }[a, b]\}, \\
& \int_{a}^{b} f=\sup \{\underline{S}(f, D) ; D \text { is a partition of }[a, b]\} .
\end{aligned}
$$

Definition. We say that a function $f$ has the Riemann integral over the interval $[a, b]$ if $\overline{\int_{a}^{b}} f=\underline{\int_{a}^{b}} f$. The value of the integral of $f$ over $[a, b]$ is then equal to the common value of $\overline{\int_{a}^{b}} f=\underline{\int_{a}^{b}} f$. We denote it by $\int_{a}^{b} f$. If $a>b$, then we define $\int_{a}^{b} f=-\int_{b}^{a} f$, and in case that $a=b$ we put $\int_{a}^{b} f=0$.
Remark. Let $D, D^{\prime}$ be partitions of $[a, b], D^{\prime}$ refines $D$, and let $f$ be a bounded function on $[a, b]$. Then

$$
\underline{S}(f, D) \leq \underline{S}\left(f, D^{\prime}\right) \leq \bar{S}\left(f, D^{\prime}\right) \leq \bar{S}(f, D) .
$$

Suppose that $D_{1}, D_{2}$ are partitions of $[a, b]$ and a partition $D^{\prime}$ refines both $D_{1}$ and $D_{2}$. Then

$$
\underline{S}\left(f, D_{1}\right) \leq \underline{S}\left(f, D^{\prime}\right) \leq \bar{S}\left(f, D^{\prime}\right) \leq \bar{S}\left(f, D_{2}\right) .
$$

It easily follows that $\underline{\int_{a}^{b}} f \leq \overline{\int_{a}^{b}} f$.
Lemma 65 (criterion for the existence of the Riemann integral). Let $f$ be a function bounded on an interval $[a, b]$.
(a) $\int_{a}^{b} f=I \in \mathbb{R}$ if and only iffor each $\varepsilon \in \mathbb{R}, \varepsilon>0$ there exists a partition $D$ of $[a, b]$ such that

$$
I-\varepsilon<\underline{S}(f, D) \leq \bar{S}(f, D)<I+\varepsilon
$$

(b) $f$ has the Riemann integral over $[a, b]$ if and only iffor each $\varepsilon \in \mathbb{R}, \varepsilon>0$ there exists a partition $D$ of $[a, b]$ such that

$$
\bar{S}(f, D)-\underline{S}(f, D)<\varepsilon .
$$

Theorem 66. (i) Suppose that $f$ has the Riemann integral over $[a, b]$ and let $[c, d] \subset[a, b]$. Then $f$ has the Riemann integral also over $[c, d]$.
(ii) Suppose that $c \in(a, b)$ and $f$ has the Riemann integral over the intervals $[a, c]$ and $[c, b]$. Then $f$ has the Riemann integral over $[a, b]$ and

$$
\begin{equation*}
\int_{a}^{b} f=\int_{a}^{c} f+\int_{c}^{b} f \tag{1}
\end{equation*}
$$

Remark. The formula (1) holds for all $a, b, c \in \mathbb{R}$ if the integral of $f$ exists over the interval $[\min \{a, b, c\}, \max \{a, b, c\}]$.
Theorem 67 (linearity of the Riemann integral). Let $f$ and $g$ be functions with Riemann integral over $[a, b]$ and let $\alpha \in \mathbb{R}$. Then
(i) the function $\alpha f$ has the Riemann integral over $[a, b]$ and

$$
\int_{a}^{b} \alpha f=\alpha \int_{a}^{b} f
$$

(ii) the function $f+g$ has the Riemann integral over $[a, b]$ and

$$
\int_{a}^{b} f+g=\int_{a}^{b} f+\int_{a}^{b} g
$$

Theorem 68. Let $a, b \in \mathbb{R}, a<b$, and let $f$ and $g$ be functions with Riemann integral over $[a, b]$. Then:
(i) If $f(x) \leq g(x)$ for each $x \in[a, b]$, then

$$
\int_{a}^{b} f \leq \int_{a}^{b} g
$$

(ii) The function $|f|$ has the Riemann integral over $[a, b]$ and

$$
\left|\int_{a}^{b} f\right| \leq \int_{a}^{b}|f| .
$$

Definition. We say that a function $f$ is uniformly continuous on an interval $I$ if

$$
\forall \varepsilon \in \mathbb{R}, \varepsilon>0 \exists \delta \in \mathbb{R}, \delta>0 \forall x, y \in I,|x-y|<\delta:|f(x)-f(y)|<\varepsilon .
$$

Theorem 69. If $f$ is continuous on a closed bounded interval $[a, b]$, then it is uniformly continuous on $[a, b]$.
Theorem 70. Let $f$ be a function continuous on an interval $[a, b], a, b \in \mathbb{R}$. Then $f$ has the Riemann integral on $[a, b]$.
Theorem 71. Let $f$ be a function continuous on an interval $(a, b)$ and let $c \in(a, b)$. If we denote $F(x)=\int_{c}^{x} f(t) \mathrm{d} t$ for $x \in(a, b)$, then $F^{\prime}(x)=f(x)$ for each $x \in(a, b)$. In other words, $F$ is an antiderivative of $f$ on $(a, b)$.

Theorem 72 (Newton-Leibniz formula). Let $f$ be a function continuous on an interval $[a, b], a, b \in \mathbb{R}, a<b$, and let $F$ be an antiderivative of $f$ on $(a, b)$. Then the limits $\lim _{x \rightarrow a+} F(x), \lim _{x \rightarrow b-} F(x)$ exist, are finite, and

$$
\int_{a}^{b} f(x) \mathrm{d} x=\lim _{x \rightarrow b-} F(x)-\lim _{x \rightarrow a+} F(x)
$$

Remark. Let us denote

$$
[F]_{a}^{b}= \begin{cases}\lim _{x \rightarrow b-} F(x)-\lim _{x \rightarrow a+} F(x) & \text { for } a<b \\ \lim _{x \rightarrow b+} F(x)-\lim _{x \rightarrow a-} F(x) & \text { for } b<a\end{cases}
$$

Then the Newton-Leibniz formula can be written as

$$
\int_{a}^{b} f=[F]_{a}^{b}
$$

even for $b<a$.
Theorem 73 (integration by parts). Suppose that the functions $f, g, f^{\prime}$ a $g^{\prime}$ are continuous on an interval $[a, b]$. Then

$$
\int_{a}^{b} f^{\prime} g=[f g]_{a}^{b}-\int_{a}^{b} f g^{\prime}
$$

Theorem 74 (substitution). Let the function $f$ be continuous on an interval $[a, b]$. Suppose that the function $\varphi$ has a continuous derivative on $[\alpha, \beta]$ and $\varphi$ maps $[\alpha, \beta]$ into the interval $[a, b]$. Then

$$
\int_{\alpha}^{\beta} f(\varphi(x)) \varphi^{\prime}(x) \mathrm{d} x=\int_{\varphi(\alpha)}^{\varphi(\beta)} f(t) \mathrm{d} t .
$$

Theorem (logarithm). There exist a unique function $\log$ with the following properties:
(L1) $D_{\log }=(0,+\infty)$,
(L2) the function $\log$ is increasing on $(0,+\infty)$,
(L3) $\forall x, y \in(0,+\infty): \log x y=\log x+\log y$,
(L4) $\lim _{x \rightarrow 1} \frac{\log x}{x-1}=1$.

