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Elementary Introduction to Stochastic Finance in Discrete Time

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Summary. This article gives an elementary introduction to stochastic finance (in discrete time). A formalization of random variables is given and some elements of Borel sets are considered. Furthermore, special functions (for buying a present portfolio and the value of a portfolio in the future) and some statements about the relation between these functions are introduced. For details see: [8] (p. 185), [7] (pp. 12, 20), [6] (pp. 3–6).

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The notation and terminology used in this paper have been introduced in the following papers: [15], [2], [1], [3], [4], [11], [10], [9], [5], [14], [12], and [13].

We use the following convention: $O_1, O_2$ are non empty sets, $S_1, F$ are $\sigma$-fields of subsets of $O_1$, and $S_2, F_2$ are $\sigma$-fields of subsets of $O_2$.

Let $a, r$ be real numbers. We introduce the halfline finance of $a$ and $r$ as a synonym of $[a, r[. Then the halfline finance of $a$ and $r$ is a subset of $\mathbb{R}$.

We now state two propositions:

1) For every real number $k$ holds $\mathbb{R} \setminus [k, +\infty[ = ]-\infty, k[.$

2) For every real number $k$ holds $\mathbb{R} \setminus ]-\infty, k[ = [k, +\infty[.$

Let $a, b$ be real numbers. The half open sets of $a$ and $b$ yields a sequence of subsets of $\mathbb{R}$ and is defined by the conditions (Def. 1).

(Def. 1)(i) (The half open sets of $a$ and $b$)(0) = the halfline finance of $a$ and $b + 1$, and

(ii) for every element $n$ of $\mathbb{N}$ holds (the half open sets of $a$ and $b$)(n+1) = the halfline finance of $a$ and $b + \frac{1}{n+1}$.

A sequence of real numbers is said to be a price function if:
(Def. 2) It(0) = 1 and for every element \( n \) of \( \mathbb{N} \) holds \( \text{it}(n) \geq 0 \).

Let \( p_1, j_1 \) be sequences of real numbers. We introduce the elements of buy portfolio of \( p_1 \) and \( j_1 \) as a synonym of \( p_1 \cdot j_1 \). Then the elements of buy portfolio of \( p_1 \) and \( j_1 \) is a sequence of real numbers.

Let \( d \) be a natural number. The buy portfolio extension of \( p_1, j_1 \), and \( d \) yields an element of \( \mathbb{R} \) and is defined as follows:

(Def. 3) The buy portfolio extension of \( p_1, j_1 \), and \( d = (\sum_{\alpha=0}^{\kappa} (\text{the elements of buy portfolio of } p_1 \text{ and } j_1)(\alpha))_{\kappa \in \mathbb{N}}(d) \).

The buy portfolio of \( p_1, j_1 \), and \( d \) yielding an element of \( \mathbb{R} \) is defined as follows:

(Def. 4) The buy portfolio of \( p_1, j_1 \), and \( d = (\sum_{\alpha=0}^{\kappa} (\text{the elements of buy portfolio of } p_1 \text{ and } j_1)^{\uparrow} 1)(\alpha))_{\kappa \in \mathbb{N}}(d - 1) \).

Let \( O_1, O_2 \) be sets, let \( S_1 \) be a \( \sigma \)-field of subsets of \( O_1 \), let \( S_2 \) be a \( \sigma \)-field of subsets of \( O_2 \), and let \( X \) be a function. We say that \( X \) is random variable on \( S_1 \) and \( S_2 \) if and only if:

(Def. 5) For every element \( x \) of \( S_2 \) holds \( \{ y \in O_1 : X(y) \text{ is an element of } x \} \) is an element of \( S_1 \).

Let \( O_1, O_2 \) be sets, let \( F \) be a \( \sigma \)-field of subsets of \( O_1 \), and let \( F_2 \) be a \( \sigma \)-field of subsets of \( O_2 \). The set of random variables on \( F \) and \( F_2 \) is defined by:

(Def. 6) The set of random variables on \( F \) and \( F_2 = \{ f : O_1 \rightarrow O_2 : f \text{ is random variable on } F \text{ and } F_2 \} \).

Let us consider \( O_1, O_2, F, F_2 \). One can check that the set of random variables on \( F \) and \( F_2 \) is non empty.

Let \( O_1, O_2 \) be non empty sets, let \( F \) be a \( \sigma \)-field of subsets of \( O_1 \), let \( F_2 \) be a \( \sigma \)-field of subsets of \( O_2 \), and let \( X \) be a set. Let us assume that \( X = \) the set of random variables on \( F \) and \( F_2 \). Let \( k \) be an element of \( X \). The change element to function \( F, F_2, \) and \( k \) yielding a function from \( O_1 \) into \( O_2 \) is defined by:

(Def. 7) The change element to function \( F, F_2, \) and \( k = k \).

Let \( O_1 \) be a non empty set, let \( F \) be a \( \sigma \)-field of subsets of \( O_1 \), let \( X \) be a non empty set, and let \( k \) be an element of \( X \). The random variables for future elements of portfolio value of \( F \) and \( k \) yields a function from \( O_1 \) into \( \mathbb{R} \) and is defined by the condition (Def. 8).

(Def. 8) Let \( w \) be an element of \( O_1 \). Then (the random variables for future elements of portfolio value of \( F \) and \( k \)) \((w) = (\text{the change element to function } F, \text{ the Borel sets, and } k)(w) \).

Let \( p \) be a natural number, let \( O_1, O_2 \) be non empty sets, let \( F \) be a \( \sigma \)-field of subsets of \( O_1 \), let \( F_2 \) be a \( \sigma \)-field of subsets of \( O_2 \), and let \( X \) be a set. Let us assume that \( X = \) the set of random variables on \( F \) and \( F_2 \). Let \( G \) be a function from \( \mathbb{N} \) into \( X \). The element of \( F, F_2, G, \) and \( p \) yields a function from \( O_1 \) into \( O_2 \) and is defined as follows:

(Def. 9) The element of \( F, F_2, G, \) and \( p = G(p) \).
Let $r$ be a real number, let $O_1$ be a non empty set, let $F$ be a $\sigma$-field of subsets of $O_1$, let $X$ be a non empty set, let $w$ be an element of $O_1$, let $G$ be a function from $\mathbb{N}$ into $X$, and let $p_1$ be a sequence of real numbers. The future elements of portfolio value of $r$, $p_1$, $F$, $w$, and $G$ yields a sequence of real numbers and is defined by the condition (Def. 10).

(Def. 10) Let $n$ be an element of $\mathbb{N}$. Then (the future elements of portfolio value of $r$, $p_1$, $F$, $w$, and $G$)\( (n) = (\text{the random variables for future elements of portfolio value of } r, p_1, F, w, \text{ and } G)(n) = (w) \cdot p_1(n). \)

Let $r$ be a real number, let $d$ be a natural number, let $p_1$ be a sequence of real numbers, let $O_1$ be a non empty set, let $F$ be a $\sigma$-field of subsets of $O_1$, let $X$ be a non empty set, let $G$ be a function from $\mathbb{N}$ into $X$, and let $w$ be an element of $O_1$. The future portfolio value extension of $r$, $d$, $p_1$, $F$, $G$, and $w$ yields an element of $\mathbb{R}$ and is defined by the condition (Def. 11).

(Def. 11) The future portfolio value extension of $r$, $d$, $p_1$, $F$, $G$, and $w = (\sum_{\alpha=0}^{\kappa} (\text{the future elements of portfolio value of } r, p_1, F, w, \text{ and } G)(\alpha))_{\kappa \in \mathbb{N}}(d)$.

The future portfolio value of $r$, $d$, $p_1$, $F$, $G$, and $w$ yields an element of $\mathbb{R}$ and is defined by the condition (Def. 12).

(Def. 12) The future portfolio value of $r$, $d$, $p_1$, $F$, $G$, and $w = (\sum_{\alpha=0}^{\kappa} ((\text{the future elements of portfolio value of } r, p_1, F, w, \text{ and } G) \uparrow 1)(\alpha))_{\kappa \in \mathbb{N}}(d - 1)$.

Let us observe that there exists an element of the Borel sets which is non empty.

One can prove the following propositions:

(3) For every real number $k$ holds $[k, +\infty]$ is an element of the Borel sets and $]-\infty, k]$ is an element of the Borel sets.

(4) For all real numbers $k_1$, $k_2$ holds $[k_2, k_1]$ is an element of the Borel sets.

(5) For all real numbers $a$, $b$ holds Intersection (the half open sets of $a$ and $b$) is an element of the Borel sets.

(6) For all real numbers $a$, $b$ holds Intersection (the half open sets of $a$ and $b$) = $[a, b]$.

(7) Let $a$, $b$ be real numbers and $n$ be a natural number. Then (the partial intersections of the half open sets of $a$ and $b$)(n) is an element of the Borel sets.

(8) For all real numbers $k_1$, $k_2$ holds $[k_2, k_1]$ is an element of the Borel sets.

(9) Let $X$ be a function from $O_1$ into $\mathbb{R}$. Suppose $X$ is random variable on $S_1$ and the Borel sets. Then for every real number $k$ holds $\{w \in O_1: X(w) \geq k\}$ is an element of $S_1$ and $\{w \in O_1: X(w) < k\}$ is an element of $S_1$ and for all real numbers $k_1$, $k_2$ such that $k_1 < k_2$ holds $\{w \in O_1: k_1 \leq X(w) \land X(w) < k_2\}$ is an element of $S_1$ and for all real numbers $k_1$, $k_2$ such that $k_1 \leq k_2$ holds $\{w \in O_1: k_1 \leq X(w) \land X(w) \leq k_2\}$ is an element of $S_1$. 


element of $S_1$ and for every real number $r$ holds $\text{LE-dom}(X, r) = \{w \in O_1: X(w) < r\}$ and for every real number $r$ holds $\text{GET-dom}(X, r) = \{w \in O_1: X(w) \geq r\}$ and for every real number $r$ holds $\text{EQ-dom}(X, r) = \{w \in O_1: X(w) = r\}$ and for every real number $r$ holds $\text{EQ-dom}(X, r)$ is an element of $S_1$.

(10) For every real number $s$ holds $O_1 \xrightarrow{s} s$ is random variable on $S_1$ and the Borel sets.

(11) Let $p_1$ be a sequence of real numbers, $j_1$ be a price function, and $d$ be a natural number. Suppose $d > 0.$ Then the buy portfolio extension of $p_1$, $j_1$, and $d = p_1(0) +$ the buy portfolio of $p_1$, $j_1$, and $d$.

(12) Let $d$ be a natural number. Suppose $d > 0.$ Let $r$ be a real number, $p_1$ be a sequence of real numbers, and $G$ be a function from $\mathbb{N}$ into the set of random variables on $F$ and the Borel sets. Suppose the element of $F$, the Borel sets, $G$, and $0 = O_1 \xrightarrow{r} 1 + r.$ Let $w$ be an element of $O_1.$ Then the future portfolio value extension of $r$, $d$, $p_1$, $F$, $G$, and $w = (1+r) \cdot p_1(0) +$ the future portfolio value of $r$, $d$, $p_1$, $F$, $G$, and $w$.

(13) Let $d$ be a natural number. Suppose $d > 0.$ Let $r$ be a real number. Suppose $r > -1.$ Let $p_1$ be a sequence of real numbers, $j_1$ be a price function, and $G$ be a function from $\mathbb{N}$ into the set of random variables on $F$ and the Borel sets. Suppose the element of $F$, the Borel sets, $G$, and $0 = O_1 \xrightarrow{r} 1 + r.$ Let $w$ be an element of $O_1.$ Suppose the buy portfolio extension of $p_1$, $j_1$, and $d \leq 0.$ Then the future portfolio value extension of $r$, $d$, $p_1$, $F$, $G$, and $w \leq$ (the future portfolio value of $r$, $d$, $p_1$, $F$, $G$, and $w$) $- (1+r) \cdot$ the buy portfolio of $p_1$, $j_1$, and $d$.

(14) Let $d$ be a natural number. Suppose $d > 0.$ Let $r$ be a real number. Suppose $r > -1.$ Let $p_1$ be a sequence of real numbers, $j_1$ be a price function, and $G$ be a function from $\mathbb{N}$ into the set of random variables on $F$ and the Borel sets. Suppose the element of $F$, the Borel sets, $G$, and $0 = O_1 \xrightarrow{r} 1 + r.$ Suppose the buy portfolio extension of $p_1$, $j_1$, and $d \leq 0.$ Then

(i) $\{w \in O_1: \text{the future portfolio value extension of } r, d, p_1, F, G, \text{ and } w \geq 0\} \subseteq \{w \in O_1: \text{the future portfolio value of } r, d, p_1, F, G, \text{ and } w \geq (1+r) \cdot \text{the buy portfolio of } p_1, j_1, \text{ and } d\},$ and

(ii) $\{w \in O_1: \text{the future portfolio value extension of } r, d, p_1, F, G, \text{ and } w > 0\} \subseteq \{w \in O_1: \text{the future portfolio value of } r, d, p_1, F, G, \text{ and } w > (1+r) \cdot \text{the buy portfolio of } p_1, j_1, \text{ and } d\}.$

(15) Let $f$ be a function from $O_1$ into $\mathbb{R}.$ Suppose $f$ is random variable on $S_1$ and the Borel sets. Then $f$ is measurable on $\Omega(S_1)$ and $f$ is a real-valued random variable on $S_1$.

(16) The set of random variables on $S_1$ and the Borel sets $\subseteq$ the real-valued random variables set on $S_1$. 


References


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Valuation Theory. Part I

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Summary. In the article we introduce a valuation function over a field [1]. Ring of non negative elements and its ideal of positive elements have been also defined.

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The notation and terminology used here have been introduced in the following papers: [11], [19], [4], [15], [20], [8], [21], [10], [9], [16], [3], [5], [7], [18], [17], [13], [14], [6], [2], and [12].

1. Extended Reals

We use the following convention: $x, y, z, s$ are extended real numbers, $i$ is an integer, and $n, m$ are natural numbers.

The following propositions are true:

(1) If $x = -x$, then $x = 0$.

(2) If $x + x = 0$, then $x = 0$. 
(3) If $0 \leq x \leq y$ and $0 \leq s \leq z$, then $x \cdot s \leq y \cdot z$.
(4) If $y \neq +\infty$ and $0 < x$ and $0 < y$, then $0 < \frac{x}{y}$.
(5) If $y \neq +\infty$ and $x < 0 < y$, then $\frac{x}{y} < 0$.
(6) If $y \neq -\infty$ and $0 < x$ and $y < 0$, then $\frac{x}{y} < 0$.
(7) If $x, y \in \mathbb{R}$ or $z \in \mathbb{R}$, then $\frac{x + y}{z} = \frac{x}{z} + \frac{y}{z}$.
(8) If $y \neq +\infty$ and $y \neq -\infty$ and $y \neq 0$, then $\frac{x}{y} \cdot y = x$.
(9) If $y \neq -\infty$ and $y \neq +\infty$ and $x \neq 0$ and $y \neq 0$, then $\frac{x}{y} \neq 0$.

Let $x$ be a number. We say that $x$ is extended integer if and only if:

(Def. 1) $x$ is integer or $x = +\infty$.

Let us mention that every number which is extended integer is also extended real.

One can verify the following observations:

* $+\infty$ is extended integer,
* $-\infty$ is non extended integer,
* $1$ is extended integer, positive, and real,
* every number which is integer is also extended integer, and
* every number which is real and extended integer is also integer.

Let us observe that there exists an element of $\mathbb{R}$ which is real, extended integer, and positive and there exists an extended integer number which is positive.

An extended integer is an extended integer number.

In the sequel $x, y, v$ denote extended integers.

One can prove the following propositions:

(10) If $x < y$, then $x + 1 \leq y$.
(11) $-\infty < x$.

Let $X$ be an extended real-membered set. Let us assume that there exists a positive extended integer $i_0$ such that $i_0 \in X$. The functor least-positive $X$ yielding a positive extended integer is defined by:

(Def. 2) least-positive $X \in X$ and for every positive extended integer $i$ such that $i \in X$ holds least-positive $X \leq i$.

Let $f$ be a binary relation. We say that $f$ is extended integer valued if and only if:

(Def. 3) For every set $x$ such that $x \in \text{rng} f$ holds $x$ is extended integer.

Let us note that there exists a function which is extended integer valued.

Let $A$ be a set. Note that there exists a function from $A$ into $\mathbb{R}$ which is extended integer valued.

Let $f$ be an extended integer valued function and let $x$ be a set. Note that $f(x)$ is extended integer.
2. Structures

One can prove the following proposition

(12) Let $K$ be a distributive left unital add-associative right zeroed right complementable non empty double loop structure. Then $-1_K \cdot -1_K = 1_K$.

Let $K$ be a non empty double loop structure, let $S$ be a subset of $K$, and let $n$ be a natural number. The functor $S^n$ yielding a subset of $K$ is defined by:

(Def. 4)(i) $S^n$ = the carrier of $K$ if $n = 0$,
(ii) there exists a finite sequence $f$ of elements of $2$the carrier of $K$ such that $S^n = f(\text{len } f)$ and $f(1) = S$ and for every natural number $i$ such that $i, i+1 \in \text{dom } f$ holds $f(i+1) = S \cdot f_i$, otherwise.

In the sequel $A$ denotes a subset of $D$. The following propositions are true:

(13) $A^1 = A$.
(14) $A^2 = A \cdot A$.

Let $R$ be a ring, let $S$ be an ideal of $R$, and let $n$ be a natural number. Observe that $S^n$ is non empty, add closed, left ideal, and right ideal.

Let $G$ be a non empty double loop structure, let $g$ be an element of $G$, and let $i$ be an integer. The functor $g^i$ yielding an element of $G$ is defined as follows:

(Def. 5) $g^i = \left\{ \begin{array}{ll} \text{power}_G(g, |i|), & \text{if } 0 \leq i, \\ \text{power}_G(g, |i|)^{-1}, & \text{otherwise.} \end{array} \right.$

Let $G$ be a non empty double loop structure, let $g$ be an element of $G$, and let $n$ be a natural number. Then $g^n$ can be characterized by the condition:

(Def. 6) $g^n = \text{power}_G(g, n)$.

In the sequel $K$ is a field-like non degenerated associative add-associative right zeroed right complementable distributive Abelian non empty double loop structure and $a, b, c$ are elements of $K$. We now state two propositions:

(15) $a^{n+m} = a^n \cdot a^m$.
(16) If $a \neq 0_K$, then $a^i \neq 0_K$.

3. Valuation

Let $K$ be a double loop structure. We say that $K$ has a valuation if and only if the condition (Def. 7) is satisfied.

(Def. 7) There exists an extended integer valued function $f$ from $K$ into $\mathbb{R}$ such that

(i) $f(0_K) = +\infty$,
(ii) for every element $a$ of $K$ such that $a \neq 0_K$ holds $f(a) \in \mathbb{Z}$,
(iii) for all elements $a, b$ of $K$ holds $f(a \cdot b) = f(a) + f(b)$,
(iv) for every element $a$ of $K$ such that $0 \leq f(a)$ holds $0 \leq f(1_K + a)$, and
(v) there exists an element $a$ of $K$ such that $f(a) \neq 0$ and $f(a) \neq +\infty$. 

Let $K$ be a double loop structure. Let us assume that $K$ has a valuation. An extended integer valued function from $K$ into $\mathbb{N}$ is said to be a valuation of $K$ if it satisfies the conditions (Def. 8).

(Def. 8) (i) $\text{It}(0_K) = +\infty$,
(ii) for every element $a$ of $K$ such that $a \neq 0_K$ holds $\text{it}(a) \in \mathbb{Z}$,
(iii) for all elements $a$, $b$ of $K$ holds $\text{it}(a \cdot b) = \text{it}(a) + \text{it}(b)$,
(iv) for every element $a$ of $K$ such that $0 \leq \text{it}(a)$ holds $0 \leq \text{it}(1_K + a)$, and
(v) there exists an element $a$ of $K$ such that $\text{it}(a) \neq 0$ and $\text{it}(a) \neq +\infty$.

In the sequel $v$ denotes a valuation of $K$.

One can prove the following propositions:

(17) If $K$ has a valuation, then $v(1_K) = 0$.
(18) If $K$ has a valuation and $a \neq 0_K$, then $v(a) \neq +\infty$.
(19) If $K$ has a valuation, then $v(-1_K) = 0$.
(20) If $K$ has a valuation, then $v(-a) = v(a)$.
(21) If $K$ has a valuation and $a \neq 0_K$, then $v(a^{-1}) = -v(a)$.
(22) If $K$ has a valuation and $b \neq 0_K$, then $v(\frac{a}{b}) = v(a) - v(b)$.
(23) If $K$ has a valuation and $a \neq 0_K$ and $b \neq 0_K$, then $v(\frac{a}{b}) = -v(\frac{b}{a})$.
(24) If $K$ has a valuation and $b \neq 0_K$ and $0 \leq v(\frac{a}{b})$, then $v(b) \leq v(a)$.
(25) If $K$ has a valuation and $a \neq 0_K$ and $b \neq 0_K$ and $v(\frac{a}{b}) \leq 0$, then $0 \leq v(\frac{b}{a})$.
(26) If $K$ has a valuation and $b \neq 0_K$ and $v(\frac{a}{b}) \leq 0$, then $v(a) \leq v(b)$.
(27) If $K$ has a valuation, then $\min(v(a), v(b)) \leq v(a + b)$.
(28) If $K$ has a valuation and $v(a) < v(b)$, then $v(a) = v(a + b)$.
(29) If $K$ has a valuation and $a \neq 0_K$, then $v(a^i) = i \cdot v(a)$.
(30) If $K$ has a valuation and $0 \leq v(1_K + a)$, then $0 \leq v(a)$.
(31) If $K$ has a valuation and $0 \leq v(1_K - a)$, then $0 \leq v(a)$.
(32) If $K$ has a valuation and $a \neq 0_K$ and $v(a) \leq v(b)$, then $0 \leq v(\frac{b}{a})$.
(33) If $K$ has a valuation, then $+\infty \in \text{rng } v$.
(34) If $v(a) = 1$, then least-positive $\text{rng } v = 1$.
(35) If $K$ has a valuation, then least-positive $\text{rng } v$ is integer.
(36) If $K$ has a valuation, then for every element $x$ of $K$ such that $x \neq 0_K$ there exists an integer $i$ such that $v(x) = i \cdot \text{least-positive } \text{rng } v$.

Let us consider $K$, $v$. Let us assume that $K$ has a valuation. The functor $\text{Pgenerator } v$ yielding an element of $K$ is defined as follows:

(Def. 9) $\text{Pgenerator } v =$ the element of $v^{-1}(\{\text{least-positive } \text{rng } v\})$.

Let us consider $K$, $v$. Let us assume that $K$ has a valuation. The functor $\text{normal-valuation } v$ yields a valuation of $K$ and is defined by:

(Def. 10) $v(a) = (\text{normal-valuation } v)(a) \cdot \text{least-positive } \text{rng } v$. 
We now state a number of propositions:

(37) If $K$ has a valuation, then $v(a) = 0$ iff (normal-valuation $v)(a) = 0$.
(38) If $K$ has a valuation, then $v(a) = +\infty$ iff (normal-valuation $v)(a) = +\infty$.
(39) If $K$ has a valuation, then $v(a) = v(b)$ iff (normal-valuation $v)(a) = (normal-valuation v)(b)$.
(40) If $K$ has a valuation, then $v(a)$ is positive iff (normal-valuation $v)(a)$ is positive.
(41) If $K$ has a valuation, then $0 \leq v(a)$ iff $0 \leq (normal-valuation v)(a)$.
(42) If $K$ has a valuation, then $v(a)$ is non negative iff (normal-valuation $v)(a)$ is non negative.
(43) If $K$ has a valuation, then $(normal-valuation v)(P_{generator v}) = 1$.
(44) If $K$ has a valuation and $0 \leq v(a)$, then $(normal-valuation v)(a) \leq v(a)$.
(45) If $K$ has a valuation and $v(a) = 1$, then normal-valuation $v = v$.
(46) If $K$ has a valuation, then normal-valuation(normal-valuation $v) = normal-valuation v$.

4. Valuation Ring

Let $K$ be a non empty double loop structure and let $v$ be a valuation of $K$. The functor NonNegElements $v$ is defined as follows:

(Def. 11) \(\text{NonNegElements } v = \{ x \in K : 0 \leq v(x) \}\).

The following four propositions are true:

(47) Let $K$ be a non empty double loop structure, $v$ be a valuation of $K$, and $a$ be an element of $K$. Then $a \in \text{NonNegElements } v$ if and only if $0 \leq v(a)$.
(48) For every non empty double loop structure $K$ and for every valuation $v$ of $K$ holds $\text{NonNegElements } v \subseteq \text{the carrier of } K$.
(49) For every non empty double loop structure $K$ and for every valuation $v$ of $K$ such that $K$ has a valuation holds $0_K \in \text{NonNegElements } v$.
(50) If $K$ has a valuation, then $1_K \in \text{NonNegElements } v$.

Let us consider $K$, $v$. Let us assume that $K$ has a valuation. The functor ValuatRing $v$ yields a strict commutative non degenerated ring and is defined by the conditions (Def. 12).

(Def. 12)(i) The carrier of ValuatRing $v = \text{NonNegElements } v$,
(ii) the addition of ValuatRing $v = (\text{the addition of } K)\mid (\text{NonNegElements } v \times \text{NonNegElements } v)$,
(iii) the multiplication of ValuatRing $v = (\text{the multiplication of } K)\mid (\text{NonNegElements } v \times \text{NonNegElements } v)$,
(iv) the zero of ValuatRing $v = 0_K$, and
(v) the one of ValuatRing $v = 1_K$. 
The following propositions are true:

(51) If $K$ has a valuation, then every element of $\text{ValuatRing} \, v$ is an element of $K$.

(52) If $K$ has a valuation, then $0 \leq v(a)$ iff $a$ is an element of $\text{ValuatRing} \, v$.

(53) If $K$ has a valuation, then for every subset $S$ of $\text{ValuatRing} \, v$ holds $0$ is a lower bound of $v^o S$.

(54) Suppose $K$ has a valuation. Let $x$, $y$ be elements of $K$ and $x_1$, $y_1$ be elements of $\text{ValuatRing} \, v$. If $x = x_1$ and $y = y_1$, then $x + y = x_1 + y_1$.

(55) Suppose $K$ has a valuation. Let $x$, $y$ be elements of $K$ and $x_1$, $y_1$ be elements of $\text{ValuatRing} \, v$. If $x = x_1$ and $y = y_1$, then $x \cdot y = x_1 \cdot y_1$.

(56) If $K$ has a valuation, then $0_{\text{ValuatRing} \, v} = 0_K$.

(57) If $K$ has a valuation, then $1_{\text{ValuatRing} \, v} = 1_K$.

(58) If $K$ has a valuation, then for every element $x$ of $K$ and for every element $y$ of $\text{ValuatRing} \, v$ such that $x = y$ holds $-x = -y$.

(59) If $K$ has a valuation, then $\text{ValuatRing} \, v$ is integral domain-like.

(60) If $K$ has a valuation, then for every element $y$ of $\text{ValuatRing} \, v$ holds $\text{power}_K(y, n) = \text{power}_{\text{ValuatRing} \, v}(y, n)$.

Let us consider $K$, $v$. Let us assume that $K$ has a valuation. The functor $\text{PosElements} \, v$ yields an ideal of $\text{ValuatRing} \, v$ and is defined as follows:

(Def. 13) $\text{PosElements} \, v = \{x \in K : 0 < v(x)\}$.

Let us consider $K$, $v$. We introduce $\text{vp} \, v$ as a synonym of $\text{PosElements} \, v$.

Next we state three propositions:

(61) If $K$ has a valuation, then $a \in \text{vp} \, v$ iff $0 < v(a)$.

(62) If $K$ has a valuation, then $0_K \in \text{vp} \, v$.

(63) If $K$ has a valuation, then $1_K \notin \text{vp} \, v$.

Let us consider $K$, $v$ and let $S$ be a non empty subset of $K$. Let us assume that $K$ has a valuation and $S$ is a subset of $\text{ValuatRing} \, v$. The functor $\text{min}(S, v)$ yielding a subset of $\text{ValuatRing} \, v$ is defined as follows:

(Def. 14) $\text{min}(S, v) = v^{-1}(\{\inf(v \circ S)\}) \cap S$.

The following four propositions are true:

(64) For every non empty subset $S$ of $K$ such that $K$ has a valuation and $S$ is a subset of $\text{ValuatRing} \, v$ holds $\text{min}(S, v) \subseteq S$.

(65) Let $S$ be a non empty subset of $K$. Suppose $K$ has a valuation and $S$ is a subset of $\text{ValuatRing} \, v$. Let $x$ be an element of $K$. Then $x \in \text{min}(S, v)$ if and only if the following conditions are satisfied:

(i) $x \in S$, and

(ii) for every element $y$ of $K$ such that $y \in S$ holds $v(x) \leq v(y)$. 
(66) Suppose $K$ has a valuation. Let $I$ be a non empty subset of $K$ and $x$ be an element of $\text{ValuatRing} v$. If $I$ is an ideal of $\text{ValuatRing} v$ and $x \in \min(I, v)$, then $I = \{x\}$–ideal.

(67) For every non empty double loop structure $R$ holds every add closed non empty subset of $R$ is a set closed w.r.t. the addition of $R$.

Let $R$ be a ring and let $P$ be a right ideal of $R$. A submodule of $\text{RightMod}(R)$ is called a submodule of $P$ if:

(Def. 15) The carrier of it = $P$.

Let $R$ be a ring and let $P$ be a right ideal of $R$. Note that there exists a submodule of $P$ which is strict. Next we state the proposition

(68) Let $R$ be a ring, $P$ be an ideal of $R$, $M$ be a submodule of $P$, $a$ be a binary operation on $P$, $z$ be an element of $P$, and $m$ be a function from $P \times$ the carrier of $R$ into $P$. Suppose $a = (\text{the addition of } R)(P \times P)$ and $m = (\text{the multiplication of } R)(P \times \text{the carrier of } R)$ and $z = \text{the zero of } R$. Then the right module structure of $M = \langle P, a, z, m \rangle$.

Let $R$ be a ring, let $M_1, M_2$ be right modules over $R$, and let $h$ be a function from $M_1$ into $M_2$. We say that $h$ is scalar linear if and only if:

(Def. 16) For every element $x$ of $M_1$ and for every element $r$ of $R$ holds $h(x \cdot r) = h(x) \cdot r$.

Let $R$ be a ring, let $M_1$ be a right module over $R$, and let $M_2$ be a submodule of $M_1$. Observe that $\text{incl}(M_2, M_1)$ is additive and scalar linear.

Next we state a number of propositions:

(69) If $K$ has a valuation and $b$ is an element of $\text{ValuatRing} v$, then $v(a) \leq v(a) + v(b)$.

(70) If $K$ has a valuation and $a$ is an element of $\text{ValuatRing} v$, then $\text{power}_K(a, n)$ is an element of $\text{ValuatRing} v$.

(71) If $K$ has a valuation, then for every element $x$ of $\text{ValuatRing} v$ such that $x \neq 0_K$ holds $\text{power}_K(x, n) \neq 0_K$.

(72) If $K$ has a valuation and $v(a) = 0$, then $a$ is an element of $\text{ValuatRing} v$ and $a^{-1}$ is an element of $\text{ValuatRing} v$.

(73) If $K$ has a valuation and $a \neq 0_K$ and $a$ is an element of $\text{ValuatRing} v$ and $a^{-1}$ is an element of $\text{ValuatRing} v$, then $v(a) = 0$.

(74) If $K$ has a valuation and $v(a) = 0$, then for every ideal $I$ of $\text{ValuatRing} v$ holds $a \in I$ iff $I$ = the carrier of $\text{ValuatRing} v$.

(75) If $K$ has a valuation, then $P_{\text{generator}} v$ is an element of $\text{ValuatRing} v$.

(76) If $K$ has a valuation, then $v p v$ is proper.

(77) If $K$ has a valuation, then for every element $x$ of $\text{ValuatRing} v$ such that $x \notin v p v$ holds $v(x) = 0$.

(78) If $K$ has a valuation, then $v p v$ is prime.
If $K$ has a valuation, then for every proper ideal $I$ of ValuatRing $v$ holds $I \subseteq \text{vp } v$.

If $K$ has a valuation, then $\text{vp } v$ is maximal.

If $K$ has a valuation, then for every maximal ideal $I$ of ValuatRing $v$ holds $I = \text{vp } v$.

If $K$ has a valuation, then $\text{NonNegElements normal-valuation } v = \text{NonNegElements } v$.

If $K$ has a valuation, then $\text{ValuatRing normal-valuation } v = \text{ValuatRing } v$.

References


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Functional Space $C(\Omega)$, $C_0(\Omega)$

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**Summary.** In this article, first we give a definition of a functional space which is constructed from all complex-valued continuous functions defined on a compact topological space. We prove that this functional space is a Banach algebra. Next, we give a definition of a function space which is constructed from all complex-valued continuous functions with bounded support. We also prove that this function space is a complex normed space.

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The terminology and notation used here have been introduced in the following articles: [6], [24], [25], [1], [26], [5], [4], [2], [21], [15], [3], [18], [19], [23], [22], [17], [7], [11], [12], [9], [10], [13], [8], [14], [20], and [16].

Let $X$ be a topological structure and let $f$ be a function from the carrier of $X$ into $\mathbb{C}$. We say that $f$ is continuous if and only if:

(Def. 1) For every subset $Y$ of $\mathbb{C}$ such that $Y$ is closed holds $f^{-1}(Y)$ is closed.

Let $X$ be a 1-sorted structure and let $y$ be a complex number. The functor $X \mapsto y$ yielding a function from the carrier of $X$ into $\mathbb{C}$ is defined by:

(Def. 2) $X \mapsto y = (\text{the carrier of } X) \mapsto y$.

One can prove the following proposition

(1) Let $X$ be a non empty topological space, $y$ be a complex number, and $f$ be a function from the carrier of $X$ into $\mathbb{C}$. If $f = X \mapsto y$, then $f$ is continuous.

Let $X$ be a non empty topological space and let $y$ be a complex number. Observe that $X \mapsto y$ is continuous.

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Let $X$ be a non empty topological space. One can verify that there exists a function from the carrier of $X$ into $\mathbb{C}$ which is continuous.

The following propositions are true:

(2) Let $X$ be a non empty topological space and $f$ be a function from the carrier of $X$ into $\mathbb{C}$. Then $f$ is continuous if and only if for every subset $Y$ of $\mathbb{C}$ such that $Y$ is open holds $f^{-1}(Y)$ is open.

(3) Let $X$ be a non empty topological space and $f$ be a function from the carrier of $X$ into $\mathbb{C}$. Then $f$ is continuous if and only if for every point $x$ of $X$ and for every subset $V$ of $\mathbb{C}$ such that $f(x) \in V$ and $V$ is open there exists a subset $W$ of $X$ such that $x \in W$ and $W$ is open and $f^{-1}W \subseteq V$.

(4) Let $X$ be a non empty topological space and $f, g$ be continuous functions from the carrier of $X$ into $\mathbb{C}$. Then $f + g$ is a continuous function from the carrier of $X$ into $\mathbb{C}$.

(5) Let $X$ be a non empty topological space, $a$ be a complex number, and $f$ be a continuous function from the carrier of $X$ into $\mathbb{C}$. Then $a \cdot f$ is a continuous function from the carrier of $X$ into $\mathbb{C}$.

(6) Let $X$ be a non empty topological space and $f, g$ be continuous functions from the carrier of $X$ into $\mathbb{C}$. Then $f - g$ is a continuous function from the carrier of $X$ into $\mathbb{C}$.

(7) Let $X$ be a non empty topological space and $f, g$ be continuous functions from the carrier of $X$ into $\mathbb{C}$. Then $f \cdot g$ is a continuous function from the carrier of $X$ into $\mathbb{C}$.

(8) Let $X$ be a non empty topological space and $f$ be a continuous function from the carrier of $X$ into $\mathbb{C}$. Then $|f|$ is a function from the carrier of $X$ into $\mathbb{R}$ and $|f|$ is continuous.

Let $X$ be a non empty topological space. The $\mathbb{C}$-continuous functions of $X$ yields a subset of $\mathbb{C}$-Algebra(the carrier of $X$) and is defined by:

(Def. 3) The $\mathbb{C}$-continuous functions of $X = \{ f : f$ ranges over continuous functions from the carrier of $X$ into $\mathbb{C} \}$.

Let $X$ be a non empty topological space. Observe that the $\mathbb{C}$-continuous functions of $X$ is non empty.

Let $X$ be a non empty topological space. Observe that the $\mathbb{C}$-continuous functions of $X$ is $\mathbb{C}$-additively linearly closed and multiplicatively closed.

Let $X$ be a non empty topological space. The $\mathbb{C}$-algebra of continuous functions of $X$ yielding a complex algebra is defined by the condition (Def. 4).

(Def. 4) The $\mathbb{C}$-algebra of continuous functions of $X = (\mathbb{C}$-continuous functions of $X$, $\text{mult}(\mathbb{C}$-continuous functions of $X$, $\mathbb{C}$-Algebra(the carrier of $X$)), $\text{Add}(\mathbb{C}$-continuous functions of $X$, $\mathbb{C}$-Algebra(the carrier of $X$)), $\text{Mult}(\mathbb{C}$-continuous functions of $X$, $\mathbb{C}$-Algebra(the carrier of $X$)), $\text{One}(\mathbb{C}$-continuous functions of $X$, $\mathbb{C}$-Algebra(the carrier of $X$)))).
Next we state the proposition

(9) Let $X$ be a non empty topological space. Then the $\mathbb{C}$-algebra of continuous functions of $X$ is a complex subalgebra of $\mathbb{C}$-algebra\((\text{the carrier of } X)\)).

Let $X$ be a non empty topological space. Observe that the $\mathbb{C}$-algebra of continuous functions of $X$ is strict and non empty.

Let $X$ be a non empty topological space. One can check that the $\mathbb{C}$-algebra of continuous functions of $X$ is Abelian, add-associative, right zeroed, right complementable, vector distributive, scalar distributive, scalar associative, scalar unital, commutative, associative, right unital, right distributive, vector distributive, scalar associative, and vector associative.

Next we state several propositions:

(10) Let $X$ be a non empty topological space, $F, G, H$ be vectors of the $\mathbb{C}$-algebra of continuous functions of $X$, and $f, g, h$ be functions from the carrier of $X$ into $\mathbb{C}$. Suppose $f = F$ and $g = G$ and $h = H$. Then $H = F + G$ if and only if for every element $x$ of the carrier of $X$ holds $h(x) = f(x) + g(x)$.

(11) Let $X$ be a non empty topological space, $F, G$ be vectors of the $\mathbb{C}$-algebra of continuous functions of $X$, $f, g$ be functions from the carrier of $X$ into $\mathbb{C}$, and $a$ be a complex number. Suppose $f = F$ and $g = G$. Then $G = a \cdot F$ if and only if for every element $x$ of $X$ holds $g(x) = a \cdot f(x)$.

(12) Let $X$ be a non empty topological space, $F, G, H$ be vectors of the $\mathbb{C}$-algebra of continuous functions of $X$, and $f, g, h$ be functions from the carrier of $X$ into $\mathbb{C}$. Suppose $f = F$ and $g = G$ and $h = H$. Then $H = F \cdot G$ if and only if for every element $x$ of the carrier of $X$ holds $h(x) = f(x) \cdot g(x)$.

(13) For every non empty topological space $X$ holds \(0\) the $\mathbb{C}$-algebra of continuous functions of $X = X \mapsto 0_{\mathbb{C}}$.

(14) For every non empty topological space $X$ holds \(1\) the $\mathbb{C}$-algebra of continuous functions of $X = X \mapsto 1_{\mathbb{C}}$.

(15) Let $A$ be a complex algebra and $A_1, A_2$ be complex subalgebras of $A$. Suppose the carrier of $A_1 \subseteq$ the carrier of $A_2$. Then $A_1$ is a complex subalgebra of $A_2$.

(16) Let $X$ be a non empty compact topological space. Then the $\mathbb{C}$-algebra of continuous functions of $X$ is a complex subalgebra of the $\mathbb{C}$-algebra of bounded functions of the carrier of $X$.

Let $X$ be a non empty compact topological space. The $\mathbb{C}$-continuous functions norm of $X$ yields a function from the $\mathbb{C}$-continuous functions of $X$ into $\mathbb{R}$ and is defined by:

\[ X \mapsto 0_{\mathbb{C}}. \]
(Def. 5) The C-continuous functions norm of $X = (\text{C-BoundedFunctionsNorm (the carrier of } X)) | \text{the C-continuous functions of } X$.

Let $X$ be a non empty compact topological space. The C-normed algebra of continuous functions of $X$ yields a normed complex algebra structure and is defined by the condition (Def. 6).

(Def. 6) The C-normed algebra of continuous functions of $X = \langle \text{the C-continuous functions of } X, \text{mult}(\text{the C-continuous functions of } X, C-\text{Algebra}(\text{the carrier of } X)), \text{Add}(\text{the C-continuous functions of } X, C-\text{Algebra}(\text{the carrier of } X)), \text{Mult}(\text{the C-continuous functions of } X, C-\text{Algebra}(\text{the carrier of } X)), \text{One}(\text{the C-continuous functions of } X, C-\text{Algebra}(\text{the carrier of } X)), \text{Zero}(\text{the C-continuous functions of } X, C-\text{Algebra}(\text{the carrier of } X)), \text{the C-continuous functions norm of } X \rangle$.

Let $X$ be a non empty compact topological space. Note that the C-normed algebra of continuous functions of $X$ is non empty and strict.

Let $X$ be a non empty compact topological space. Observe that the C-normed algebra of continuous functions of $X$ is unital.

Next we state the proposition

(17) Let $X$ be a non empty compact topological space. Then the C-normed algebra of continuous functions of $X$ is a complex algebra.

Let $X$ be a non empty compact topological space. One can check that the C-normed algebra of continuous functions of $X$ is right complementable, Abelian, add-associative, right zeroed, vector distributive, scalar distributive, scalar associative, associative, commutative, right distributive, right unital, and vector associative.

One can prove the following proposition

(18) Let $X$ be a non empty compact topological space and $F$ be a point of the C-normed algebra of continuous functions of $X$. Then $(\text{Mult}(\text{the C-continuous functions of } X, C-\text{Algebra}(\text{the carrier of } X))(1_C, F) = F$.

Let $X$ be a non empty compact topological space. Observe that the C-normed algebra of continuous functions of $X$ is vector distributive, scalar distributive, scalar associative, and scalar unital.

We now state a number of propositions:

(19) Let $X$ be a non empty compact topological space. Then the C-normed algebra of continuous functions of $X$ is a complex linear space.

(20) Let $X$ be a non empty compact topological space. Then $X \mapsto 0 = _{0\text{the C-normed algebra of continuous functions of } X}$.

(21) Let $X$ be a non empty compact topological space and $F$ be a point of the C-normed algebra of continuous functions of $X$. Then $0 \leq \|F\|$.  

(22) Let $X$ be a non empty compact topological space, $f$, $g$, $h$ be functions from the carrier of $X$ into $C$, and $F$, $G$, $H$ be points of the C-normed
algebra of continuous functions of $X$. Suppose $f = F$ and $g = G$ and $h = H$. Then $H = F + G$ if and only if for every element $x$ of $X$ holds $h(x) = f(x) + g(x)$.

(23) Let $a$ be a complex number, $X$ be a non empty compact topological space, $f, g$ be functions from the carrier of $X$ into $\mathbb{C}$, and $F, G, H$ be points of the $\mathbb{C}$-normed algebra of continuous functions of $X$. Suppose $f = F$ and $g = G$. Then $H = a \cdot F$ if and only if for every element $x$ of $X$ holds $h(x) = a \cdot f(x)$.

(24) Let $X$ be a non empty compact topological space, $f, g, h$ be functions from the carrier of $X$ into $\mathbb{C}$, and $F, G, H$ be points of the $\mathbb{C}$-normed algebra of continuous functions of $X$. Suppose $f = F$ and $g = G$ and $h = H$. Then $H = F \cdot G$ if and only if for every element $x$ of $X$ holds $h(x) = f(x) \cdot g(x)$.

(25) Let $X$ be a non empty compact topological space. Then $\|0\|_{\mathbb{C}} = 0$.

(26) Let $X$ be a non empty compact topological space and $F$ be a point of the $\mathbb{C}$-normed algebra of continuous functions of $X$. Suppose $\|F\| = 0$. Then $F = 0$.

(27) Let $a$ be a complex number, $X$ be a non empty compact topological space, and $F$ be a point of the $\mathbb{C}$-normed algebra of continuous functions of $X$. Then $\|a \cdot F\| = |a| \cdot \|F\|$.

(28) Let $X$ be a non empty compact topological space and $F, G$ be points of the $\mathbb{C}$-normed algebra of continuous functions of $X$. Then $\|F + G\| \leq \|F\| + \|G\|$.

Let $X$ be a non empty compact topological space. Observe that the $\mathbb{C}$-normed algebra of continuous functions of $X$ is discernible, reflexive, and complex normed space-like.

The following propositions are true:

(29) Let $X$ be a non empty compact topological space, $f, g, h$ be functions from the carrier of $X$ into $\mathbb{C}$, and $F, G, H$ be points of the $\mathbb{C}$-normed algebra of continuous functions of $X$. Suppose $f = F$ and $g = G$ and $h = H$. Then $H = F - G$ if and only if for every element $x$ of $X$ holds $h(x) = f(x) - g(x)$.

(30) Let $X$ be a complex Banach space, $Y$ be a subset of $X$, and $s_1$ be a sequence of $X$. Suppose $Y$ is closed and $\text{rng } s_1 \subseteq Y$ and $s_1$ is $\mathbb{C}$-Cauchy. Then $s_1$ is convergent and $\lim s_1 \in Y$.

(31) Let $X$ be a non empty compact topological space and $Y$ be a subset of the $\mathbb{C}$-normed algebra of bounded functions of the carrier of $X$. If $Y$ is the $\mathbb{C}$-continuous functions of $X$, then $Y$ is closed.

(32) Let $X$ be a non empty compact topological space and $s_1$ be a sequence
of the \(\mathbb{C}\)-normed algebra of continuous functions of \(X\). If \(s_1\) is \(\mathbb{C}\)-Cauchy, then \(s_1\) is convergent.

Let \(X\) be a non empty compact topological space. One can verify that the \(\mathbb{C}\)-normed algebra of continuous functions of \(X\) is complete.

Let \(X\) be a non empty compact topological space. Observe that the \(\mathbb{C}\)-normed algebra of continuous functions of \(X\) is Banach Algebra-like.

Next we state three propositions:

(33) For every non empty topological space \(X\) and for all functions \(f, g\) from the carrier of \(X\) into \(\mathbb{C}\) holds \(\text{support}(f + g) \subseteq \text{support } f \cup \text{support } g\).

(34) Let \(X\) be a non empty topological space, \(a\) be a complex number, and \(f\) be a function from the carrier of \(X\) into \(\mathbb{C}\). Then \(\text{support}(a \cdot f) \subseteq \text{support } f\).

(35) For every non empty topological space \(X\) and for all functions \(f, g\) from the carrier of \(X\) into \(\mathbb{C}\) holds \(\text{support}(f \cdot g) \subseteq \text{support } f \cup \text{support } g\).

Let \(X\) be a non empty topological space. The \(CC_0\)-functions of \(X\) yielding a non empty subset of the \(\mathbb{C}\)-vector space of the carrier of \(X\) is defined by the condition (Def. 7).

(Def. 7) The \(CC_0\)-functions of \(X = \{f; f\ \text{ranges over functions from the carrier of } X \text{ into } \mathbb{C}; f \text{ is continuous } \land \bigvee_{Y; \text{non empty subset of } X} (Y \text{ is compact } \land \bigwedge_{A; \text{subset of } X} (A = \text{support } f \Rightarrow \overline{A} \text{ is a subset of } Y))\}.

The following propositions are true:

(36) Let \(X\) be a non empty topological space. Then the \(CC_0\)-functions of \(X\) is a non empty subset of \(\mathbb{C}\)-Algebra(\(X\)).

(37) Let \(X\) be a non empty topological space and \(W\) be a non empty subset of \(\mathbb{C}\)\(-\text{Algebra}(\(X\)). Suppose \(W = \text{the } CC_0\)-functions of \(X\). Then \(W\) is \(\mathbb{C}\)-additively linearly closed.

(38) For every non empty topological space \(X\) holds the \(CC_0\)-functions of \(X\) is add closed.

(39) For every non empty topological space \(X\) holds the \(CC_0\)-functions of \(X\) is linearly closed.

Let \(X\) be a non empty topological space. Observe that the \(CC_0\)-functions of \(X\) is non empty and linearly closed.

The following propositions are true:

(40) Let \(V\) be a complex linear space and \(V_1\) be a subset of \(V\). Suppose \(V_1\) is linearly closed and \(V_1\) is not empty. Then \(\langle V_1, \text{Zero}(V_1, V), \text{Add}(V_1, V), \text{Mult}(V_1, V) \rangle\) is a subspace of \(V\).

(41) Let \(X\) be a non empty topological space. Then \(\langle \text{the } CC_0\)-functions of \(X, \text{Zero}(\text{the } CC_0\)-functions of \(X), \text{the } \mathbb{C}\)-vector space of the carrier of \(X), \text{Add}(\text{the } CC_0\)-functions of \(X), \text{the } \mathbb{C}\)-vector space of the carrier of \(X), \text{Mult}(\text{the } CC_0\)-functions of \(X), \text{the } \mathbb{C}\)-vector space of the carrier of \(X) \rangle\) is a subspace of the \(\mathbb{C}\)-vector space of the carrier of \(X\).
Let $X$ be a non empty topological space. The $\mathbb{C}$-vector space of $C_0$-functions of $X$ yielding a complex linear space is defined by the condition (Def. 8).

(Def. 8) The $\mathbb{C}$-vector space of $C_0$-functions of $X = \langle$ the $C_0$-functions of $X$, Zero(the $C_0$-functions of $X$, the $\mathbb{C}$-vector space of the carrier of $X$), Add(the $C_0$-functions of $X$, the $\mathbb{C}$-vector space of the carrier of $X$), Mult(the $C_0$-functions of $X$, the $\mathbb{C}$-vector space of the carrier of $X$)$\rangle$.

Next we state the proposition

(42) Let $X$ be a non empty topological space and $x$ be a set. If $x \in$ the $C_0$-functions of $X$, then $x \in \mathbb{C}$-BoundedFunctions (the carrier of $X$).

Let $X$ be a non empty topological space. The $C_0$-functions norm of $X$ yielding a function from the $C_0$-functions of $X$ into $\mathbb{R}$ is defined by:

(Def. 9) The $C_0$-functions norm of $X = \langle$ $C$-BoundedFunctionsNorm (the carrier of $X$)$\rangle$.

Let $X$ be a non empty topological space. The $\mathbb{C}$-normed space of $C_0$-functions of $X$ yielding a complex normed space structure is defined by the condition (Def. 10).

(Def. 10) The $\mathbb{C}$-normed space of $C_0$-functions of $X = \langle$ the $C_0$-functions of $X$, Zero(the $C_0$-functions of $X$, the $\mathbb{C}$-vector space of the carrier of $X$), Add(the $C_0$-functions of $X$, the $\mathbb{C}$-vector space of the carrier of $X$), Mult(the $C_0$-functions of $X$, the $\mathbb{C}$-vector space of the carrier of $X$), the $C_0$-functions norm of $X$)$\rangle$.

Let $X$ be a non empty topological space. One can check that the $\mathbb{C}$-normed space of $C_0$-functions of $X$ is strict and non empty.

One can prove the following propositions:

(43) Let $X$ be a non empty topological space and $x$ be a set. Suppose $x \in$ the $C_0$-functions of $X$. Then $x \in \mathbb{C}$-continuous functions of $X$.

(44) For every non empty topological space $X$ holds

$0$ the $\mathbb{C}$-vector space of $C_0$-functions of $X = X \mapsto 0$.

(45) For every non empty topological space $X$ holds

$0$ the $\mathbb{C}$-normed space of $C_0$-functions of $X = X \mapsto 0$.

(46) Let $a$ be a complex number, $X$ be a non empty topological space, and $F, G$ be points of the $\mathbb{C}$-normed space of $C_0$-functions of $X$. Then $\|F\| = 0$ iff $F = 0$ the $\mathbb{C}$-normed space of $C_0$-functions of $X$ and $\|a \cdot F\| = |a| \cdot \|F\|$ and $\|F + G\| \leq \|F\| + \|G\|$.

Let $X$ be a non empty topological space. Note that the $\mathbb{C}$-normed space of $C_0$-functions of $X$ is reflexive, discernible, complex normed space-like, vector distributive, scalar distributive, scalar associative, scalar unital, Abelian, add-associative, right zeroed, and right complementable.

The following proposition is true
Let $X$ be a non empty topological space. Then the $C$-normed space of $C_0$-functions of $X$ is a complex normed space.

References


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Summary. We introduce length-preserving linear transformations of Euclidean topological spaces. We also introduce rotation which preserves orientation (proper rotation) and reverses orientation (improper rotation). We show that every rotation that preserves orientation can be represented as a composition of base proper rotations. And finally, we show that every rotation that reverses orientation can be represented as a composition of proper rotations and one improper rotation.

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The papers [11], [35], [36], [8], [10], [9], [3], [7], [14], [2], [30], [4], [19], [12], [31], [24], [34], [13], [22], [17], [1], [20], [15], [16], [40], [38], [33], [25], [28], [37], [23], [6], [39], [18], [21], [32], [5], [26], [29], and [27] provide the terminology and notation for this paper.

1. Preliminaries

We adopt the following rules: $x$, $X$ are sets, $\alpha$, $\alpha_1$, $\alpha_2$, $r$, $s$ are real numbers, and $i$, $j$, $k$, $m$, $n$ are natural numbers.

We now state three propositions:

1. Let $K$ be a field, $M$ be a square matrix over $K$ of dimension $n$, and $P$ be a permutation of $\text{Seg } n$. Then $\det((M \cdot P)^T \cdot P) = \det M$ and for all $i$, $j$ such that $(i, j) \in$ the indices of $M$ holds $((M \cdot P)^T \cdot P)_{i,j} = MP(i), p(j)$.

2. For every field $K$ and for every diagonal square matrix $M$ over $K$ of dimension $n$ holds $M^T = M$. 
For every real-valued finite sequence $f$ and for every $i$ such that $i \in \text{dom} f$ holds $\sum (f + (i, r)) = (\sum f - f(i))^2 + r^2$.

Let us consider $X$ and let $F$ be a function yielding function. We say that $F$ is $X$-support-yielding if and only if:

(Def. 1) For every function $f$ and for every $x$ such that $f \in \text{dom} F$ and $F(f)(x) \neq f(x)$ holds $x \in X$.

Let us consider $X$. One can check that there exists a function yielding function which is $X$-support-yielding.

Let us consider $X$ and let $Y$ be a subset of $X$. One can check that every function yielding function which is $Y$-support-yielding is also $X$-support-yielding.

Let $X, Y$ be sets. Note that every function yielding function which is $X$-support-yielding and $Y$-support-yielding is also $X \cap Y$-support-yielding. Let $f$ be an $X$-support-yielding function yielding function and let $g$ be a $Y$-support-yielding function yielding function. Note that $f \cdot g$ is $X \cup Y$-support-yielding.

Let us consider $n$. Observe that there exists a function from $E^n_T$ into $E^n_T$ which is homogeneous.

Let us consider $n, m$. Observe that every function from $E^n_T$ into $E^m_T$ is finite sequence-yielding.

Let us consider $n$ and let $A$ be a square matrix over $\mathbb{R}_F$ of dimension $n \times m$. One can check that $Mx2Tran A$ is additive.

Let us consider $n$ and let $A$ be a square matrix over $\mathbb{R}_F$ of dimension $n$. Note that $Mx2Tran A$ is homogeneous.

Let us consider $n$ and let $f, g$ be homogeneous functions from $E^n_T$ into $E^n_T$. Note that $f \cdot g$ is homogeneous.

2. Improper Rotation

In the sequel $p, q$ are points of $E^n_T$.

Let us consider $n, i$. Let us assume that $i \in \text{Seg} n$. The axial symmetry of $i$ and $n$ yields an invertible square matrix over $\mathbb{R}_F$ of dimension $n$ and is defined by the conditions (Def. 2).

(Def. 2)(i) (The axial symmetry of $i$ and $n$)$_{i,i} = -1_{\mathbb{R}_F}$, and

(ii) for all $k, m$ such that $(k, m) \in$ the indices of the axial symmetry of $i$ and $n$ holds if $k = m$ and $k \neq i$, then (the axial symmetry of $i$ and $n$)$_{k,k} = 1_{\mathbb{R}_F}$ and if $k \neq m$, then (the axial symmetry of $i$ and $n$)$_{k,m} = 0_{\mathbb{R}_F}$.

The following propositions are true:

(4) If $i \in \text{Seg} n$, then $\text{Det} (\text{the axial symmetry of } i \text{ and } n) = -1_{\mathbb{R}_F}$.

(5) If $i, j \in \text{Seg} n$ and $i \neq j$, then $^{(p)} (\text{the axial symmetry of } i \text{ and } n) \Box_{i,j} = p(j)$.

(6) If $i \in \text{Seg} n$, then $^{(p)} (\text{the axial symmetry of } i \text{ and } n) \Box_{i,i} = -p(i)$. 


Suppose \( i \in \text{Seg} \, n \). Then

(i) the axial symmetry of \( i \) and \( n \) is diagonal, and

(ii) \((\text{the axial symmetry of } i \text{ and } n)^\sim = \text{the axial symmetry of } i \text{ and } n\)\.

If \( i \in \text{Seg} \, n \) and \( i \neq j \), then \((\text{Mx2Tran} \, (\text{the axial symmetry of } i \text{ and } n))(p)(j) = p(j)\).

If \( i \in \text{Seg} \, n \), then \((\text{Mx2Tran} \, (\text{the axial symmetry of } i \text{ and } n))(p)(i) = -p(i)\).

If \( i \in \text{Seg} \, n \), then \((\text{Mx2Tran} \, (\text{the axial symmetry of } i \text{ and } n))(p) = p + (i, -p(i))\).

If \( i \in \text{Seg} \, n \), then \((\text{Mx2Tran} \, (\text{the axial symmetry of } i \text{ and } n))\) is \{i\}-support-yielding.

For all elements \( a, b \) of \( \mathbb{R}_F \) such that \( a = \cos \rho \) and \( b = \sin \rho \) holds

\[
\det (\text{the } 0_{\mathbb{R}_F}\text{-block diagonal of } \left(\begin{array}{cc}
a & b \\
b & a
d\end{array}\right), I_{n \times n}) = 1_{\mathbb{R}_F}.
\]

### 3. Proper Rotation

Let us consider \( n, \alpha \) and let us consider \( i, j \). Let us assume that \( 1 \leq i < j \leq n \). The functor \( \text{Rotation}(i, j, n, \alpha) \) yielding an invertible square matrix over \( \mathbb{R}_F \) of dimension \( n \) is defined by the conditions (Def. 3).

(Def. 3)(i) \((\text{Rotation}(i, j, n, \alpha))_{i,i} = \cos \alpha\),

(ii) \((\text{Rotation}(i, j, n, \alpha))_{j,j} = \cos \alpha\),

(iii) \((\text{Rotation}(i, j, n, \alpha))_{i,j} = \sin \alpha\),

(iv) \((\text{Rotation}(i, j, n, \alpha))_{j,i} = -\sin \alpha\), and

(v) for all \( k, m \) such that \( \{k, m\} \in \text{the indices of } \text{Rotation}(i, j, n, \alpha) \) holds

if \( k = m \) and \( k \neq i \) and \( k \neq j \), then \((\text{Rotation}(i, j, n, \alpha))_{k,k} = 1_{\mathbb{R}_F} \) and if \( k \neq m \) and \( \{k, m\} \neq \{i, j\} \), then \((\text{Rotation}(i, j, n, \alpha))_{k,m} = 0_{\mathbb{R}_F} \).

We now state a number of propositions:

(13) If \( 1 \leq i < j \leq n \), then \( \det \, \text{ Rotation}(i, j, n, \alpha) = 1_{\mathbb{R}_F} \).

(14) If \( 1 \leq i < j \leq n \) and \( k \in \text{Seg} \, n \) and \( k \neq i \) and \( k \neq j \), then \((^\alpha p) \cdot (\text{Rotation}(i, j, n, \alpha)) \square k = p(k)\).

(15) If \( 1 \leq i < j \leq n \), then \((^\alpha p) \cdot (\text{Rotation}(i, j, n, \alpha)) \square i = p(i) \cdot \cos \alpha + p(j) \cdot -\sin \alpha\).

(16) If \( 1 \leq i < j \leq n \), then \((^\alpha p) \cdot (\text{Rotation}(i, j, n, \alpha)) \square j = p(i) \cdot \sin \alpha + p(j) \cdot \cos \alpha\).

(17) If \( 1 \leq i < j \leq n \), then \((\text{Rotation}(i, j, n, \alpha_1) \cdot \text{Rotation}(i, j, n, \alpha_2) = \text{Rotation}(i, j, n, \alpha_1 + \alpha_2)\).

(18) If \( 1 \leq i < j \leq n \), then \((\text{Rotation}(i, j, n, 0) = I_{n \times n} \).

(19) If \( 1 \leq i < j \leq n \), then \((\text{Rotation}(i, j, n, \alpha) \) is orthogonal and \((\text{Rotation}(i, j, n, \alpha))^\sim = \text{Rotation}(i, j, n, -\alpha)\).
(20) If $1 \leq i < j \leq n$ and $k \neq i$ and $k \neq j$, then

$$(\text{Mx2Tran Rotation}(i, j, n, \alpha))(p)(k) = p(k).$$

(21) If $1 \leq i < j \leq n$, then $(\text{Mx2Tran Rotation}(i, j, n, \alpha))(p)(i) = p(i) \cdot \cos \alpha + p(j) \cdot -\sin \alpha$.

(22) If $1 \leq i < j \leq n$, then $(\text{Mx2Tran Rotation}(i, j, n, \alpha))(p)(j) = p(i) \cdot \sin \alpha + p(j) \cdot \cos \alpha$.

(23) If $1 \leq i < j \leq n$, then $(\text{Mx2Tran Rotation}(i, j, n, \alpha))(p)(i) = p(i) \cdot \cos \alpha + p(j) \cdot -\sin \alpha$.

(24) If $1 \leq i < j \leq n$, then $(\text{Mx2Tran Rotation}(i, j, n, \alpha))(p)(j) = p(i) \cdot \sin \alpha + p(j) \cdot \cos \alpha$.

(25) If $1 \leq i < j \leq n$, then $s^2 \leq p(i)^2 + p(j)^2$, then there exists $\alpha$ such that

$$(\text{Mx2Tran Rotation}(i, j, n, \alpha))(p)(i) = s.$$ 

(26) If $1 \leq i < j \leq n$, then $\text{Mx2Tran Rotation}(i, j, n, \alpha)$ is $\{i, j\}$-support-yielding.

4. LENGTH-PRESERVING LINEAR TRANSFORMATIONS

Let us consider $n$ and let $f$ be a function from $E^n_T$ into $E^n_T$. We say that $f$ is rotation if and only if:

(Def. 4) $|p| = |f(p)|$.

One can prove the following proposition

(27) If $i \in \text{Seg} n$, then Mx2Tran (the axial symmetry of $i$ and $n$) is rotation.

Let us consider $n$ and let $f$ be a function from $E^n_T$ into $E^n_T$. We say that $f$ is base rotation if and only if the condition (Def. 5) is satisfied.

(Def. 5) There exists a finite sequence $F$ of elements of the semigroup of functions onto the carrier of $E^n_T$ such that $f = \prod F$ and for every $k$ such that $k \in \text{dom} F$ there exist $i, j, r$ such that $1 \leq i < j \leq n$ and $F(k) = \text{Mx2Tran Rotation}(i, j, n, r)$.

Let us consider $n$. One can check that $\text{id}_{E^n_T}$ is base rotation.

Let us consider $n$. One can check that there exists a function from $E^n_T$ into $E^n_T$ which is base rotation.

Let us consider $n$ and let $f, g$ be base rotation functions from $E^n_T$ into $E^n_T$. One can check that $f \cdot g$ is base rotation.

Next we state the proposition

(28) If $1 \leq i < j \leq n$, then $\text{Mx2Tran Rotation}(i, j, n, r)$ is base rotation.

Let us consider $n$. Observe that every function from $E^n_T$ into $E^n_T$ which is base rotation is also homogeneous, additive, rotation, and homeomorphism.

Let us consider $n$ and let $f$ be a base rotation function from $E^n_T$ into $E^n_T$. Note that $f^{-1}$ is base rotation.
Let us consider $n$ and let $f, g$ be rotation functions from $E^n_T$ into $E^n_T$. One can check that $f \cdot g$ is rotation.

In the sequel $f, f_1, f_2$ are homogeneous additive functions from $E^n_T$ into $E^n_T$.

Let us consider $n$ and let us consider $f$. The functor $\text{AutMt} f$ yields a square matrix over $\mathbb{R}_F$ of dimension $n$ and is defined as follows:

(Def. 6) $f = Mx2T\text{ran} \text{AutMt} f$.

Next we state several propositions:

(29) $\text{AutMt}(f_1 \cdot f_2) = \text{AutMt} f_2 \cdot \text{AutMt} f_1$.

(30) Suppose $k \in X$ and $k \in \text{Seg} n$. Then there exists $f$ such that

(i) $f$ is $X$-support-yielding and base rotation,

(ii) if $X \cap \text{Seg} n > 1$, then $f(p)(k) \geq 0$, and

(iii) for every $i$ such that $i \in X \cap \text{Seg} n$ and $i \neq k$ holds $f(p)(i) = 0$.

(31) For every subset $A$ of $E^n_T$ such that $f|A = \text{id}_A$ holds $f|\text{Lin}(A) = \text{id}_{\text{Lin}(A)}$.

(32) Let $A$ be a subset of $E^n_T$. Suppose $f$ is rotation and $f|A = \text{id}_A$. Let given $i$. Suppose $i \in \text{Seg} n$ and the base finite sequence of $n$ and $i \in \text{Lin}(A)$. Then $f(p)(i) = p(i)$.

(33) Let $f$ be a rotation function from $E^n_T$ into $E^n_T$. Suppose $f$ is $X$-support-yielding and for every $i$ such that $i \in X \cap \text{Seg} n$ holds $p(i) = 0$. Then $f(p) = p$.

(34) If $i \in \text{Seg} n$ and $n \geq 2$, then there exists $f$ such that $f$ is base rotation and $f(p) = p + (i, -p(i))$.

(35) If $f$ is \{i\}-support-yielding and rotation, then $\text{AutMt} f = \text{the axial symmetry of } i$ and $n$ or $\text{AutMt} f = I_{\mathbb{R}_F}$.

(36) If $f_1$ is rotation, then there exists $f_2$ such that $f_2$ is base rotation and $f_2 \cdot f_1$ is \{n\}-support-yielding.

5. Rotation Matrix Classification

The following three propositions are true:

(37) If $f$ is rotation, then $\text{Det AutMt } f = 1_{\mathbb{R}_F}$ iff $f$ is base rotation.

(38) If $f$ is rotation, then $\text{Det AutMt } f = 1_{\mathbb{R}_F}$ or $\text{Det AutMt } f = -1_{\mathbb{R}_F}$.

(39) If $f_1$ is rotation and $\text{Det AutMt } f_1 = -1_{\mathbb{R}_F}$ and $i \in \text{Seg} n$ and $\text{AutMt } f_2 = \text{the axial symmetry of } i$ and $n$, then $f_1 \cdot f_2$ is base rotation.

Let us consider $n$ and let $f$ be a rotation homogeneous additive function from $E^n_T$ into $E^n_T$. One can check that $\text{AutMt } f$ is orthogonal.

Let us consider $n$. One can verify that every function from $E^n_T$ into $E^n_T$ which is homogeneous, additive, and rotation is also homeomorphism.
6. The Rotation Mapping a Given Point to Another Point

One can prove the following propositions:

(40) Suppose $n = 1$ and $|p| = |q|$. Then there exists $f$ such that $f$ is rotation and $f(p) = q$ either $\text{AutMt} f = \text{the axial symmetry of } n$ and $n$ or $\text{AutMt} f = I_{S_F}^{n \times n}$.

(41) If $n \neq 1$ and $|p| = |q|$, then there exists $f$ such that $f$ is base rotation and $f(p) = q$.

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Differentiable Functions on Normed Linear Spaces

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Summary. In this article, we formalize differentiability of functions on normed linear spaces. Partial derivative, mean value theorem for vector-valued functions, continuous differentiability, etc. are formalized. As it is well known, there is no exact analog of the mean value theorem for vector-valued functions. However a certain type of generalization of the mean value theorem for vector-valued functions is obtained as follows: If $\|f'(x + t \cdot h)\|$ is bounded for $t$ between 0 and 1 by some constant $M$, then $\|f(x + t \cdot h) - f(x)\| \leq M \cdot \|h\|$. This theorem is called the mean value theorem for vector-valued functions. By this theorem, the relation between the (total) derivative and the partial derivatives of a function is derived [23].

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The notation and terminology used here have been introduced in the following papers: [28], [29], [9], [4], [30], [12], [10], [25], [11], [1], [2], [26], [7], [3], [5], [8], [17], [22], [20], [27], [21], [31], [14], [24], [18], [16], [15], [19], [13], and [6].

1. Preliminaries

In this paper $r$ is a real number and $S$, $T$ are non trivial real normed spaces. Next we state several propositions:

1. Let $R$ be a function from $\mathbb{R}$ into $S$. Then $R$ is rest-like if and only if for every real number $r$ such that $r > 0$ there exists a real number $d$ such that $d > 0$ and for every real number $z$ such that $z \neq 0$ and $|z| < d$ holds $|z|^{-1} \cdot \|Rz\| < r$.

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(2) Let $R$ be a rest of $S$. Suppose $R_0 = 0_S$. Let $e$ be a real number. Suppose $e > 0$. Then there exists a real number $d$ such that $d > 0$ and for every real number $h$ such that $|h| < d$ holds $\|R_h\| \leq e \cdot |h|$.

(3) For every rest $R$ of $S$ and for every bounded linear operator $L$ from $S$ into $T$ holds $L \cdot R$ is a rest of $T$.

(4) Let $R_1$ be a rest of $S$. Suppose $(R_1)_0 = 0_S$. Let $R_2$ be a rest of $S$, $T$. If $(R_2)_0 = 0_T$, then for every linear $L$ of $S$ holds $R_2 \cdot (L + R_1)$ is a rest of $T$.

(5) Let $R_1$ be a rest of $S$. Suppose $(R_1)_0 = 0_S$. Let $R_2$ be a rest of $S$, $T$. Suppose $(R_2)_0 = 0_T$. Let $L_1$ be a linear of $S$ and $L_2$ be a bounded linear operator from $S$ into $T$. Then $L_2 \cdot R_1 + R_2 \cdot (L_1 + R_1)$ is a rest of $T$.

(6) Let $x_0$ be an element of $\mathbb{R}$ and $g$ be a partial function from $\mathbb{R}$ to the carrier of $S$. Suppose $g$ is differentiable in $x_0$. Let $f$ be a partial function from the carrier of $S$ to the carrier of $T$. Suppose $f$ is differentiable in $g_{x_0}$. Then $f \cdot g$ is differentiable in $x_0$ and $(f \cdot g)'(x_0) = f'(g_{x_0})(g'(x_0))$.

(7) Let $S$ be a real normed space, $x_1$ be a finite sequence of elements of $S$, and $y_1$ be a finite sequence of elements of $\mathbb{R}$. Suppose $\text{len} \cdot x_1 = \text{len} \cdot y_1$ and for every element $i$ of $\mathbb{N}$ such that $i \in \text{dom} \cdot x_1$ holds $y_1(i) = \|x_1(i)\|$. Then $\sum \|x_1\| \leq \sum \|y_1\|$.  

(8) Let $S$ be a real normed space, $x$ be a point of $S$, and $N_1$, $N_2$ be neighbourhoods of $x$. Then $N_1 \cap N_2$ is a neighbourhood of $x$.

(9) For every non-empty finite sequence $X$ and for every set $x$ such that $x \in \prod X$ holds $x$ is a finite sequence.

Let $G$ be a real norm space sequence. One can verify that $\prod G$ is constituted finite sequences.

Let $G$ be a real linear space sequence, let $z$ be an element of $\prod \overline{G}$, and let $j$ be an element of $\text{dom} \cdot G$. Then $z(j)$ is an element of $G(j)$.

One can prove the following propositions:

(10) The carrier of $\prod G = \prod \overline{G}$.

(11) Let $i$ be an element of $\text{dom} \cdot G$, $r$ be a set, and $x$ be a function. If $r \in \text{carrier of } G(i)$ and $x \in \prod \overline{G}$, then $x \cdot (i, r) \in \text{carrier of } \prod G$.

Let $G$ be a real norm space sequence. We say that $G$ is nontrivial if and only if:

(Def. 1) For every element $j$ of $\text{dom} \cdot G$ holds $G(j)$ is non trivial.

Let us mention that there exists a real norm space sequence which is non-trivial.

Let $G$ be a nontrivial real norm space sequence and let $i$ be an element of $\text{dom} \cdot G$. Note that $G(i)$ is non trivial.

Let $G$ be a nontrivial real norm space sequence. Note that $\prod G$ is non trivial.

The following propositions are true:
(12) Let $G$ be a real norm space sequence, $p$, $q$ be points of $\prod G$, and $r_0$, $p_0$, $q_0$ be elements of $\prod G$. Suppose $p = p_0$ and $q = q_0$. Then $p + q = r_0$ if and only if for every element $i$ of dom $G$ holds $r_0(i) = p_0(i) + q_0(i)$.

(13) Let $G$ be a real norm space sequence, $p$ be a point of $\prod G$, $r$ be a real number, and $r_0$, $p_0$ be elements of $\prod G$. Suppose $p = p_0$. Then $r \cdot p = r_0$ if and only if for every element $i$ of dom $G$ holds $r_0(i) = r \cdot p_0(i)$.

(14) Let $G$ be a real norm space sequence and $p_0$ be an element of $\prod G$. Then $0\prod G = p_0$ if and only if for every element $i$ of dom $G$ holds $p_0(i) = 0_{G(i)}$.

(15) Let $G$ be a real norm space sequence, $p$, $q$ be points of $\prod G$, and $r_0$, $p_0$, $q_0$ be elements of $\prod G$. Suppose $p = p_0$ and $q = q_0$. Then $p - q = r_0$ if and only if for every element $i$ of dom $G$ holds $r_0(i) = p_0(i) - q_0(i)$.

2. Mean Value Theorem for Vector-Valued Functions

Let $S$ be a real linear space and let $p$, $q$ be points of $S$. The functor $\{p, q\}$ yielding a subset of $S$ is defined as follows:

(Def. 2) $\{p, q\} = \{p + t \cdot (q - p) ; t \text{ ranges over real numbers: } 0 < t \land t < 1\}$.

Let $S$ be a real linear space and let $p$, $q$ be points of $S$. We introduce $[p, q]$ as a synonym of $L(p, q)$.

Next we state several propositions:

(16) For every real linear space $S$ and for all points $p$, $q$ of $S$ holds $\{p, q\} \subseteq [p, q]$.

(17) Let $T$ be a non trivial real normed space and $R$ be a partial function from $\mathbb{R}$ to $T$. Suppose $R$ is total. Then $R$ is rest-like if and only if for every real number $r$ such that $r > 0$ there exists a real number $d$ such that $d > 0$ and for every real number $z$ such that $z \neq 0$ and $|z| < d$ holds $\frac{\|R(z)\|}{|z|} < r$.

(18) Let $R$ be a function from $\mathbb{R}$ into $\mathbb{R}$. Then $R$ is rest-like if and only if for every real number $r$ such that $r > 0$ there exists a real number $d$ such that $d > 0$ and for every real number $z$ such that $z \neq 0$ and $|z| < d$ holds $\frac{|R(z)|}{|z|} < r$.

(19) Let $S$, $T$ be non trivial real normed spaces, $f$ be a partial function from $S$ to $T$, $p$, $q$ be points of $S$, and $M$ be a real number. Suppose that

(i) $[p, q] \subseteq \text{dom } f$,

(ii) for every point $x$ of $S$ such that $x \in [p, q]$ holds $f$ is continuous in $x$,

(iii) for every point $x$ of $S$ such that $x \in [p, q]$ holds $f$ is differentiable in $x$, and

(iv) for every point $x$ of $S$ such that $x \in [p, q]$ holds $\|f'(x)\| \leq M$.

Then $\|f_q - f_p\| \leq M \cdot \|q - p\|$. 
(20) Let $S$, $T$ be nontrivial real normed spaces, $f$ be a partial function from $S$ to $T$, $p$, $q$ be points of $S$, $M$ be a real number, and $L$ be a point of the real norm space of bounded linear operators from $S$ into $T$. Suppose that

(i) $[p, q] \subseteq \text{dom } f$,
(ii) for every point $x$ of $S$ such that $x \in [p, q]$ holds $f$ is continuous in $x$,
(iii) for every point $x$ of $S$ such that $x \in ]p, q[$ holds $f$ is differentiable in $x$, and
(iv) for every point $x$ of $S$ such that $x \in ]p, q[\, ]$ holds $\|f'(x) - L\| \leq M$.

Then $\|f_q - f_p - L(q - p)\| \leq M \cdot \|q - p\|$.

3. Partial Derivative of a Function of Several Variables

Let $G$ be a real norm space sequence and let $i$ be an element of $\text{dom } G$. The projection onto $i$ yielding a function from $\prod G$ into $G(i)$ is defined by:

(Def. 3) For every element $x$ of $\prod G$ holds $(\text{the projection onto } i)(x) = x(i)$.

Let $G$ be a real norm space sequence, let $i$ be an element of $\text{dom } G$, and let $x$ be an element of $\prod G$. The functor $\text{reproj}(i, x)$ yielding a function from $G(i)$ into $\prod G$ is defined by:

(Def. 4) For every element $r$ of $G(i)$ holds $\text{reproj}(i, x)(r) = x + \cdot (i, r)$.

Let $G$ be a nontrivial real norm space sequence and let $j$ be a set. Let us assume that $j \in \text{dom } G$. The functor $\text{modetrans}(G, j)$ yields an element of $\text{dom } G$ and is defined by:

(Def. 5) $\text{modetrans}(G, j) = j$.

Let $G$ be a nontrivial real norm space sequence, let $F$ be a non trivial real normed space, let $i$ be a set, let $f$ be a partial function from $\prod G$ to $F$, and let $x$ be an element of $\prod G$. We say that $f$ is partially differentiable in $x$ w.r.t. $i$ if and only if:

(Def. 6) $f \cdot \text{reproj}(\text{modetrans}(G, i), x)$ is differentiable in (the projection onto $\text{modetrans}(G, i))(x)$.

Let $G$ be a nontrivial real norm space sequence, let $F$ be a non trivial real normed space, let $i$ be a set, let $f$ be a partial function from $\prod G$ to $F$, and let $x$ be a point of $\prod G$. The functor $\text{partdiff}(f, x, i)$ yielding a point of the real norm space of bounded linear operators from $G(\text{modetrans}(G, i))$ into $F$ is defined as follows:

(Def. 7) $\text{partdiff}(f, x, i) = (f \cdot \text{reproj}(\text{modetrans}(G, i), x))'(\text{(the projection onto } \text{modetrans}(G, i))(x))$. 
4. Linearity of Partial Differential Operator

For simplicity, we adopt the following rules: $G$ denotes a nontrivial real norm space sequence, $F$ denotes a non trivial real normed space, $i$ denotes an element of $\text{dom} \, G$, $f$, $f_1$, $f_2$ denote partial functions from $\Pi G$ to $F$, $x$ denotes a point of $\Pi G$, and $X$ denotes a set.

Let $G$ be a nontrivial real norm space sequence, let $F$ be a non trivial real normed space, let $i$ be a set, let $f$ be a partial function from $\Pi G$ to $F$, and let $X$ be a set. We say that $f$ is partially differentiable on $X$ w.r.t. $i$ if and only if:

(Def. 8) $X \subseteq \text{dom} \, f$ and for every point $x$ of $\Pi G$ such that $x \in X$ holds $f|_X$ is partially differentiable in $x$ w.r.t. $i$.

Next we state several propositions:

(21) For every element $x_2$ of $G(i)$ holds $\|((\text{reproj}(i,0\Pi G))(x_2))\| = \|x_2\|$.

(22) Let $G$ be a nontrivial real norm space sequence, $i$ be an element of $\text{dom} \, G$, $x$ be a point of $\Pi G$, and $r$ be a point of $G(i)$. Then $(\text{reproj}(i,x))(r) - x = (\text{reproj}(i,0\Pi G))(r - (\text{the projection onto } i)(x))$ and $x - (\text{reproj}(i,x))(r) = (\text{reproj}(i,0\Pi G))(i) - (\text{the projection onto } i)(x) - r$.

(23) Let $G$ be a nontrivial real norm space sequence, $i$ be an element of $\text{dom} \, G$, $x$ be a point of $\Pi G$, and $Z$ be a subset of $\Pi G$. Suppose $Z$ is open and $x \in Z$. Then there exists a neighbourhood $N$ of $(\text{the projection onto } i)(x)$ such that for every point $z$ of $G(i)$ if $z \in N$, then $(\text{reproj}(i,x))(z) \in Z$.

(24) Let $G$ be a nontrivial real norm space sequence, $T$ be a non trivial real normed space, $i$ be a set, $f$ be a partial function from $\Pi G$ to $T$, and $Z$ be a subset of $\Pi G$. Suppose $Z$ is open. Then $f$ is partially differentiable on $Z$ w.r.t. $i$ if and only if $Z \subseteq \text{dom} \, f$ and for every point $x$ of $\Pi G$ such that $x \in Z$ holds $f$ is partially differentiable in $x$ w.r.t. $i$.

(25) For every set $i$ such that $i \in \text{dom} \, G$ and $f$ is partially differentiable on $X$ w.r.t. $i$ holds $X$ is a subset of $\Pi G$.

Let $G$ be a nontrivial real norm space sequence, let $S$ be a non trivial real normed space, and let $i$ be a set. Let us assume that $i \in \text{dom} \, G$. Let $f$ be a partial function from $\Pi G$ to $S$ and let $X$ be a set. Let us assume that $f$ is partially differentiable on $X$ w.r.t. $i$. The functor $f|_X$ yields a partial function from $\Pi G$ to the real norm space of bounded linear operators from $G(\text{modetrans}(G,i))$ into $S$ and is defined by:

(Def. 9) $\text{dom}(f|_X) = X$ and for every point $x$ of $\Pi G$ such that $x \in X$ holds $\langle f|_X \rangle_x = \text{partdiff}(f, x, i)$.

One can prove the following propositions:

(26) For every set $i$ such that $i \in \text{dom} \, G$ holds $(f_1 + f_2) \cdot \text{reproj}(\text{modetrans}(G, i), x) = f_1 \cdot \text{reproj}(\text{modetrans}(G, i), x) + f_2 \cdot$
reproj(modetrans(G, i), x) and \((f_1 - f_2) \cdot \text{reproj(modetrans}(G, i), x) = f_1 \cdot \text{reproj(modetrans}(G, i), x) - f_2 \cdot \text{reproj(modetrans}(G, i), x)\).

(27) For every set \(i\) such that \(i \in \text{dom } G\) holds \(r \cdot (f \cdot \text{reproj(modetrans}(G, i), x)) = (r \cdot f) \cdot \text{reproj(modetrans}(G, i), x)\).

(28) Let \(i\) be a set. Suppose \(i \in \text{dom } G \) and \(f_1\) is partially differentiable in \(x\) w.r.t. \(i\) and \(f_2\) is partially differentiable in \(x\) w.r.t. \(i\). Then \(f_1 + f_2\) is partially differentiable in \(x\) w.r.t. \(i\) and \(\text{partdiff}(f_1 + f_2, x, i) = \text{partdiff}(f_1, x, i) + \text{partdiff}(f_2, x, i)\).

(29) Let \(i\) be a set. Suppose \(i \in \text{dom } G \) and \(f_1\) is partially differentiable in \(x\) w.r.t. \(i\) and \(f_2\) is partially differentiable in \(x\) w.r.t. \(i\). Then \(f_1 - f_2\) is partially differentiable in \(x\) w.r.t. \(i\) and \(\text{partdiff}(f_1 - f_2, x, i) = \text{partdiff}(f_1, x, i) - \text{partdiff}(f_2, x, i)\).

(30) Let \(i\) be a set. Suppose \(i \in \text{dom } G \) and \(f\) is partially differentiable in \(x\) w.r.t. \(i\). Then \(r \cdot f\) is partially differentiable in \(x\) w.r.t. \(i\) and \(\text{partdiff}(r \cdot f, x, i) = r \cdot \text{partdiff}(f, x, i)\).

5. Continuous Differentiability of Partial Derivative

Next we state the proposition

(31) \(\| (\text{the projection onto } i) (x) \| \leq \| x \|\).

Let \(G\) be a nontrivial real norm space sequence. One can verify that every point of \(\prod G\) is \(\text{len } G\)-element.

We now state a number of propositions:

(32) Let \(G\) be a nontrivial real norm space sequence, \(T\) be a non trivial real normed space, \(i\) be a set, \(Z\) be a subset of \(\prod G\), and \(f\) be a partial function from \(\prod G\) to \(T\). Suppose \(Z\) is open. Then \(f\) is partially differentiable on \(Z\) w.r.t. \(i\) if and only if \(Z \subseteq \text{dom } f\) and for every point \(x\) of \(\prod G\) such that \(x \in Z\) holds \(f\) is partially differentiable in \(x\) w.r.t. \(i\).

(33) Let \(i, j\) be elements of \(\text{dom } G\), \(x\) be a point of \(G(i)\), and \(z\) be an element of \(\prod G\) such that \(z = (\text{reproj}(i, 0_{\prod G}))(x)\). Then

(i) if \(i = j\), then \(z(j) = x\), and

(ii) if \(i \neq j\), then \(z(j) = 0_{G(j)}\).

(34) For all points \(x, y\) of \(G(i)\) holds \((\text{reproj}(i, 0_{\prod G}))(x + y) = (\text{reproj}(i, 0_{\prod G}))(x) + (\text{reproj}(i, 0_{\prod G}))(y)\).

(35) Let \(x, y\) be points of \(\prod G\). Then \((\text{the projection onto } i)(x + y) = (\text{the projection onto } i)(x) + (\text{the projection onto } i)(y)\).

(36) For all points \(x, y\) of \(G(i)\) holds \((\text{reproj}(i, 0_{\prod G}))(x - y) = (\text{reproj}(i, 0_{\prod G}))(x) - (\text{reproj}(i, 0_{\prod G}))(y)\).
Let $x, y$ be points of $\prod G$. Then (the projection onto $i$)$(x - y) = (the$
projection onto $i$)$((x - y) - (the projection onto $i$)$y$).

For every point $x$ of $G(i)$ such that $x \neq 0_{G(i)}$ holds (reproj$((0, 0)_{G(i)}(x) \neq$
0$\prod G$).

For every point $x$ of $G(i)$ and for every element $a$ of $\mathbb{R}$ holds
(reproj$((i, 0)_{G(i)}(a \cdot x) = a \cdot (reproj((i, 0)_{G(i)}(x))$.

Let $x$ be a point of $\prod G$ and $a$ be an element of $\mathbb{R}$. Then (the projection
onto $i$)$((a \cdot x) = a \cdot (the projection onto $i$)$x$).

Let $G$ be a nontrivial real norm space sequence, $S$ be a non trivial
real normed space, $f$ be a partial function from $\prod G$ to $S$, $x$ be a
point of $\prod G$, and $i$ be a set. Suppose $f$ is differentiable in $x$. Then
$f$ is partially differentiable in $x$ w.r.t. $i$ and partdiff$((f, x, i) = f'(x) \cdot$
reproj$\modetrans(G, i), 0_{\prod G}$).

Let $S$ be a real normed space and $h, g$ be finite sequences of elements
of $S$. Suppose $\text{len } h = \text{len } g + 1$ and for every natural number $i$ such that
$i \in \text{dom } g$ holds $g_i = h_i - h_{i+1}$. Then $h_1 - h_{\text{len } h} = \sum g$.

Let $G$ be a nontrivial real norm space sequence, $x, y$ be elements of $\prod G$, and $Z$ be a set. Then $x + y|Z$ is an element of $\prod G$.

Let $G$ be a nontrivial real norm space sequence, $x, y$ be points of $\prod G$, $Z, x_0$ be elements of $\prod G$, and $X$ be a set. If $Z = 0_{\prod G}$ and $x_0 = x$ and
$y = Z + x_0|X$, then $\|y\| \leq \|x\|$.

Let $G$ be a nontrivial real norm space sequence, $S$ be a non trivial real
normed space, $f$ be a partial function from $\prod G$ to $S$, and $x, y$ be points
of $\prod G$. Then there exists a finite sequence $h$ of elements of $\prod G$ and there
exists a finite sequence $g$ of elements of $S$ and there exist elements $Z, y_0$ of $\prod G$ such that
$y_0 = y$ and $Z = 0_{\prod G}$ and len $h = \text{len } G + 1$ and len $g = \text{len } G$ and for every
natural number $i$ such that $i \in \text{dom } h$ holds $h_i = Z \cdot y_0 | \text{Seg}((\text{len } G + 1) - i)$
and for every natural number $i$ such that $i \in \text{dom } g$ holds $g_i = f_{x + h_i} - f_{x + h_{i+1}}$, and for every natural number $i$ and for every point $h_1$ of $\prod G$ such that
$i \in \text{dom } h$ and $h_i = h_1$ holds $\|h_1\| \leq \|y\|$ and $f_{x + y} - f_x = \sum g$.

Let $G$ be a nontrivial real norm space sequence, $i$ be an element of $\text{dom } G$, $x, y$ be points of $\prod G$, and $x_2$ be a point of $G(i)$. If $y = (\text{reproj}(i, x))(x_2)$, then (the projection onto $i$)$y = x_2$.

Let $G$ be a nontrivial real norm space sequence, $i$ be an element of $\text{dom } G$, $y$ be a point of $\prod G$, and $q$ be a point of $G(i)$. If $q = (\text{the projection onto}$
i$)(y)$, then $y = (\text{reproj}(i, y))(q)$.

Let $G$ be a nontrivial real norm space sequence, $i$ be an element of $\text{dom } G$, $x, y$ be points of $\prod G$, and $x_2$ be a point of $G(i)$. If $y = (\text{reproj}(i, x))(x_2)$, then $\text{reproj}(i, x) = \text{reproj}(i, y)$.
(49) Let $G$ be a nontrivial real norm space sequence, $i, j$ be elements of $\text{dom } G$, $x, y$ be points of $\prod G$, and $x_2$ be a point of $G(i)$. Suppose $y = \text{reproj}(i, x)(x_2)$ and $i \neq j$. Then $(\text{the projection onto } j)(x) = (\text{the projection onto } j)(y)$.

(50) Let $G$ be a nontrivial real norm space sequence, $F$ be a non trivial real normed space, $i$ be an element of $\text{dom } G$, $x$ be a point of $\prod G$, $x_2$ be a point of $G(i)$, $f$ be a partial function from $\prod G$ to $F$, and $g$ be a partial function from $G(i)$ to $F$. If $(\text{the projection onto } i)(x) = x_2$ and $g = f \cdot \text{reproj}(i, x)$, then $g(x_2) = \text{partdiff}(f, x, i)$.

(51) Let $G$ be a nontrivial real norm space sequence, $F$ be a non trivial real normed space, $f$ be a partial function from $\prod G$ to $F$, $x$ be a point of $\prod G$, $i$ be a set, $M$ be a real number, $L$ be a point of the real norm space of bounded linear operators from $G(\text{modetrans}(G, i))$ into $F$, and $p, q$ be points of $G(\text{modetrans}(G, i))$. Suppose that

(i) $i \in \text{dom } G$,

(ii) for every point $h$ of $G(\text{modetrans}(G, i))$ such that $h \in [p, q]$ holds $\|\text{partdiff}(f, \text{reproj}(\text{modetrans}(G, i), x))(h), i) - L\| \leq M$,

(iii) for every point $h$ of $G(\text{modetrans}(G, i))$ such that $h \in [p, q]$ holds $(\text{reproj}(\text{modetrans}(G, i), x))(h) \in \text{dom } f$, and

(iv) for every point $h$ of $G(\text{modetrans}(G, i))$ such that $h \in [p, q]$ holds $f$ is partially differentiable in $(\text{reproj}(\text{modetrans}(G, i), x))(h)$ w.r.t. $i$.

Then $\|f(\text{reproj}(\text{modetrans}(G, i), x)) - f(\text{reproj}(\text{modetrans}(G, i), x))(p) - L(q - p)\| \leq M \cdot \|q - p\|$.

(52) Let $G$ be a nontrivial real norm space sequence, $x, y, z, w$ be points of $\prod G$, $i$ be an element of $\text{dom } G$, $d$ be a real number, and $p, q, r$ be points of $G(i)$. Suppose $\|y - x\| < d$ and $\|z - x\| < d$ and $p = (\text{the projection onto } i)(y)$ and $z = (\text{reproj}(i, y))(q)$ and $r \in [p, q]$ and $w = (\text{reproj}(i, y))(r)$. Then $\|w - x\| < d$.

(53) Let $G$ be a nontrivial real norm space sequence, $S$ be a non trivial real normed space, $f$ be a partial function from $\prod G$ to $S$, $X$ be a subset of $\prod G$, $x, y, z$ be points of $\prod G$, $i$ be a set, $p, q$ be points of $G(\text{modetrans}(G, i))$, and $d, r$ be real numbers. Suppose that $i \in \text{dom } G$ and $X$ is open and $x \in X$ and $\|y - x\| < d$ and $\|z - x\| < d$ and $X \subseteq \text{dom } f$ and for every point $x$ of $\prod G$ such that $x \in X$ holds $f$ is partially differentiable in $x$ w.r.t. $i$ and for every point $z$ of $\prod G$ such that $\|z - x\| < d$ holds $\|\text{partdiff}(f, z, i) - \text{partdiff}(f, x, i)\| \leq r$ and $z = (\text{reproj}(\text{modetrans}(G, i), y))(p)$ and $q = (\text{the projection onto } \text{modetrans}(G, i))(y)$. Then $\|f_z - f_y - (\text{partdiff}(f, x, i))(p - q)\| \leq \|p - q\| \cdot r$.

(54) Let $G$ be a nontrivial real norm space sequence, $h$ be a finite sequence of elements of $\prod G$, $y, x$ be points of $\prod G$, $y_0, Z$ be elements of $\prod G$, and $j$ be an element of $\mathbb{N}$. Suppose $y = y_0$ and $Z = 0 \prod G$ and
len \( h \) = len \( G \) + 1 and \( 1 \leq j \leq \text{len} \( G \) \) and for every natural number \( i \) such that \( i \in \text{dom} \( h \) \) holds \( h_i = Z + y_0 | \text{Seg}((\text{len} \( G \) + 1) - j) \). Then
\[
x + h_j = (\text{reproj}(\text{modetrans}(G, (\text{len} \( G \) + 1) - j), x + h_{j+1}))((\text{the projection onto modetrans}(G, (\text{len} \( G \) + 1) - j))(x + y)).
\]

(55) Let \( G \) be a nontrivial real norm space sequence, \( h \) be a finite sequence of elements of \( \prod G \), \( y \) be points of \( \prod G \), \( y_0 \) \( Z \) be elements of \( \prod G \), and \( j \) be an element of \( \mathbb{N} \). Suppose \( y = y_0 \) and \( Z = 0 \prod G \) and \( \text{len} \( h \) = \text{len} \( G \) + 1 \) and \( 1 \leq j \leq \text{len} \( G \) \) and for every natural number \( i \) such that \( i \in \text{dom} \( h \) \) holds \( h_i = Z + y_0 | \text{Seg}((\text{len} \( G \) + 1) - i) \). Then (the projection onto modetrans\((G, (\text{len} \( G \) + 1) - j))\((x + y) - (\text{the projection onto modetrans}(G, (\text{len} \( G \) + 1) - j))(x + h_{j+1}) = (\text{the projection onto modetrans}\((G, (\text{len} \( G \) + 1) - j))(y)\).

(56) Let \( G \) be a nontrivial real norm space sequence, \( S \) be a non trivial real normed space, \( f \) be a partial function from \( \prod G \) to \( S \), \( X \) be a subset of \( \prod G \), and \( x \) be a point of \( \prod G \). Suppose that

(i) \( X \) is open,

(ii) \( x \in X \), and

(iii) for every set \( i \) such that \( i \in \text{dom} \( G \) \) holds \( f \) is partially differentiable on \( X \) w.r.t. \( i \) and \( f^i X \) is continuous on \( X \).

Then

(iv) \( f \) is differentiable in \( x \), and

(v) for every point \( h \) of \( \prod G \) there exists a finite sequence \( w \) of elements of \( S \) such that \( \text{dom} \( w \) = \text{dom} \( G \) \) and for every set \( i \) such that \( i \in \text{dom} \( G \) \) holds \( w(i) = (\text{partdiff}(f, x, i))((\text{the projection onto modetrans}(G, i))(h)) \) and \( f'(x)(h) = \sum w. \)

(57) Let \( G \) be a nontrivial real norm space sequence, \( F \) be a non trivial real normed space, \( f \) be a partial function from \( \prod G \) to \( F \), and \( X \) be a subset of \( \prod G \). Suppose \( X \) is open. Then for every set \( i \) such that \( i \in \text{dom} \( G \) \) holds \( f \) is partially differentiable on \( X \) w.r.t. \( i \) and \( f^i X \) is continuous on \( X \) if and only if \( f \) is differentiable on \( X \) and \( f^{i}|X \) is continuous on \( X \).

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Planes and Spheres as Topological Manifolds. Stereographic Projection

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Summary. The goal of this article is to show some examples of topological manifolds: planes and spheres in Euclidean space. In doing it, the article introduces the stereographic projection [25].

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The papers [29], [34], [9], [14], [40], [41], [11], [10], [4], [2], [18], [13], [31], [20], [21], [30], [32], [16], [17], [35], [26], [1], [22], [38], [36], [24], [19], [37], [28], [6], [15], [8], [27], [39], [3], [42], [12], [23], [7], [5], and [33] provide the notation and terminology for this paper.

1. Preliminaries

Let us observe that $\emptyset$ is $\emptyset$-valued and $\emptyset$ is onto.

Next we state three propositions:

1. For every function $f$ and for every set $Y$ holds $\text{dom}(Y|f) = f^{-1}(Y)$.
2. For every function $f$ and for all sets $Y_1, Y_2$ such that $Y_2 \subseteq Y_1$ holds $(Y_1|f)^{-1}(Y_2) = f^{-1}(Y_2)$.
3. Let $S, T$ be topological structures and $f$ be a function from $S$ into $T$. If $f$ is homeomorphism, then $f^{-1}$ is homeomorphism.

Let $S, T$ be topological structures. Let us note that the predicate $S$ and $T$ are homeomorphic is symmetric.

For simplicity, we use the following convention: $T_1, T_2, T_3$ denote topological spaces, $A_1$ denotes a subset of $T_1$, $A_2$ denotes a subset of $T_2$, and $A_3$ denotes a subset of $T_3$. 
Next we state several propositions:

(4) Let \( f \) be a function from \( T_1 \) into \( T_2 \). Suppose \( f \) is homeomorphism. Let \( g \) be a function from \( T_1 | f^{-1}(A_2) \) into \( T_2 | A_2 \). If \( g = A_2 | f \), then \( g \) is homeomorphism.

(5) For every function \( f \) from \( T_1 \) into \( T_2 \) such that \( f \) is homeomorphism holds \( f^{-1}(A_2) \) and \( A_2 \) are homeomorphic.

(6) If \( A_1 \) and \( A_2 \) are homeomorphic, then \( A_2 \) and \( A_1 \) are homeomorphic.

(7) If \( A_1 \) and \( A_2 \) are homeomorphic, then \( A_1 \) is empty iff \( A_2 \) is empty.

(8) If \( A_1 \) and \( A_2 \) are homeomorphic and \( A_2 \) and \( A_3 \) are homeomorphic, then \( A_1 \) and \( A_3 \) are homeomorphic.

(9) If \( T_1 \) is second-countable and \( T_1 \) and \( T_2 \) are homeomorphic, then \( T_2 \) is second-countable.

In the sequel \( n, k \) are natural numbers and \( M, N \) are non empty topological spaces.

The following propositions are true:

(10) If \( M \) is Hausdorff and \( M \) and \( N \) are homeomorphic, then \( N \) is Hausdorff.

(11) If \( M \) is \( n \)-locally Euclidean and \( M \) and \( N \) are homeomorphic, then \( N \) is \( n \)-locally Euclidean.

(12) If \( M \) is \( n \)-manifold and \( M \) and \( N \) are homeomorphic, then \( N \) is \( n \)-manifold.

(13) Let \( x_1, x_2 \) be finite sequences of elements of \( \mathbb{R} \) and \( i \) be an element of \( \mathbb{N} \). If \( i \in \text{dom}(x_1 \cdot x_2) \), then \( (x_1 \cdot x_2)(i) = (x_1)_i \cdot (x_2)_i \) and \( (x_1 \cdot x_2)_i = (x_1)_i \cdot (x_2)_i \).

(14) For all finite sequences \( x_1, x_2, y_1, y_2 \) of elements of \( \mathbb{R} \) such that \( \text{len} \ x_1 = \text{len} \ x_2 \) and \( \text{len} \ y_1 = \text{len} \ y_2 \) holds \( x_1 \triangleq y_1 \cdot x_2 \triangleq y_2 = (x_1 \cdot x_2) \triangleq (y_1 \cdot y_2) \).

(15) For all finite sequences \( x_1, x_2, y_1, y_2 \) of elements of \( \mathbb{R} \) such that \( \text{len} \ x_1 = \text{len} \ x_2 \) and \( \text{len} \ y_1 = \text{len} \ y_2 \) holds \( |(x_1 \triangleq y_1, x_2 \triangleq y_2)| = |(x_1, x_2)| + |(y_1, y_2)| \).

In the sequel \( p, q, p_1 \) are points of \( \mathcal{E}_T^n \) and \( r \) is a real number.

One can prove the following propositions:

(16) If \( k \in \text{Seg} \ n \), then \( (p_1 + p_2)(k) = p_1(k) + p_2(k) \).

(17) For every set \( X \) holds \( X \) is a linear combination of \( \mathbb{R}^{\text{Seg} \ n}_\mathbb{R} \) iff \( X \) is a linear combination of \( \mathcal{E}_T^n \).

(18) Let \( F \) be a finite sequence of elements of \( \mathcal{E}_T^n \), \( f_1 \) be a function from \( \mathcal{E}_T^n \) into \( \mathbb{R} \), \( F_1 \) be a finite sequence of elements of \( \mathbb{R}^{\text{Seg} \ n}_\mathbb{R} \), and \( f_2 \) be a function from \( \mathbb{R}^{\text{Seg} \ n}_\mathbb{R} \) into \( \mathbb{R} \). If \( f_1 = f_2 \) and \( F = F_1 \), then \( f_1 \cdot F = f_2 \cdot F_1 \).

(19) Let \( F \) be a finite sequence of elements of \( \mathcal{E}_T^n \) and \( F_1 \) be a finite sequence of elements of \( \mathbb{R}^{\text{Seg} \ n}_\mathbb{R} \). If \( F_1 = F \), then \( \sum F = \sum F_1 \).

(20) For every linear combination \( L_2 \) of \( \mathbb{R}^{\text{Seg} \ n}_\mathbb{R} \) and for every linear combination \( L_1 \) of \( \mathcal{E}_T^n \) such that \( L_1 = L_2 \) holds \( \sum L_1 = \sum L_2 \).
(21) Let $A_4$ be a subset of $\mathbb{R}^n_{\mathbb{R}}$ and $A_5$ be a subset of $\mathcal{E}^n_T$. Suppose $A_4 = A_5$. Then $A_4$ is linearly independent if and only if $A_5$ is linearly independent.

(22) For every subset $V$ of $\mathcal{E}^n_T$ such that $V = \mathbb{R}^n$ there exists a linear combination $l$ of $V$ such that $p = \sum l$.

(23) $\mathbb{R}^n$ is a basis of $\mathcal{E}^n_T$.

(24) Let $V$ be a subset of $\mathcal{E}^n_T$. Then $V \in \text{the topology of } \mathcal{E}^n_T$ if and only if for every $p$ such that $p \in V$ there exists $r > 0$ and $\text{Ball}(p, r) \subseteq V$.

Let $n$ be a natural number and let $p$ be a point of $\mathcal{E}^n_T$.

The functor $\text{InnerProduct}_p$ yields a function from $\mathcal{E}^n_T$ into $\mathbb{R}^1$ and is defined by:

(Def. 1) For every point $q$ of $\mathcal{E}^n_T$ holds $\text{InnerProduct}_p(q) = |(p, q)|$.

Let us consider $n$, $p$. Note that $\text{InnerProduct}_p$ is continuous.

2. Planes

Let us consider $n$ and let us consider $p$, $q$. The functor $\text{Plane}(p, q)$ yielding a subset of $\mathcal{E}^n_T$ is defined as follows:

(Def. 2) $\text{Plane}(p, q) = \{y; y \text{ ranges over points of } \mathcal{E}^n_T: |(p, y - q)| = 0\}$.

The following propositions are true:

(25) $\text{(trans}(p_1, \mathcal{E}^n_T)) \circ \text{Plane}(p, p_2) = \text{Plane}(p, p_1 + p_2)$.

(26) If $p \neq 0_{\mathcal{E}^n_T}$, then there exists a linearly independent subset $A$ of $\mathcal{E}^n_T$ such that $\overline{A} = n - 1$ and $\Omega_{\text{Lin}(A)} = \text{Plane}(p, 0_{\mathcal{E}^n_T})$.

(27) If $p_1 \neq 0_{\mathcal{E}^n_T}$ and $p_2 \neq 0_{\mathcal{E}^n_T}$, then there exists a function $R$ from $\mathcal{E}^n_T$ into $\mathcal{E}^n_T$ such that $R$ is homeomorphism and $R^\circ \text{Plane}(p_1, 0_{\mathcal{E}^n_T}) = \text{Plane}(p_2, 0_{\mathcal{E}^n_T})$.

Let us consider $n$ and let us consider $p$, $q$. The functor $\text{TPlane}(p, q)$ yields a non empty subspace of $\mathcal{E}^n_T$ and is defined by:

(Def. 3) $\text{TPlane}(p, q) = \mathcal{E}^n_T | \text{Plane}(p, q)$.

The following three propositions are true:

(28) The base finite sequence of $n + 1$ and $n + 1 = (0_{\mathcal{E}^n_T}) \sim (1)$.

(29) For all points $p, q$ of $\mathcal{E}^{n+1}_T$ such that $p \neq 0_{\mathcal{E}^{n+1}_T}$ holds $\mathcal{E}^n_T$ and $\text{TPlane}(p, q)$ are homeomorphic.

(30) For all points $p, q$ of $\mathcal{E}^{n+1}_T$ such that $p \neq 0_{\mathcal{E}^{n+1}_T}$ holds $\text{TPlane}(p, q)$ is $n$-manifold.
3. Spheres

Let us consider \( n \). The functor \( S^n \) yields a topological space and is defined by:

(Def. 4) \( S^n = \text{TopUnitCircle}(n + 1) \).

Let us consider \( n \). Note that \( S^n \) is non empty.

Let us consider \( n, p \) and let \( S \) be a subspace of \( E^n \). Let us assume that \( p \in \text{Sphere}((0_{E^n}), 1) \). The functor \( \sigma_{S,p} \) yielding a function from \( S \) into \( \text{TPlane}(p, 0_{E^n}) \) is defined as follows:

(Def. 5) For every \( q \) such that \( q \in S \) holds (
\[
\sigma_{S,p}(q) = \frac{1}{1 - |(q,p)|} \cdot (q - |(q,p)| \cdot p).
\]
)

Next we state the proposition

(31) For every subspace \( S \) of \( E^n \) such that \( \Omega_S = \text{Sphere}((0_{E^n}), 1) \setminus \{p\} \) and \( p \in \text{Sphere}((0_{E^n}), 1) \) holds \( \sigma_{S,p} \) is homeomorphism.

Let us consider \( n \). One can verify the following observations:

\( S^n \) is second-countable,
\( S^n \) is \( n \)-locally Euclidean, and
\( S^n \) is \( n \)-manifold.

References

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**Z-modules**

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**Summary.** In this article, we formalize \(\mathbb{Z}\)-module, that is a module over integer ring. \(\mathbb{Z}\)-module is necessary for lattice problems, LLL (Lenstra-Lenstra-Lovász) base reduction algorithm and cryptographic systems with lattices [11].

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The papers [10], [17], [18], [7], [2], [9], [14], [8], [6], [13], [5], [1], [15], [4], [3], [19], [16], and [12] provide the terminology and notation for this paper.

1. **Definition of \(\mathbb{Z}\)-module**

We introduce \(\mathbb{Z}\)-module structures which are extensions of additive loop structure and are systems

\[
\langle \text{a carrier, a zero, an addition, an external multiplication} \rangle,
\]

where the carrier is a set, the zero is an element of the carrier, the addition is a binary operation on the carrier, and the external multiplication is a function from \(\mathbb{Z} \times \) the carrier into the carrier.

Let us mention that there exists a \(\mathbb{Z}\)-module structure which is non empty.

Let \(V\) be a \(\mathbb{Z}\)-module structure. A vector of \(V\) is an element of \(V\).

In the sequel \(V\) denotes a non empty \(\mathbb{Z}\)-module structure and \(v\) denotes a vector of \(V\).

Let us consider \(V, v\) and let \(a\) be an integer number. The functor \(a \cdot v\) yields an element of \(V\) and is defined by:

\[(\text{Def. 1}) \quad a \cdot v = (\text{the external multiplication of } V)(a, v).\]
Let $Z_1$ be a non empty set, let $O$ be an element of $Z_1$, let $F$ be a binary operation on $Z_1$, and let $G$ be a function from $Z \times Z_1$ into $Z_1$. One can verify that $\langle \langle Z_1, O, F, G \rangle \rangle$ is non empty.

Let $I_1$ be a non empty $Z$-module structure. We say that $I_1$ is vector distributive if and only if:

(Def. 2) For every $a$ and for all vectors $v, w$ of $I_1$ holds $a \cdot (v + w) = a \cdot v + a \cdot w$.

We say that $I_1$ is scalar distributive if and only if:

(Def. 3) For all $a, b$ and for every vector $v$ of $I_1$ holds $(a + b) \cdot v = a \cdot v + b \cdot v$.

We say that $I_1$ is scalar associative if and only if:

(Def. 4) For all $a, b$ and for every vector $v$ of $I_1$ holds $(a \cdot b) \cdot v = a \cdot (b \cdot v)$.

We say that $I_1$ is scalar unital if and only if:

(Def. 5) For every vector $v$ of $I_1$ holds $1 \cdot v = v$.

The strict $Z$-module structure the trivial structure of $Z$-module is defined as follows:

(Def. 6) The trivial structure of $Z$-module $= \langle \langle 1, \text{op}_0, \text{op}_2, \pi_2(Z \times 1) \rangle \rangle$.

Let us observe that the trivial structure of $Z$-module is trivial and non empty. Let us observe that there exists a non empty $Z$-module structure which is strict, Abelian, add-associative, right zeroed, right complementable, scalar distributive, vector distributive, scalar associative, and scalar unital.

A $Z$-module is an Abelian add-associative right zeroed right complementable scalar distributive vector distributive scalar associative scalar unital non empty $Z$-module structure.

In the sequel $v, w$ denote vectors of $V$.

Let $I_1$ be a non empty $Z$-module structure. We say that $I_1$ inherits cancelable on multiplication if and only if:

(Def. 7) For every $a$ and for every vector $v$ of $I_1$ such that $a \cdot v = 0_{(I_1)}$ holds $a = 0$ or $v = 0_{(I_1)}$.

The following propositions are true:

1. If $a = 0$ or $v = 0_V$, then $a \cdot v = 0_V$.
2. $-v = (-1) \cdot v$.
3. If $V$ inherits cancelable on multiplication and $v = -v$, then $v = 0_V$.
4. If $V$ inherits cancelable on multiplication and $v + v = 0_V$, then $v = 0_V$.
5. $a \cdot -v = (-a) \cdot v$.
6. $a \cdot -v = -a \cdot v$.
7. $(-a) \cdot -v = a \cdot v$.
8. $a \cdot (v - w) = a \cdot v - a \cdot w$.
9. $(a - b) \cdot v = a \cdot v - b \cdot v$.
10. If $V$ inherits cancelable on multiplication and $a \neq 0$ and $a \cdot v = a \cdot w$, then $v = w$. 
(11) If \( V \) inherits cancelable on multiplication and \( v \neq 0_V \) and \( a \cdot v = b \cdot v \), then \( a = b \).

For simplicity, we follow the rules: \( V \) is a \( \mathbb{Z} \)-module, \( u, v, w \) are vectors of \( V, F, G, H, I \) are finite sequences of elements of \( V \), \( j, k, n \) are elements of \( \mathbb{N} \), and \( f_9 \) is a function from \( \mathbb{N} \) into the carrier of \( V \).

Next we state several propositions:

(12) If \( \text{len } F = \text{len } G \) and for all \( k, v \) such that \( k \in \text{dom } F \) and \( v = G(k) \) holds \( F(k) = a \cdot v \), then \( \sum F = a \cdot \sum G \).

(13) For every \( \mathbb{Z} \)-module \( V \) and for every integer \( a \) holds \( a \cdot \sum (\epsilon \text{(the carrier of } V)) = 0_V \).

(14) For every \( \mathbb{Z} \)-module \( V \) and for every integer \( a \) and for all vectors \( v, u \) of \( V \) holds \( a \cdot \sum \langle v, u \rangle = a \cdot v + a \cdot u \).

(15) For every \( \mathbb{Z} \)-module \( V \) and for every integer \( a \) and for all vectors \( v, u, w \) of \( V \) holds \( a \cdot \sum \langle v, u, w \rangle = a \cdot v + a \cdot u + a \cdot w \).

(16) \( (-a) \cdot v = -a \cdot v \).

(17) If \( \text{len } F = \text{len } G \) and for every \( k \) such that \( k \in \text{dom } F \) holds \( G(k) = a \cdot F(k) \), then \( \sum G = a \cdot \sum F \).

2. Submodules and Cosets of Submodules in \( \mathbb{Z} \)-module

We use the following convention: \( V, X \) are \( \mathbb{Z} \)-modules, \( V_1, V_2, V_3 \) are subsets of \( V \), and \( x \) is a set.

Let us consider \( V, V_1 \). We say that \( V_1 \) is linearly closed if and only if:

(Def. 8) For all \( v, u \) such that \( v, u \in V_1 \) holds \( v + u \in V_1 \) and for all \( a, v \) such that \( v \in V_1 \) holds \( a \cdot v \in V_1 \).

One can prove the following propositions:

(18) If \( V_1 \neq \emptyset \) and \( V_1 \) is linearly closed, then \( 0_V \in V_1 \).

(19) If \( V_1 \) is linearly closed, then for every \( v \) such that \( v \in V_1 \) holds \( -v \in V_1 \).

(20) If \( V_1 \) is linearly closed, then for all \( v, u \) such that \( v, u \in V_1 \) holds \( v - u \in V_1 \).

(21) If the carrier of \( V = V_1 \), then \( V_1 \) is linearly closed.

(22) If \( V_1 \) is linearly closed and \( V_2 \) is linearly closed and \( V_3 = \{ v + u : v \in V_1 \land u \in V_2 \} \), then \( V_3 \) is linearly closed.

Let us consider \( V \). Observe that \( \{0_V\} \) is linearly closed.

Let us consider \( V \). Note that there exists a subset of \( V \) which is linearly closed.

Let us consider \( V \) and let \( V_1, V_2 \) be linearly closed subsets of \( V \). Note that \( V_1 \cap V_2 \) is linearly closed.

Let us consider \( V \). A \( \mathbb{Z} \)-module is called a submodule of \( V \) if it satisfies the conditions (Def. 9).
(Def. 9)(i) The carrier of it \( \subseteq \) the carrier of \( V \),
(ii) \( 0_t = 0_V \),
(iii) the addition of it = (the addition of \( V \)) \| (the carrier of it), and
(iv) the external multiplication of it = (the external multiplication of \( V \))\| (Z \times \text{the carrier of it}).

In the sequel \( W_2 \) denotes a submodule of \( V \) and \( w, w_1, w_2 \) denote vectors of \( W \).

We now state a number of propositions:

(23) If \( x \in W_1 \) and \( W_1 \) is a submodule of \( W_2 \), then \( x \in W_2 \).
(24) If \( x \in W \), then \( x \in V \).
(25) \( w \) is a vector of \( V \).
(26) \( 0_W = 0_V \).
(27) \( 0_{(W_1)} = 0_{(W_2)} \).
(28) If \( w_1 = v \) and \( w_2 = u \), then \( w_1 + w_2 = v + u \).
(29) If \( w = v \), then \( a \cdot w = a \cdot v \).
(30) If \( w = v \), then \( -v = -w \).
(31) If \( w_1 = v \) and \( w_2 = u \), then \( w_1 - w_2 = v - u \).
(32) \( V \) is a submodule of \( V \).
(33) \( 0_V \in W \).
(34) \( 0_{(W_1)} \in W_2 \).
(35) \( 0_W \in V \).
(36) If \( u, v \in W \), then \( u + v \in W \).
(37) If \( v \in W \), then \( a \cdot v \in W \).
(38) If \( v \in W \), then \( -v \in W \).
(39) If \( u, v \in W \), then \( u - v \in W \).

In the sequel \( d_1 \) is an element of \( D \), \( A \) is a binary operation on \( D \), and \( M \) is a function from \( Z \times D \) into \( D \).

We now state several propositions:

(40) Suppose \( V_1 = D \) and \( d_1 = 0_V \) and \( A = (\text{the addition of } V) \| (V_1) \) and \( M = (\text{the external multiplication of } V)\| (Z \times V_1) \). Then \( \langle D, d_1, A, M \rangle \) is a submodule of \( V \).
(41) For all strict \( Z \)-modules \( V, X \) such that \( V \) is a submodule of \( X \) and \( X \) is a submodule of \( V \) holds \( V = X \).
(42) If \( V \) is a submodule of \( X \) and \( X \) is a submodule of \( Y \), then \( V \) is a submodule of \( Y \).
(43) If the carrier of \( W_1 \subseteq \) the carrier of \( W_2 \), then \( W_1 \) is a submodule of \( W_2 \).
(44) If for every \( v \) such that \( v \in W_1 \) holds \( v \in W_2 \), then \( W_1 \) is a submodule of \( W_2 \).
Let us consider $V$. Note that there exists a submodule of $V$ which is strict.

Next we state several propositions:

(45) For all strict submodules $W_1$, $W_2$ of $V$ such that the carrier of $W_1$ = the carrier of $W_2$ holds $W_1 = W_2$.

(46) For all strict submodules $W_1$, $W_2$ of $V$ such that for every $v$ holds $v \in W_1$ iff $v \in W_2$ holds $W_1 = W_2$.

(47) Let $V$ be a strict $\mathbb{Z}$-module and $W$ be a strict submodule of $V$. If the carrier of $W$ = the carrier of $V$, then $W = V$.

(48) Let $V$ be a strict $\mathbb{Z}$-module and $W$ be a strict submodule of $V$. If for every vector $v$ of $V$ holds $v \in W$ iff $v \in V$, then $W = V$.

(49) If the carrier of $W$ = $V_1$, then $V_1$ is linearly closed.

(50) If $V_1 \neq \emptyset$ and $V_1$ is linearly closed, then there exists a strict submodule $W$ of $V$ such that $V_1$ = the carrier of $W$.

Let us consider $V$. The functor $0_V$ yielding a strict submodule of $V$ is defined by:

(Def. 10) The carrier of $0_V = \{0_V\}$.

Let us consider $V$. The functor $\Omega_V$ yields a strict submodule of $V$ and is defined by:

(Def. 11) $\Omega_V = \text{the } \mathbb{Z}\text{-module structure of } V$.

We now state several propositions:

(51) $0_W = 0_V$.

(52) $0_{(W_1)} = 0_{(W_2)}$.

(53) $0_W$ is a submodule of $V$.

(54) $0_V$ is a submodule of $W$.

(55) $0_{(W_1)}$ is a submodule of $W_2$.

(56) Every strict $\mathbb{Z}$-module $V$ is a submodule of $\Omega_V$.

Let us consider $V$, $v$, $W$. The functor $v + W$ yields a subset of $V$ and is defined as follows:

(Def. 12) $v + W = \{v + u : u \in W\}$.

Let us consider $V$, $W$. A subset of $V$ is called a coset of $W$ if:

(Def. 13) There exists $v$ such that it = $v + W$.

In the sequel $B$, $C$ are cosets of $W$.

The following propositions are true:

(57) $0_V \in v + W$ iff $v \in W$.

(58) $v \in v + W$.

(59) $0_V + W$ = the carrier of $W$.

(60) $v + 0_V = \{v\}$.

(61) $v + \Omega_V$ = the carrier of $V$. 

(62) \(0_V \in v + W \) iff \( v + W = \text{the carrier of } W \).
(63) \( v \in W \) iff \( v + W = \text{the carrier of } W \).
(64) If \( v \in W \), then \( a \cdot v + W = \text{the carrier of } W \).
(65) \( u \in W \) iff \( v + W = v + u + W \).
(66) \( u \in W \) iff \( v + W = (v - u) + W \).
(67) \( v \in u + W \) iff \( u + W = v + W \).
(68) If \( u \in v_1 + W \) and \( u \in v_2 + W \), then \( v_1 + W = v_2 + W \).
(69) If \( v \in W \), then \( a \cdot v \in v + W \).
(70) \( u + v \in v + W \) iff \( u \in W \).
(71) \( v - u \in v + W \) iff \( u \in W \).
(72) \( u \in v + W \) iff there exists \( v_1 \) such that \( v_1 \in W \) and \( u = v + v_1 \).
(73) \( u \in v + W \) iff there exists \( v_1 \) such that \( v_1 \in W \) and \( u = v - v_1 \).
(74) There exists \( v \) such that \( v_1, v_2 \in v + W \) iff \( v_1 - v_2 \in W \).
(75) If \( v + W = u + W \), then there exists \( v_1 \) such that \( v_1 \in W \) and \( v + v_1 = u \).
(76) If \( v + W = u + W \), then there exists \( v_1 \) such that \( v_1 \in W \) and \( v - v_1 = u \).
(77) For all strict submodules \( W_1, W_2 \) of \( V \) such that \( v + W_1 = v + W_2 \) holds \( W_1 = W_2 \).
(78) For all strict submodules \( W_1, W_2 \) of \( V \) such that \( v + W_1 = u + W_2 \) holds \( W_1 = W_2 \).
(79) \( C \) is linearly closed iff \( C = \text{the carrier of } W \).
(80) For all strict submodules \( W_1, W_2 \) of \( V \) and for every coset \( C_1 \) of \( W_1 \) and for every coset \( C_2 \) of \( W_2 \) such that \( C_1 = C_2 \) holds \( W_1 = W_2 \).
(81) \( \{v\} \) is a coset of \( 0_V \).
(82) If \( V_1 \) is a coset of \( 0_V \), then there exists \( v \) such that \( V_1 = \{v\} \).
(83) The carrier of \( W \) is a coset of \( W \).
(84) The carrier of \( V \) is a coset of \( \Omega_V \).
(85) If \( V_1 \) is a coset of \( \Omega_V \), then \( V_1 = \text{the carrier of } V \).
(86) \( 0_V \in C \) iff \( C = \text{the carrier of } W \).
(87) \( u \in C \) iff \( C = u + W \).
(88) If \( u, v \in C \), then there exists \( v_1 \) such that \( v_1 \in W \) and \( u + v_1 = v \).
(89) If \( u, v \in C \), then there exists \( v_1 \) such that \( v_1 \in W \) and \( u - v_1 = v \).
(90) There exists \( C \) such that \( v_1, v_2 \in C \) iff \( v_1 - v_2 \in W \).
(91) If \( u \in B \) and \( u \in C \), then \( B = C \).
3. Operations on Submodules in \( \mathbb{Z} \)-module

For simplicity, we use the following convention: \( V \) is a \( \mathbb{Z} \)-module, \( W, W_1, W_2, W_3 \) are submodules of \( V, u, u_1, u_2, v, v_1, v_2 \) are vectors of \( V, a, a_1, a_2 \) are integer numbers, and \( X, Y, y, y_1, y_2 \) are sets.

Let us consider \( V, W_1, W_2 \). The functor \( W_1 + W_2 \) yielding a strict submodule of \( V \) is defined by:

(Def. 14) The carrier of \( W_1 + W_2 = \{ v + u : v \in W_1 \land u \in W_2 \} \).

Let us notice that the functor \( W_1 + W_2 \) is commutative.

Let us consider \( V, W_1, W_2 \). The functor \( W_1 \cap W_2 \) yields a strict submodule of \( V \) and is defined as follows:

(Def. 15) The carrier of \( W_1 \cap W_2 = (\text{the carrier of } W_1) \cap (\text{the carrier of } W_2) \).

Let us observe that the functor \( W_1 \cap W_2 \) is commutative.

We now state a number of propositions:

(92) \( x \in W_1 + W_2 \) iff there exist \( v_1, v_2 \) such that \( v_1 \in W_1 \) and \( v_2 \in W_2 \) and \( x = v_1 + v_2 \).

(93) If \( v \in W_1 \) or \( v \in W_2 \), then \( v \in W_1 + W_2 \).

(94) \( x \in W_1 \cap W_2 \) iff \( x \in W_1 \) and \( x \in W_2 \).

(95) For every strict submodule \( W \) of \( V \) holds \( W + W = W \).

(96) \( W_1 + (W_2 + W_3) = (W_1 + W_2) + W_3 \).

(97) \( W_1 \) is a submodule of \( W_1 + W_2 \).

(98) For every strict submodule \( W_2 \) of \( V \) holds \( W_1 \) is a submodule of \( W_2 \) iff \( W_1 + W_2 = W_2 \).

(99) For every strict submodule \( W \) of \( V \) holds \( 0_V + W = W \).

(100) \( 0_V + \Omega_V = \) the \( \mathbb{Z} \)-module structure of \( V \).

(101) \( \Omega_V + W = \) the \( \mathbb{Z} \)-module structure of \( V \).

(102) For every strict \( \mathbb{Z} \)-module \( V \) holds \( \Omega_V + \Omega_V = V \).

(103) For every strict submodule \( W \) of \( V \) holds \( W \cap W = W \).

(104) \( W_1 \cap (W_2 \cap W_3) = (W_1 \cap W_2) \cap W_3 \).

(105) \( W_1 \cap W_2 \) is a submodule of \( W_1 \).

(106) For every strict submodule \( W_1 \) of \( V \) holds \( W_1 \) is a submodule of \( W_2 \) iff \( W_1 \cap W_2 = W_1 \).

(107) \( 0_V \cap W = 0_V \).

(108) \( 0_V \cap \Omega_V = 0_V \).

(109) For every strict submodule \( W \) of \( V \) holds \( \Omega_V \cap W = W \).

(110) For every strict \( \mathbb{Z} \)-module \( V \) holds \( \Omega_V \cap \Omega_V = V \).

(111) \( W_1 \cap W_2 \) is a submodule of \( W_1 + W_2 \).

(112) For every strict submodule \( W_2 \) of \( V \) holds \( W_1 \cap W_2 + W_2 = W_2 \).
For every strict submodule $W_1$ of $V$ holds $W_1 \cap (W_1 + W_2) = W_1$.

$W_1 \cap W_2 + W_2 \cap W_3$ is a submodule of $W_2 \cap (W_1 + W_3)$.

If $W_1$ is a submodule of $W_2$, then $W_2 \cap (W_1 + W_3) = W_1 \cap W_2 + W_2 \cap W_3$.

$W_2 + W_1 \cap W_3$ is a submodule of $(W_1 + W_2) \cap (W_2 + W_3)$.

If $W_1$ is a submodule of $W_2$, then $W_2 + W_1 \cap W_3 = (W_1 + W_2) \cap (W_2 + W_3)$.

If $W_1$ is a strict submodule of $W_3$, then $W_1 + W_2 \cap W_3 = (W_1 + W_2) \cap W_3$.

For all strict submodules $W_1, W_2$ of $V$ holds $W_1 + W_2 = W_2$ iff $W_1 \cap W_2 = W_1$.

For all strict submodules $W_1, W_2, W_3$ of $V$ such that $W_1$ is a submodule of $W_2$ holds $W_1 + W_3$ is a submodule of $W_2 + W_3$.

There exists $W$ such that the carrier of $W = (\text{the carrier of } W_1) \cup (\text{the carrier of } W_2)$ if and only if $W_1$ is a submodule of $W_2$ or $W_2$ is a submodule of $W_1$.

Let us consider $V$. The functor $\text{Sub}(V)$ yields a set and is defined by:

- For every $x$ holds $x \in \text{Sub}(V)$ iff $x$ is a strict submodule of $V$.
- Let us consider $V$. One can verify that $\text{Sub}(V)$ is non empty.
- We now state the proposition

For every strict $\mathbb{Z}$-module $V$ holds $V \in \text{Sub}(V)$.

Let us consider $V, W_1, W_2$. We say that $V$ is the direct sum of $W_1$ and $W_2$ if and only if:

- The $\mathbb{Z}$-module structure of $V = W_1 + W_2$ and $W_1 \cap W_2 = 0_V$.
- Let $V$ be a $\mathbb{Z}$-module and let $W$ be a submodule of $V$. We say that $W$ has linear complement if and only if:

There exists a submodule $C$ of $V$ such that $V$ is the direct sum of $C$ and $W$.

Let $V$ be a $\mathbb{Z}$-module. Observe that there exists a submodule of $V$ which has linear complement.

Let $V$ be a $\mathbb{Z}$-module and let $W$ be a submodule of $V$. Let us assume that $W$ has linear complement. A submodule of $V$ is called a linear complement of $W$ if:

$V$ is the direct sum of it and $W$.

One can prove the following propositions:

- Let $V$ be a $\mathbb{Z}$-module and $W_1, W_2$ be submodules of $V$. Suppose $V$ is the direct sum of $W_1$ and $W_2$. Then $W_2$ is a linear complement of $W_1$.
- Let $V$ be a $\mathbb{Z}$-module, $W$ be a submodule of $V$ with linear complement, and $L$ be a linear complement of $W$. Then $V$ is the direct sum of $L$ and $W$ and the direct sum of $W$ and $L$. 
(125) Let $V$ be a $\mathbb{Z}$-module, $W$ be a submodule of $V$ with linear complement, and $L$ be a linear complement of $W$. Then $W + L = \text{the } Z$-module structure of $V$.

(126) Let $V$ be a $\mathbb{Z}$-module, $W$ be a submodule of $V$ with linear complement, and $L$ be a linear complement of $W$. Then $W \cap L = 0_V$.

(127) If $V$ is the direct sum of $W_1$ and $W_2$, then $V$ is the direct sum of $W_2$ and $W_1$.

(128) Let $V$ be a $\mathbb{Z}$-module, $W$ be a submodule of $V$ with linear complement, and $L$ be a linear complement of $W$. Then $W$ is a linear complement of $L$.

(129) Every $\mathbb{Z}$-module $V$ is the direct sum of $0_V$ and $\Omega_V$ and the direct sum of $\Omega_V$ and $0_V$.

(130) For every $\mathbb{Z}$-module $V$ holds $0_V$ is a linear complement of $\Omega_V$ and $\Omega_V$ is a linear complement of $0_V$.

In the sequel $C$ is a coset of $W$, $C_1$ is a coset of $W_1$, and $C_2$ is a coset of $W_2$.

Next we state several propositions:

(131) If $C_1$ meets $C_2$, then $C_1 \cap C_2$ is a coset of $W_1 \cap W_2$.

(132) Let $V$ be a $\mathbb{Z}$-module and $W_1, W_2$ be submodules of $V$. Then $V$ is the direct sum of $W_1$ and $W_2$ if and only if for every coset $C_1$ of $W_1$ and for every $v$ of $V$ such that $C_1 \cap C_2 = \{v\}$.

(133) Let $V$ be a $\mathbb{Z}$-module and $W_1, W_2$ be submodules of $V$. Then $W_1 + W_2 = \text{the } Z$-module structure of $V$ if and only if for every vector $v$ of $V$ there exist vectors $v_1, v_2$ of $V$ such that $v_1 \in W_1$ and $v_2 \in W_2$ and $v = v_1 + v_2$.

(134) If $V$ is the direct sum of $W_1$ and $W_2$ and $v_1 + v_2 = u_1 + u_2$ and $v_1, u_1 \in W_1$ and $v_2, u_2 \in W_2$, then $v_1 = u_1$ and $v_2 = u_2$.

(135) Suppose $V = W_1 + W_2$ and there exists $v$ such that for all $v_1, v_2, u_1, u_2$ such that $v_1 + v_2 = u_1 + u_2$ and $v_1, u_1 \in W_1$ and $v_2, u_2 \in W_2$ holds $v_1 = u_1$ and $v_2 = u_2$. Then $V$ is the direct sum of $W_1$ and $W_2$.

Let us consider $V, v, W_1, W_2$. Let us assume that $V$ is the direct sum of $W_1$ and $W_2$. The functor $v_{\langle W_1, W_2 \rangle}$ yields an element of $(\text{the carrier of } V) \times (\text{the carrier of } V)$ and is defined as follows:

$$v = (v_{\langle W_1, W_2 \rangle})_1 + (v_{\langle W_1, W_2 \rangle})_2 \text{ and } (v_{\langle W_1, W_2 \rangle})_1 \in W_1 \text{ and } (v_{\langle W_1, W_2 \rangle})_2 \in W_2.$$ 

Next we state several propositions:

(136) If $V$ is the direct sum of $W_1$ and $W_2$, then $(v_{\langle W_1, W_2 \rangle})_1 = (v_{\langle W_2, W_1 \rangle})_2$.

(137) If $V$ is the direct sum of $W_1$ and $W_2$, then $(v_{\langle W_1, W_2 \rangle})_2 = (v_{\langle W_2, W_1 \rangle})_1$.

(138) Let $V$ be a $\mathbb{Z}$-module, $W$ be a submodule of $V$ with linear complement, $L$ be a linear complement of $W$, $v$ be a vector of $V$, and $t$ be an element
of \((\text{the carrier of } V) \times (\text{the carrier of } V)\). If \(t_1 + t_2 = v\) and \(t_1 \in W\) and \(t_2 \in L\), then \(t = \langle v_{W,L} \rangle\).

(139) Let \(V\) be a \(Z\)-module, \(W\) be a submodule of \(V\) with linear complement, \(L\) be a linear complement of \(W\), and \(v\) be a vector of \(V\). Then \((v_{W,L})_1 + (v_{W,L})_2 = v\).

(140) Let \(V\) be a \(Z\)-module, \(W\) be a submodule of \(V\) with linear complement, \(L\) be a linear complement of \(W\), and \(v\) be a vector of \(V\). Then \((v_{W,L})_1 \in W\) and \((v_{W,L})_2 \in L\).

(141) Let \(V\) be a \(Z\)-module, \(W\) be a submodule of \(V\) with linear complement, \(L\) be a linear complement of \(W\), and \(v\) be a vector of \(V\). Then \((v_{W,L})_1 = (v_{L,W})_2\).

(142) Let \(V\) be a \(Z\)-module, \(W\) be a submodule of \(V\) with linear complement, \(L\) be a linear complement of \(W\), and \(v\) be a vector of \(V\). Then \((v_{W,L})_2 = (v_{L,W})_1\).

In the sequel \(A_1, A_2, B\) are elements of \(\text{Sub}(V)\).

Let us consider \(V\). The functor \(\text{SubJoin}_V\) yielding a binary operation on \(\text{Sub}(V)\) is defined by:

(Def. 21) For all \(A_1, A_2, W_1, W_2\) such that \(A_1 = W_1\) and \(A_2 = W_2\) holds \((\text{SubJoin}_V)(A_1, A_2) = W_1 + W_2\).

Let us consider \(V\). The functor \(\text{SubMeet}_V\) yields a binary operation on \(\text{Sub}(V)\) and is defined by:

(Def. 22) For all \(A_1, A_2, W_1, W_2\) such that \(A_1 = W_1\) and \(A_2 = W_2\) holds \((\text{SubMeet}_V)(A_1, A_2) = W_1 \cap W_2\).

One can prove the following proposition

(143) \((\text{Sub}(V), \text{SubJoin}_V, \text{SubMeet}_V)\) is a lattice.

Let us consider \(V\). Note that \((\text{Sub}(V), \text{SubJoin}_V, \text{SubMeet}_V)\) is lattice-like.

We now state several propositions:

(144) For every \(Z\)-module \(V\) holds \((\text{Sub}(V), \text{SubJoin}_V, \text{SubMeet}_V)\) is lower-bounded.

(145) For every \(Z\)-module \(V\) holds \((\text{Sub}(V), \text{SubJoin}_V, \text{SubMeet}_V)\) is upper-bounded.

(146) For every \(Z\)-module \(V\) holds \((\text{Sub}(V), \text{SubJoin}_V, \text{SubMeet}_V)\) is a bound lattice.

(147) For every \(Z\)-module \(V\) holds \((\text{Sub}(V), \text{SubJoin}_V, \text{SubMeet}_V)\) is modular.

(148) Let \(V\) be a \(Z\)-module and \(W_1, W_2, W_3\) be strict submodules of \(V\). If \(W_1\) is a submodule of \(W_2\), then \(W_1 \cap W_3\) is a submodule of \(W_2 \cap W_3\).

(149) Let \(V\) be a \(Z\)-module and \(W\) be a strict submodule of \(V\). Suppose that for every vector \(v\) of \(V\) holds \(v \in W\). Then \(W\) is the \(Z\)-module structure
There exists $C$ such that $v \in C$. 

4. Transformation of Abelian Group to $\mathbb{Z}$-module

Let $A_3$ be a non empty additive loop structure. The left integer multiplication of $A_3$ yielding a function from $\mathbb{Z} \times$ the carrier of $A_3$ into the carrier of $A_3$ is defined by the condition (Def. 23).

(Def. 23) Let $i$ be an element of $\mathbb{Z}$ and $a$ be an element of $A_3$. Then

(i) if $i \geq 0$, then (the left integer multiplication of $A_3)(i,a) = (\text{Nat-mult-left } A_3)(i,a)$, and

(ii) if $i < 0$, then (the left integer multiplication of $A_3)(i,a) = (\text{Nat-mult-left } A_3)(-i,-a)$.

The following propositions are true:

(151) Let $R$ be a non empty additive loop structure, $a$ be an element of $R$, $i$ be an element of $\mathbb{Z}$, and $i_1$ be an element of $\mathbb{N}$. If $i = i_1$, then (the left integer multiplication of $R)(i,a) = i_1 \cdot a$.

(152) Let $R$ be a non empty additive loop structure, $a$ be an element of $R$, and $i$ be an element of $\mathbb{Z}$. If $i = 0$, then (the left integer multiplication of $R)(i,a) = 0_R$.

(153) Let $R$ be an add-associative right zeroed right complementable non empty additive loop structure and $i$ be an element of $\mathbb{N}$. Then $(\text{Nat-mult-left } R)(i,0_R) = 0_R$.

(154) Let $R$ be an add-associative right zeroed right complementable non empty additive loop structure and $i$ be an element of $\mathbb{Z}$. Then (the left integer multiplication of $R)(i,0_R) = 0_R$.

(155) Let $R$ be a right zeroed non empty additive loop structure, $a$ be an element of $R$, and $i$ be an element of $\mathbb{Z}$. If $i = 1$, then (the left integer multiplication of $R)(i,a) = a$.

(156) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a$ be an element of $R$, and $i, j, k$ be elements of $\mathbb{N}$. If $i \leq j$ and $k = j - i$, then $(\text{Nat-mult-left } R)(k,a) = (\text{Nat-mult-left } R)(j,a) - (\text{Nat-mult-left } R)(i,a)$.

(157) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a$ be an element of $R$, and $i$ be an element of $\mathbb{N}$. Then $-(\text{Nat-mult-left } R)(i,a) = (\text{Nat-mult-left } R)(i,-a)$.

(158) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a$ be an element of $R$, and $i, j$ be elements of $\mathbb{Z}$. Suppose $i \in \mathbb{N}$ and $j \notin \mathbb{N}$. Then (the left integer multiply-
cation of $R)(i + j, a) = (\text{the left integer multiplication of } R)(i, a) + (\text{the left integer multiplication of } R)(j, a)$.

(159) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a$ be an element of $R$, and $i, j$ be elements of $\mathbb{Z}$. Then $(\text{the left integer multiplication of } R)(i + j, a) = (\text{the left integer multiplication of } R)(i, a) + (\text{the left integer multiplication of } R)(j, a)$.

(160) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a, b$ be elements of $R$, and $i$ be an element of $\mathbb{N}$. Then $(\text{Nat-mult-left } R)(i, a + b) = (\text{Nat-mult-left } R)(i, a) + (\text{Nat-mult-left } R)(i, b)$.

(161) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a, b$ be elements of $R$, and $i$ be an element of $\mathbb{Z}$. Then $(\text{the left integer multiplication of } R)(i, a + b) = (\text{the left integer multiplication of } R)(i, a) + (\text{the left integer multiplication of } R)(i, b)$.

(162) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a$ be an element of $R$, and $i, j$ be elements of $\mathbb{N}$. Then $(\text{Nat-mult-left } R)(i \cdot j, a) = (\text{Nat-mult-left } R)(i, (\text{Nat-mult-left } R)(j, a))$.

(163) Let $R$ be an Abelian right zeroed add-associative right complementable non empty additive loop structure, $a$ be an element of $R$, and $i, j$ be elements of $\mathbb{Z}$. Then $(\text{the left integer multiplication of } R)(i \cdot j, a) = (\text{the left integer multiplication of } R)(i, (\text{the left integer multiplication of } R)(j, a))$.

(164) Let $A_3$ be a non empty Abelian add-associative right zeroed right complementable additive loop structure. Then $\langle\text{the carrier of } A_3, \text{the zero of } A_3, \text{the addition of } A_3, \text{the left integer multiplication of } A_3\rangle$ is a $\mathbb{Z}$-module.

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Morphology for Image Processing. Part I

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Summary. In this article we defined mathematical morphology image processing with set operations. First, we defined Minkowski set operations and proved their properties. Next, we defined basic image processing, dilation and erosion proving basic fact about them [5], [8].

The terminology and notation used in this paper have been introduced in the following papers: [10], [7], [1], [2], [6], [9], [4], and [3].

1. Minkowski Set Operations

Let $E$ be a non empty RLS structure. A binary image of $E$ is a subset of $E$. In the sequel $E$ denotes a real linear space and $A$ denotes a binary image of $E$.

Let $E$ be a real linear space and let $A$, $B$ be binary images of $E$. The functor $A \oplus B$ yielding a binary image of $E$ is defined as follows:

\[(\text{Def. 1}) \quad A \oplus B = \{ z \in E : \land_{b: \text{element of } E} (b \in B \Rightarrow z - b \in A) \}.\]

Let $a$ be a real number, let $E$ be a real linear space, and let $A$ be a subset of $E$. We introduce $a \cdot A$ as a synonym of $a \circ A$. The following propositions are true:

1. Let $E$ be a real linear space and $A$, $B$ be subsets of $E$. If $B = \emptyset$, then $A \oplus B = B$ and $B \oplus A = B$ and $A \circ B = \text{the carrier of } E$.
2. For every real linear space $E$ and for all subsets $A$, $B$ of $E$ such that $A \neq \emptyset$ and $B = \emptyset$ holds $B \circ A = B$.

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Let $E$ be a real linear space and $A, B$ be subsets of $E$. If $B$ = the carrier of $E$ and $A \neq \emptyset$, then $A \oplus B = B$ and $B \oplus A = B$.

For every real linear space $E$ and for all subsets $A, B$ of $E$ such that $B$ = the carrier of $E$ holds $B \ominus A = B$.

$A \oplus B = \bigcup\{b + A; b \text{ ranges over elements of } E; b \in B\}$.

Let $E$ be a non empty RLS structure. A binary image family of $E$ is a family of subsets of the carrier of $E$.

We follow the rules: $F, G$ are binary image families of $E$ and $A, B, C$ are non empty binary images of $E$. We now state four propositions:

(6) $A \ominus B = \bigcap\{b + A; b \text{ ranges over elements of } E; b \in B\}$.

(7) $A \ominus B = \{v \in E: (v + (-1) \cdot B) \cap A \neq \emptyset\}$.

(8) $A \ominus B = \{v \in E: v + (-1) \cdot B \subseteq A\}$.

(9) ((The carrier of $E) \setminus A \ominus B = \text{(the carrier of } E) \setminus A \oplus B \text{ and ((the carrier of } E) \setminus A \ominus B) \setminus B = \text{(the carrier of } E) \setminus A \oplus B$.)

Let $E$ be a non empty Abelian additive loop structure and let $A, B$ be subsets of $E$. Let us note that the functor $A \oplus B$ is commutative.

One can prove the following propositions:

(10) For every non empty add-associative additive loop structure $E$ and for all subsets $A, B, C$ of $E$ holds $(A + B) + C = A + (B + C)$.

(11) $(A \oplus B) \oplus C = A \oplus (B \oplus C)$.

(12) $\bigcup F \ominus B = \bigcup\{X \ominus B; X \text{ ranges over binary images of } E; X \in F\}$.

(13) $A \ominus \bigcup F = \bigcup\{A \ominus X; X \text{ ranges over binary images of } E; X \in F\}$.

(14) $\bigcap F \ominus B \subseteq \bigcap\{X \ominus B; X \text{ ranges over binary images of } E; X \in F\}$.

(15) $A \ominus \bigcap F \subseteq \bigcap\{A \ominus X; X \text{ ranges over binary images of } E; X \in F\}$.

(16) For every non empty additive loop structure $E$ and for all subsets $A, B, C$ of $E$ such that $B \subseteq C$ holds $A + B \subseteq A + C$.

(17) $(v + A) \ominus B = A \ominus (v + B)$ and $(v + A) \ominus B = v + A \ominus B$.

(18) $\bigcap F \ominus B = \bigcap\{X \ominus B; X \text{ ranges over binary images of } E; X \in F\}$.

(19) $\bigcap\{B \ominus X; X \text{ ranges over binary images of } E; X \in F\} \subseteq B \ominus \bigcap F$.

(20) $\bigcup\{X \ominus B; X \text{ ranges over binary images of } E; X \in F\} \subseteq \bigcup F \ominus B$.

(21) If $F \neq \emptyset$, then $B \ominus \bigcup F = \bigcap\{B \ominus X; X \text{ ranges over binary images of } E; X \in F\}$.

(22) If $A \subseteq B$, then $A \ominus C \subseteq B \ominus C$.

(23) If $A \subseteq B$, then $C \ominus B \subseteq C \ominus A$.

(24) $(v + A) \ominus B = A \ominus (v + B)$ and $(v + A) \ominus B = v + A \ominus B$.

(25) $A \ominus B \ominus C = A \ominus (B \ominus C)$. 
2. Dilation and Erosion

Let $E$ be a real linear space and let $B$ be a binary image of $E$. The functor dilation $B$ yields a function from $2^{\text{carrier of } E}$ into $2^{\text{carrier of } E}$ and is defined as follows:

(Def. 2) For every binary image $A$ of $E$ holds $(\text{dilation } B)(A) = A \oplus B$.

Let $E$ be a real linear space and let $B$ be a binary image of $E$. The functor erosion $B$ yields a function from $2^{\text{carrier of } E}$ into $2^{\text{carrier of } E}$ and is defined by:

(Def. 3) For every binary image $A$ of $E$ holds $(\text{erosion } B)(A) = A \ominus B$.

The following propositions are true:

(26) $(\text{dilation } B)(\bigcup F) = \bigcup \{ (\text{dilation } B)(X); X \text{ ranges over binary images of } E: X \in F \}$.

(27) If $A \subseteq B$, then $(\text{dilation } C)(A) \subseteq (\text{dilation } C)(B)$.

(28) $(\text{dilation } C)(v + A) = v + (\text{dilation } C)(A)$.

(29) $(\text{erosion } B)(\bigcap F) = \bigcap \{ (\text{erosion } B)(X); X \text{ ranges over binary images of } E: X \in F \}$.

(30) If $A \subseteq B$, then $(\text{erosion } C)(A) \subseteq (\text{erosion } C)(B)$.

(31) $(\text{erosion } C)(v + A) = v + (\text{erosion } C)(A)$.

(32) $(\text{dilation } C)((\text{the carrier of } E) \setminus A) = (\text{the carrier of } E) \setminus (\text{erosion } C)(A)$ and $(\text{erosion } C)((\text{the carrier of } E) \setminus A) = (\text{the carrier of } E) \setminus (\text{dilation } C)(A)$.

(33) $(\text{dilation } C)((\text{dilation } B)(A)) = (\text{dilation}(\text{dilation } C)(B))(A)$ and $(\text{erosion } C)((\text{erosion } B)(A)) = (\text{erosion}(\text{dilation } C)(B))(A)$.

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The Differentiable Functions from $\mathbb{R}$ into $\mathbb{R}^n$

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Summary. In control engineering, differentiable partial functions from $\mathbb{R}$ into $\mathbb{R}^n$ play a very important role. In this article, we formalized basic properties of such functions.

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The notation and terminology used in this paper are introduced in the following articles: [25], [26], [6], [2], [27], [8], [7], [24], [1], [4], [3], [5], [9], [22], [20], [28], [21], [10], [23], [17], [13], [11], [12], [15], [19], [18], [16], and [14].

Let us observe that there exists a sequence of real numbers which is convergent to 0 and non-zero.

For simplicity, we adopt the following convention: $x_0$, $r$ denote real numbers, $i$, $m$ denote elements of $\mathbb{N}$, $n$ denotes a non-empty element of $\mathbb{N}$, $Y$ denotes a subset of $\mathbb{R}$, $Z$ denotes an open subset of $\mathbb{R}$, and $f_1$, $f_2$ denote partial functions from $\mathbb{R}$ to $\mathbb{R}^n$.

The following proposition is true

(1) For all partial functions $f_1$, $f_2$ from $\mathbb{R}$ to $\mathbb{R}^m$ holds $f_1 - f_2 = f_1 + -f_2$. 

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Let \( n \) be a non empty element of \( \mathbb{N} \), let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \), and let \( x \) be a real number. We say that \( f \) is differentiable in \( x \) if and only if:

(Def. 1) There exists a partial function \( g \) from \( \mathbb{R} \) to \( \langle \mathbb{E}^n, \| \cdot \| \rangle \) such that \( f = g \) and \( g \) is differentiable in \( x \).

One can prove the following proposition

(2) Let \( n \) be a non empty element of \( \mathbb{N} \), \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \), \( h \) be a partial function from \( \mathbb{R} \) to \( \langle \mathbb{E}^n, \| \cdot \| \rangle \), and \( x \) be a real number. Suppose \( h = f \). Then \( f \) is differentiable in \( x \) if and only if \( h \) is differentiable in \( x \).

Let \( n \) be a non empty element of \( \mathbb{N} \), let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \), and let \( x \) be a real number. The functor \( f'(x) \) yields an element of \( \mathbb{R}^n \) and is defined as follows:

(Def. 2) There exists a partial function \( g \) from \( \mathbb{R} \) to \( \langle \mathbb{E}^n, \| \cdot \| \rangle \) such that \( f = g \) and \( f'(x) = g'(x) \).

One can prove the following proposition

(3) Let \( n \) be a non empty element of \( \mathbb{N} \), \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \), \( h \) be a partial function from \( \mathbb{R} \) to \( \langle \mathbb{E}^n, \| \cdot \| \rangle \), and \( x \) be a real number. If \( h = f \), then \( f'(x) = h'(x) \).

Let us consider \( n, f, X \). We say that \( f \) is differentiable on \( X \) if and only if:

(Def. 3) \( X \subseteq \text{dom } f \) and for every \( x \) such that \( x \in X \) holds \( f|_X \) is differentiable in \( x \).

The following propositions are true:

(4) If \( f \) is differentiable on \( X \), then \( X \) is a subset of \( \mathbb{R} \).

(5) \( f \) is differentiable on \( Z \) iff \( Z \subseteq \text{dom } f \) and for every \( x \) such that \( x \in Z \) holds \( f \) is differentiable in \( x \).

(6) If \( f \) is differentiable on \( Y \), then \( Y \) is open.

Let us consider \( n, f, X \). Let us assume that \( f \) is differentiable on \( X \). The functor \( f'_{|X} \) yields a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \) and is defined by:

(Def. 4) \( \text{dom}(f'_{|X}) = X \) and for every \( x \) such that \( x \in X \) holds \( f'_{|X}(x) = f'(x) \).

One can prove the following propositions:

(7) Suppose \( Z \subseteq \text{dom } f \) and there exists an element \( r \) of \( \mathbb{R}^n \) such that \( \text{rng } f = \{ r \} \). Then \( f \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( (f'_{|Z})_x = (0, \ldots, 0) \).

(8) Let \( x_0 \) be a real number, \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \), \( g \) be a partial function from \( \mathbb{R} \) to \( \langle \mathbb{E}^n, \| \cdot \| \rangle \), and \( N \) be a neighbourhood of \( x_0 \). Suppose \( f = g \) and \( f \) is differentiable in \( x_0 \) and \( N \subseteq \text{dom } f \). Let given \( h, c \).
Suppose \( \text{rng } c = \{x_0\} \) and \( \text{rng}(h+c) \subseteq N \). Then \( h^{-1} \cdot ((g_*(h+c)) - (g_*c)) \) is convergent and \( f'(x_0) = \lim(h^{-1} \cdot ((g_*(h+c)) - (g_*c))). \)

(9) If \( f \) is differentiable in \( x_0 \), then \( r \cdot f \) is differentiable in \( x_0 \) and \( (r \cdot f)'(x_0) = r \cdot f'(x_0) \).

(10) If \( f \) is differentiable in \( x_0 \), then \( -f \) is differentiable in \( x_0 \) and \( (-f)'(x_0) = -f'(x_0) \).

(11) If \( f_1 \) is differentiable in \( x_0 \) and \( f_2 \) is differentiable in \( x_0 \), then \( f_1 + f_2 \) is differentiable in \( x_0 \) and \( (f_1 + f_2)'(x_0) = f_1'(x_0) + f_2'(x_0) \).

(12) If \( f_1 \) is differentiable in \( x_0 \) and \( f_2 \) is differentiable in \( x_0 \), then \( f_1 - f_2 \) is differentiable in \( x_0 \) and \( (f_1 - f_2)'(x_0) = f_1'(x_0) - f_2'(x_0) \).

(13) Suppose \( Z \subseteq \text{dom } f \) and \( f \) is differentiable on \( Z \). Then \( r \cdot f \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( (r \cdot f)'_Z(x) = r \cdot f'(x) \).

(14) If \( Z \subseteq \text{dom } f \) and \( f \) is differentiable on \( Z \), then \( -f \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( (-f)'_Z(x) = -f'(x) \).

(15) Suppose \( Z \subseteq \text{dom}(f_1 + f_2) \) and \( f_1 \) is differentiable on \( Z \) and \( f_2 \) is differentiable on \( Z \). Then \( f_1 + f_2 \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( (f_1 + f_2)'_Z(x) = f_1'(x) + f_2'(x) \).

(16) Suppose \( Z \subseteq \text{dom}(f_1 - f_2) \) and \( f_1 \) is differentiable on \( Z \) and \( f_2 \) is differentiable on \( Z \). Then \( f_1 - f_2 \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( (f_1 - f_2)'_Z(x) = f_1'(x) - f_2'(x) \).

(17) If \( Z \subseteq \text{dom } f \) and \( f|Z \) is constant, then \( f \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( f'|_Z(x) = \begin{pmatrix} 0, \ldots, 0 \end{pmatrix} \).

(18) Let \( r, p \) be elements of \( \mathbb{R}^n \). Suppose \( Z \subseteq \text{dom } f \) and for every \( x \) such that \( x \in Z \) holds \( f_x = x \cdot r + p \). Then \( f \) is differentiable on \( Z \) and for every \( x \) such that \( x \in Z \) holds \( f'|_Z(x) = r \).

(19) For every real number \( x_0 \) such that \( f \) is differentiable in \( x_0 \) holds \( f \) is continuous in \( x_0 \).

(20) If \( f \) is differentiable on \( X \), then \( f|X \) is continuous.

(21) If \( f \) is differentiable on \( X \) and \( Z \subseteq X \), then \( f \) is differentiable on \( Z \).

Let \( n \) be a non empty element of \( \mathbb{N} \) and let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R}^n \). We say that \( f \) is differentiable if and only if:

(Def. 5) \( f \) is differentiable on \( \text{dom } f \).

Let us consider \( n \). One can check that \( \mathbb{R} \mapsto \begin{pmatrix} 0, \ldots, 0 \end{pmatrix} \) is differentiable.

Let us consider \( n \). Note that there exists a function from \( \mathbb{R} \) into \( \mathbb{R}^n \) which is differentiable.

One can prove the following proposition

(22) For every differentiable partial function \( f \) from \( \mathbb{R} \) to \( \mathbb{R}^n \) such that \( Z \subseteq \text{dom } f \) holds \( f \) is differentiable on \( Z \).
In the sequel $G_1$, $R$ are rests of $\langle \mathcal{E}^n, \| \cdot \| \rangle$ and $D_1$, $L$ are linears of $\langle \mathcal{E}^n, \| \cdot \| \rangle$.

Next we state a number of propositions:

(23) Let $R$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $R$ is total. Then $R$ is rest-like if and only if for every real number $r$ such that $r > 0$ there exists a real number $d$ such that $d > 0$ and for every real number $z$ such that $z \neq 0$ and $|z| < d$ holds $|z|^{-1} \cdot \| R_z \| < r$.

(24) Let $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$ and $x_0$ be a real number. Suppose $1 \leq i \leq n$ and $g$ is differentiable in $x_0$. Then $\text{Proj}(i, n) \cdot g$ is differentiable in $x_0$ and $(\text{Proj}(i, n))(g'(x_0)) = (\text{Proj}(i, n) \cdot g)'(x_0)$.

(25) Let $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$ and $x_0$ be a real number. Then $g$ is differentiable in $x_0$ if and only if for every element $i$ of $\mathbb{N}$ such that $1 \leq i \leq n$ holds $\text{Proj}(i, n) \cdot g$ is differentiable in $x_0$.

(26) Let $f$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$ and $x_0$ be a real number. Suppose $1 \leq i \leq n$ and $f$ is differentiable in $x_0$. Then $\text{Proj}(i, n) \cdot f$ is differentiable in $x_0$ and $(\text{Proj}(i, n))(f'(x_0)) = (\text{Proj}(i, n) \cdot f)'(x_0)$.

(27) Let $f$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$ and $x_0$ be a real number. Then $f$ is differentiable in $x_0$ if and only if for every element $i$ of $\mathbb{N}$ such that $1 \leq i \leq n$ holds $\text{Proj}(i, n) \cdot f$ is differentiable in $x_0$.

(28) Let $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $1 \leq i \leq n$ and $g$ is differentiable on $X$. Then $\text{Proj}(i, n) \cdot g$ is differentiable on $X$ and $\text{Proj}(i, n) \cdot g|_X = (\text{Proj}(i, n) \cdot g)'|_X$.

(29) Let $f$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$. Suppose $1 \leq i \leq n$ and $f$ is differentiable on $X$. Then $\text{Proj}(i, n) \cdot f$ is differentiable on $X$ and $\text{Proj}(i, n) \cdot f|_X = (\text{Proj}(i, n) \cdot f)'|_X$.

(30) Let $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Then $g$ is differentiable on $X$ if and only if for every element $i$ of $\mathbb{N}$ such that $1 \leq i \leq n$ holds $\text{Proj}(i, n) \cdot g$ is differentiable on $X$.

(31) Let $f$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$. Then $f$ is differentiable on $X$ if and only if for every element $i$ of $\mathbb{N}$ such that $1 \leq i \leq n$ holds $\text{Proj}(i, n) \cdot f$ is differentiable on $X$.

(32) For every function $J$ from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ into $\mathbb{R}$ and for every point $x_0$ of $\langle \mathcal{E}^1, \| \cdot \| \rangle$ such that $J = \text{proj}(1, 1)$ holds $J$ is continuous in $x_0$.

(33) For every function $I$ from $\mathbb{R}$ into $\langle \mathcal{E}^1, \| \cdot \| \rangle$ such that $I = \text{proj}(1, 1)^{-1}$ holds $I$ is continuous in $x_0$.

(34) Let $S$, $T$ be real normed spaces, $f_1$ be a partial function from $S$ to $\mathbb{R}$, $f_2$ be a partial function from $\mathbb{R}$ to $T$, and $x_0$ be a point of $S$. Suppose $x_0 \in \text{dom}(f_2 \cdot f_1)$ and $f_1$ is continuous in $x_0$ and $f_2$ is continuous in $(f_1)|_{x_0}$. Then $f_2 \cdot f_1$ is continuous in $x_0$.

(35) Let $J$ be a function from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ into $\mathbb{R}$, $x_0$ be a point of $\langle \mathcal{E}^1, \| \cdot \| \rangle$, $y_0$ be an element of $\mathbb{R}$, $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$, and $f$
be a partial function from \(\langle E^1, \| \cdot \| \rangle\) to \(\langle E^n, \| \cdot \| \rangle\). Suppose \(J = \text{proj}(1,1)\) and \(x_0 \in \text{dom} \ g \) and \(y_0 \in \text{dom} \ g \) and \(x_0 = \langle y_0 \rangle\) and \(f = g \cdot J\). Then \(f\) is continuous in \(x_0\) if and only if \(g\) is continuous in \(y_0\).

(36) Let \(I\) be a function from \(\mathbb{R}\) into \(\langle E^1, \| \cdot \| \rangle\), \(x_0\) be a point of \(\langle E^1, \| \cdot \| \rangle\), \(y_0\) be an element of \(\mathbb{R}\), \(g\) be a partial function from \(\mathbb{R}\) to \(\langle E^n, \| \cdot \| \rangle\), and \(f\) be a partial function from \(\langle E^1, \| \cdot \| \rangle\) to \(\langle E^n, \| \cdot \| \rangle\). Suppose \(I = \text{proj}(1,1)^{-1}\) and \(x_0 \in \text{dom} \ f\) and \(y_0 \in \text{dom} \ g\) and \(x_0 = \langle y_0 \rangle\) and \(f \cdot I = g\). Then \(f\) is continuous in \(x_0\) if and only if \(g\) is continuous in \(y_0\).

(37) For every function \(I\) from \(\mathbb{R}\) into \(\langle E^1, \| \cdot \| \rangle\) such that \(I = \text{proj}(1,1)^{-1}\) holds \(I\) is differentiable in \(x_0\) and \(I'(x_0) = (1)\).

Let \(n\) be a non empty element of \(\mathbb{N}\), let \(f\) be a partial function from \(\langle E^n, \| \cdot \| \rangle\) to \(\mathbb{R}\), and let \(x\) be a point of \(\langle E^n, \| \cdot \| \rangle\). We say that \(f\) is differentiable in \(x\) if and only if the condition (Def. 6) is satisfied.

(Def. 6) There exists a partial function \(g\) from \(\mathcal{R}^n\) to \(\mathbb{R}\) and there exists an element \(y\) of \(\mathcal{R}^n\) such that \(f = g\) and \(x = y\) and \(g\) is differentiable in \(y\).

Let \(n\) be a non empty element of \(\mathbb{N}\), let \(f\) be a partial function from \(\langle E^n, \| \cdot \| \rangle\) to \(\mathbb{R}\), and let \(x\) be a point of \(\langle E^n, \| \cdot \| \rangle\). The functor \(f'(x)\) yields a function from \(\langle E^n, \| \cdot \| \rangle\) into \(\mathbb{R}\) and is defined by:

(Def. 7) There exists a partial function \(g\) from \(\mathcal{R}^n\) to \(\mathbb{R}\) and there exists an element \(y\) of \(\mathcal{R}^n\) such that \(f = g\) and \(x = y\) and \(f'(x) = g'(y)\).

We now state several propositions:

(38) Let \(J\) be a function from \(\mathcal{R}^1\) into \(\mathbb{R}\) and \(x_0\) be an element of \(\mathcal{R}^1\). If \(J = \text{proj}(1,1)\), then \(J\) is differentiable in \(x_0\) and \(J'(x_0) = J\).

(39) Let \(J\) be a function from \(\langle E^1, \| \cdot \| \rangle\) into \(\mathbb{R}\) and \(x_0\) be a point of \(\langle E^1, \| \cdot \| \rangle\). If \(J = \text{proj}(1,1)\), then \(J\) is differentiable in \(x_0\) and \(J'(x_0) = J\).

(40) Let \(I\) be a function from \(\mathbb{R}\) into \(\langle E^1, \| \cdot \| \rangle\). Suppose \(I = \text{proj}(1,1)^{-1}\). Then

(i) for every rest \(R\) of \(\langle E^1, \| \cdot \| \rangle\), \(\langle E^n, \| \cdot \| \rangle\) holds \(R \cdot I\) is a rest of \(\langle E^n, \| \cdot \| \rangle\), and

(ii) for every linear operator \(L\) from \(\langle E^1, \| \cdot \| \rangle\) into \(\langle E^n, \| \cdot \| \rangle\) holds \(L \cdot I\) is a linear of \(\langle E^n, \| \cdot \| \rangle\).

(41) Let \(J\) be a function from \(\langle E^1, \| \cdot \| \rangle\) into \(\mathbb{R}\). Suppose \(J = \text{proj}(1,1)\). Then

(i) for every rest \(R\) of \(\langle E^n, \| \cdot \| \rangle\) holds \(R \cdot J\) is a rest of \(\langle E^1, \| \cdot \| \rangle\), \(\langle E^n, \| \cdot \| \rangle\), and

(ii) for every linear \(L\) of \(\langle E^n, \| \cdot \| \rangle\) holds \(L \cdot J\) is a bounded linear operator from \(\langle E^1, \| \cdot \| \rangle\) into \(\langle E^n, \| \cdot \| \rangle\).

(42) Let \(I\) be a function from \(\mathbb{R}\) into \(\langle E^1, \| \cdot \| \rangle\), \(x_0\) be a point of \(\langle E^1, \| \cdot \| \rangle\), \(y_0\) be an element of \(\mathbb{R}\), \(g\) be a partial function from \(\mathbb{R}\) to \(\langle E^n, \| \cdot \| \rangle\), and \(f\) be a partial function from \(\langle E^1, \| \cdot \| \rangle\) to \(\langle E^n, \| \cdot \| \rangle\). Suppose \(I = \text{proj}(1,1)^{-1}\) and \(x_0 \in \text{dom} \ f\) and \(y_0 \in \text{dom} \ g\) and \(x_0 = \langle y_0 \rangle\) and \(f \cdot I = g\) and \(f\) is
differentiable in $x_0$. Then $g$ is differentiable in $y_0$ and $g'(y_0) = f'(x_0)(1)$
and for every element $r$ of $\mathbb{R}$ holds $f'(x_0)(r) = r \cdot g'(y_0)$. 

(43) Let $I$ be a function from $\mathbb{R}$ into $\langle \mathcal{E}^1, \| \cdot \| \rangle$, $x_0$ be a point of $\langle \mathcal{E}^1, \| \cdot \| \rangle$, $y_0$ be an element of $\mathbb{R}$, $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$, and $f$ be a partial function from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $I = \text{proj}(1,1)^{-1}$
and $x_0 \in \text{dom } f$ and $y_0 \in \text{dom } g$ and $x_0 = \langle y_0 \rangle$ and $f \cdot I = g$. Then $f$ is
differentiable in $x_0$ if and only if $g$ is differentiable in $y_0$.

(44) Let $J$ be a function from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ into $\mathbb{R}$, $x_0$ be a point of $\langle \mathcal{E}^1, \| \cdot \| \rangle$,
y_0 be an element of $\mathbb{R}$, $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$, and $f$
be a partial function from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $J = \text{proj}(1,1)$
and $x_0 \in \text{dom } f$ and $y_0 \in \text{dom } g$ and $x_0 = \langle y_0 \rangle$ and $f = g \cdot J$. Then $f$ is
differentiable in $x_0$ if and only if $g$ is differentiable in $y_0$.

(45) Let $J$ be a function from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ into $\mathbb{R}$, $x_0$ be a point of $\langle \mathcal{E}^1, \| \cdot \| \rangle$,
y_0 be an element of $\mathbb{R}$, $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$, and $f$
be a partial function from $\langle \mathcal{E}^1, \| \cdot \| \rangle$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $J = \text{proj}(1,1)$
and $x_0 \in \text{dom } f$ and $y_0 \in \text{dom } g$ and $x_0 = \langle y_0 \rangle$ and $f = g \cdot J$ and $g$
is differentiable in $y_0$. Then $f$ is differentiable in $x_0$ and $g'(y_0) = f'(x_0)(1)$
and for every element $r$ of $\mathbb{R}$ holds $f'(x_0)(r) = r \cdot g'(y_0)$.

(46) Let $R$ be a rest of $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $R_0 = 0_{\langle \mathcal{E}^n, \| \cdot \| \rangle}$. Let $e$ be a real
number. Suppose $e > 0$. Then there exists a real number $d$ such that $d > 0$ and for every real number $h$
such that $|h| < d$ holds $\| R_h \| \leq e \cdot |h|$

In the sequel $m$, $n$ denote non empty elements of $\mathbb{N}$.

One can prove the following propositions:

(47) For every rest $R$ of $\langle \mathcal{E}^n, \| \cdot \| \rangle$ and for every bounded linear operator $L$
from $\langle \mathcal{E}^n, \| \cdot \| \rangle$ into $\langle \mathcal{E}^m, \| \cdot \| \rangle$ holds $L \cdot R$ is a rest of $\langle \mathcal{E}^m, \| \cdot \| \rangle$.

(48) Let $R_1$ be a rest of $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $(R_1)_0 = 0_{\langle \mathcal{E}^n, \| \cdot \| \rangle}$. Let $R_2$ be a rest
of $\langle \mathcal{E}^n, \| \cdot \| \rangle$, $\langle \mathcal{E}^m, \| \cdot \| \rangle$. Suppose $(R_2)_0 = 0_{\langle \mathcal{E}^m, \| \cdot \| \rangle}$. Let $L$ be a linear
of $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Then $R_2 \cdot (L + R_1)$ is a rest of $\langle \mathcal{E}^m, \| \cdot \| \rangle$.

(49) Let $R_1$ be a rest of $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $(R_1)_0 = 0_{\langle \mathcal{E}^n, \| \cdot \| \rangle}$. Let $R_2$ be a rest
of $\langle \mathcal{E}^n, \| \cdot \| \rangle$, $\langle \mathcal{E}^m, \| \cdot \| \rangle$. Suppose $(R_2)_0 = 0_{\langle \mathcal{E}^m, \| \cdot \| \rangle}$. Let $L_1$ be a linear
of $\langle \mathcal{E}^n, \| \cdot \| \rangle$ and $L_2$ be a bounded linear operator from $\langle \mathcal{E}^n, \| \cdot \| \rangle$
to $\langle \mathcal{E}^m, \| \cdot \| \rangle$. Then $L_2 \cdot R_1 + R_2 \cdot (L_1 + R_1)$ is a rest of $\langle \mathcal{E}^m, \| \cdot \| \rangle$.

(50) Let $x_0$ be an element of $\mathbb{R}$ and $g$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$.
Suppose $g$ is differentiable in $x_0$. Let $f$ be a partial function from $\langle \mathcal{E}^n, \| \cdot \| \rangle$
to $\langle \mathcal{E}^m, \| \cdot \| \rangle$. Suppose $f$ is differentiable in $g(x_0)$. Then $f \cdot g$ is differentiable
in $x_0$ and $(f \cdot g)'(x_0) = f'(g(x_0))g'(x_0)$.

References


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Some Basic Properties of Some Special Matrices. Part III

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Summary. This article describes definitions of subsymmetric matrix, anti-subsymmetric matrix, central symmetric matrix, symmetry circulant matrix and their basic properties.

The notation and terminology used here have been introduced in the following papers: [7], [9], [13], [6], [14], [1], [3], [18], [17], [4], [2], [8], [11], [12], [16], [15], [5], and [10].

1. Basic Properties of Subordinate Symmetric Matrices

For simplicity, we use the following convention: \( n \) denotes a natural number, \( K \) denotes a field, \( a, b \) denote elements of \( K \), \( p, q \) denote finite sequences of elements of \( K \), and \( M_1, M_2 \) denote square matrices over \( K \) of dimension \( n \). We say that \( M \) is subsymmetric if and only if:

(Def. 1) For all natural numbers \( i, j, k, l \) such that \( \langle i, j \rangle \in \) the indices of \( M \) and \( k = (n + 1) - j \) and \( l = (n + 1) - i \) holds \( M_{i,j} = M_{k,l} \).

Let us consider \( n, K, a \). Note that \( (a)^{n\times n} \) is subsymmetric.

Let us consider \( n, K \). Observe that there exists a square matrix over \( K \) of dimension \( n \) which is subsymmetric.

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Let us consider $n$, $K$ and let $M$ be a subsymmetric square matrix over $K$ of dimension $n$. Note that $-M$ is subsymmetric.

Let us consider $n$, $K$ and let $M_1$, $M_2$ be subsymmetric square matrices over $K$ of dimension $n$. One can check that $M_1 + M_2$ is subsymmetric.

Let us consider $n$, $K$, $a$ and let $M$ be a subsymmetric square matrix over $K$ of dimension $n$. Note that $a \cdot M$ is subsymmetric.

Let us consider $n$, $K$ and let $M_1$, $M_2$ be subsymmetric square matrices over $K$ of dimension $n$. One can verify that $M_1 - M_2$ is subsymmetric.

Let us consider $n$, $K$ and let $M$ be a subsymmetric square matrix over $K$ of dimension $n$. Observe that $M^T$ is subsymmetric.

Let us consider $n$, $K$. Observe that every square matrix over $K$ of dimension $n$ which is line circulant is also subsymmetric and every square matrix over $K$ of dimension $n$ which is column circulant is also subsymmetric.

Let $K$ be a field, let $n$ be a natural number, and let $M$ be a square matrix over $K$ of dimension $n$. We say that $M$ is anti-subsymmetric if and only if:

(Def. 2) For all natural numbers $i$, $j$, $k$, $l$ such that $\langle i, j \rangle \in$ the indices of $M$ and $k = (n + 1) - j$ and $l = (n + 1) - i$ holds $M_{i,j} = -M_{k,l}$.

Let us consider $n$, $K$. One can verify that there exists a square matrix over $K$ of dimension $n$ which is anti-subsymmetric.

The following proposition is true

(1) Let $K$ be a Fanoian field, $n$, $i$, $j$, $k$, $l$ be natural numbers, and $M_1$ be a square matrix over $K$ of dimension $n$. Suppose $\langle i, j \rangle \in$ the indices of $M_1$ and $i + j = n + 1$ and $k = (n + 1) - j$ and $l = (n + 1) - i$ and $M_1$ is anti-subsymmetric. Then $(M_1)_{i,j} = 0_K$.

Let us consider $n$, $K$ and let $M$ be an anti-subsymmetric square matrix over $K$ of dimension $n$. Note that $-M$ is anti-subsymmetric.

Let us consider $n$, $K$ and let $M_1$, $M_2$ be anti-subsymmetric square matrices over $K$ of dimension $n$. Observe that $M_1 + M_2$ is anti-subsymmetric.

Let us consider $n$, $K$, $a$ and let $M$ be an anti-subsymmetric square matrix over $K$ of dimension $n$. One can verify that $a \cdot M$ is anti-subsymmetric.

Let us consider $n$, $K$ and let $M_1$, $M_2$ be anti-subsymmetric square matrices over $K$ of dimension $n$. One can check that $M_1 - M_2$ is anti-subsymmetric.

Let us consider $n$, $K$ and let $M$ be an anti-subsymmetric square matrix over $K$ of dimension $n$. One can verify that $M^T$ is anti-subsymmetric.

2. Basic Properties of Central Symmetric Matrices

Let $K$ be a field, let $n$ be a natural number, and let $M$ be a square matrix over $K$ of dimension $n$. We say that $M$ is central symmetric if and only if:

(Def. 3) For all natural numbers $i$, $j$, $k$, $l$ such that $\langle i, j \rangle \in$ the indices of $M$ and $k = (n + 1) - i$ and $l = (n + 1) - j$ holds $M_{i,j} = M_{k,l}$. 


Let us consider $n$, $K$, $a$. Note that $(a)^{n\times n}$ is central symmetric.

Let us consider $n$, $K$. One can verify that there exists a square matrix over $K$ of dimension $n$ which is central symmetric.

Let us consider $n$, $K$ and let $M$ be a central symmetric square matrix over $K$ of dimension $n$. One can verify that $-M$ is central symmetric.

Let us consider $n$, $K$ and let $M_1$, $M_2$ be central symmetric square matrices over $K$ of dimension $n$. One can verify that $M_1 + M_2$ is central symmetric.

Let us consider $n$, $K$, $a$ and let $M$ be a central symmetric square matrix over $K$ of dimension $n$. Note that $a \cdot M$ is central symmetric.

Let us consider $n$, $K$ and let $M_1$, $M_2$ be central symmetric square matrices over $K$ of dimension $n$. Observe that $M_1 - M_2$ is central symmetric.

Let us consider $n$, $K$. Observe that $M^T$ is central symmetric.

Let us consider $n$, $K$. Note that every square matrix over $K$ of dimension $n$ which is symmetric and subsymmetric is also central symmetric.

3. Basic Properties of Symmetric Circulant Matrices

Let $K$ be a set, let $M$ be a matrix over $K$, and let $p$ be a finite sequence. We say that $M$ is symmetry circulant about $p$ if and only if the conditions (Def. 4) are satisfied.

(Def. 4)(i) $\text{len } p = \text{width } M$,

(ii) for all natural numbers $i$, $j$ such that $\langle i, j \rangle \in$ the indices of $M$ and $i + j \neq \text{len } p + 1$ holds $M_{i,j} = p(((i + j) - 1) \mod \text{len } p)$, and

(iii) for all natural numbers $i$, $j$ such that $\langle i, j \rangle \in$ the indices of $M$ and $i + j = \text{len } p + 1$ holds $M_{i,j} = p(\text{len } p)$.

The following propositions are true:

(2) $(a)^{n\times n}$ is symmetry circulant about $n \mapsto a$.

(3) If $M_1$ is symmetry circulant about $p$, then $a \cdot M_1$ is symmetry circulant about $a \cdot p$.

(4) If $M_1$ is symmetry circulant about $p$, then $-M_1$ is symmetry circulant about $-p$.

(5) If $M_1$ is symmetry circulant about $p$ and $M_2$ is symmetry circulant about $q$, then $M_1 + M_2$ is symmetry circulant about $p + q$.

Let $K$ be a set and let $M$ be a matrix over $K$. We say that $M$ is symmetry circulant if and only if:

(Def. 5) There exists a finite sequence $p$ of elements of $K$ such that $\text{len } p = \text{width } M$ and $M$ is symmetry circulant about $p$.

Let $K$ be a non empty set and let $p$ be a finite sequence of elements of $K$. We say that $p$ is first symmetry of circulant if and only if:
There exists a square matrix over $K$ of dimension $\text{len}\ p$ which is symmetry circulant about $p$.

Let $K$ be a non empty set and let $p$ be a finite sequence of elements of $K$. Let us assume that $p$ is first symmetry of circulant. The functor $\text{SCirc}\ p$ yielding a square matrix over $K$ of dimension $\text{len}\ p$ is defined as follows:

Let us consider $n, K, a$. Note that $(a)^{n\times n}$ is symmetry circulant.

Let us consider $n, K$. Note that there exists a square matrix over $K$ of dimension $n$ which is symmetry circulant.

In the sequel $D$ is a non empty set, $t$ is a finite sequence of elements of $D$, and $A$ is a square matrix over $D$ of dimension $n$.

We now state the proposition

(6) Let $p$ be a finite sequence of elements of $D$. Suppose $0 < n$ and $A$ is symmetry circulant about $p$. Then $A^T$ is symmetry circulant about $p$.

Let us consider $n, K, a$ and let $M_1$ be a symmetry circulant square matrix over $K$ of dimension $n$. Note that $a \cdot M_1$ is symmetry circulant.

Let us consider $n, K$ and let $M_1, M_2$ be symmetry circulant square matrices over $K$ of dimension $n$. Note that $M_1 + M_2$ is symmetry circulant.

Let us consider $n, K$ and let $M_1$ be a symmetry circulant square matrix over $K$ of dimension $n$. Note that $-M_1$ is symmetry circulant.

Let us consider $n, K$ and let $M_1, M_2$ be symmetry circulant square matrices over $K$ of dimension $n$. Observe that $M_1 - M_2$ is symmetry circulant.

The following propositions are true:

(7) If $A$ is symmetry circulant and $n > 0$, then $A^T$ is symmetry circulant.

(8) If $p$ is first symmetry of circulant, then $-p$ is first symmetry of circulant.

(9) If $p$ is first symmetry of circulant, then $\text{SCirc}(-p) = -\text{SCirc}\ p$.

(10) Suppose $p$ is first symmetry of circulant and $q$ is first symmetry of circulant and $\text{len} p = \text{len} q$. Then $p + q$ is first symmetry of circulant.

(11) If $\text{len} p = \text{len} q$ and $p$ is first symmetry of circulant and $q$ is first symmetry of circulant, then $\text{SCirc}(p + q) = \text{SCirc}\ p + \text{SCirc}\ q$.

(12) If $p$ is first symmetry of circulant, then $a \cdot p$ is first symmetry of circulant.

(13) If $p$ is first symmetry of circulant, then $\text{SCirc}(a \cdot p) = a \cdot \text{SCirc}\ p$.

(14) If $p$ is first symmetry of circulant, then $a \cdot \text{SCirc}\ p + b \cdot \text{SCirc}\ p = \text{SCirc}(a + b) \cdot p$.

(15) If $p$ is first symmetry of circulant and $q$ is first symmetry of circulant and $\text{len} p = \text{len} q$, then $a \cdot \text{SCirc}\ p + a \cdot \text{SCirc}\ q = \text{SCirc}(a \cdot (p + q))$.

(16) Suppose $p$ is first symmetry of circulant and $q$ is first symmetry of circulant and $\text{len} p = \text{len} q$. Then $a \cdot \text{SCirc}\ p + b \cdot \text{SCirc}\ q = \text{SCirc}(a \cdot p + b \cdot q)$.

(17) If $M_1$ is symmetry circulant, then $M_1^T = M_1$. 


Let us consider $n, K$. Note that every square matrix over $K$ of dimension $n$ which is symmetry circulant is also symmetric.

REFERENCES


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Riemann Integral of Functions from $\mathbb{R}$ into $n$-dimensional Real Normed Space

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Summary. In this article, we define the Riemann integral on functions $\mathbb{R}$ into $n$-dimensional real normed space and prove the linearity of this operator. As a result, the Riemann integration can be applied to the wider range. Our method refers to the [21].

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The terminology and notation used in this paper have been introduced in the following papers: [23], [24], [6], [2], [25], [8], [7], [1], [4], [3], [5], [20], [10], [14], [12], [13], [18], [22], [19], [26], [9], [11], [15], [17], and [16].

1. ON THE FUNCTIONS FROM $\mathbb{R}$ INTO $n$-DIMENSIONAL REAL SPACE

For simplicity, we adopt the following convention: $X$ denotes a set, $n$ denotes an element of $\mathbb{N}$, $a$, $b$, $c$, $d$, $e$, $r$, $x_0$ denote real numbers, $A$ denotes a non empty closed-interval subset of $\mathbb{R}$, $f$, $g$, $h$ denote partial functions from $\mathbb{R}$ to $\mathbb{R}^n$, and $E$ denotes an element of $\mathbb{R}^n$. We now state a number of propositions:

1. If $a \leq c \leq b$, then $c \in [a, b]$ and $[a, c] \subseteq [a, b]$ and $[c, b] \subseteq [a, b]$. 

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(2) If \( a \leq c \leq d \leq b \) and \([a, b] \subseteq X\), then \([c, d] \subseteq X\).
(3) If \( a \leq b \) and \( c, d \in [a, b] \) and \([a, b] \subseteq X\), then \([\min(c, d), \max(c, d)] \subseteq X\).
(4) If \( a \leq c \leq d \leq b \) and \([a, b] \subseteq \text{dom} \, f \) and \([a, b] \subseteq \text{dom} \, g\), then \([c, d] \subseteq \text{dom}(f + g)\).
(5) If \( a \leq c \leq d \leq b \) and \([a, b] \subseteq \text{dom} \, f \) and \([a, b] \subseteq \text{dom} \, g\), then \([c, d] \subseteq \text{dom}(f - g)\).
(6) Let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathbb{R} \). Suppose \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f\). Then \( r \cdot f \) is integrable on \([c, d]\) and \((r \cdot f)|_{[c, d]} \) is bounded.
(7) Let \( f, g \) be partial functions from \( \mathbb{R} \) to \( \mathbb{R} \). Suppose that \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( g \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \( g|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \([a, b] \subseteq \text{dom} \, g\). Then \( f - g \) is integrable on \([c, d]\) and \((f - g)|_{[c, d]} \) is bounded.
(8) Suppose \( a \leq b \) and \( f \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \( c \in [a, b] \). Then \( f \) is integrable on \([a, c]\) and \( f \) is integrable on \([c, b]\) and \( \int_{a}^{b} f(x) \, dx = \int_{a}^{c} f(x) \, dx + \int_{c}^{b} f(x) \, dx \).
(9) Suppose \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f\). Then \( f \) is integrable on \([c, d]\) and \( f|_{[c, d]} \) is bounded.
(10) Suppose that \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( g \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \( g|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \([a, b] \subseteq \text{dom} \, g\). Then \( f + g \) is integrable on \([c, d]\) and \((f + g)|_{[c, d]} \) is bounded.
(11) Suppose \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f\). Then \( r \cdot f \) is integrable on \([c, d]\) and \((r \cdot f)|_{[c, d]} \) is bounded.
(12) Suppose \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f\). Then \( -f \) is integrable on \([c, d]\) and \((-f)|_{[c, d]} \) is bounded.
(13) Suppose that \( a \leq c \leq d \leq b \) and \( f \) is integrable on \([a, b]\) and \( f|_{[a, b]} \) is bounded and \( g|_{[a, b]} \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \([a, b] \subseteq \text{dom} \, g\). Then \( f - g \) is integrable on \([c, d]\) and \((f - g)|_{[c, d]} \) is bounded.
(14) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( f \) be a function from \( A \) into \( \mathcal{R}^n \). Then \( f \) is bounded if and only if \(|f|\) is bounded.
(15) If \( f \) is bounded and \( A \subseteq \text{dom} \, f \), then \( f\vert_A \) is bounded.
(16) Let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \) and \( g \) be a function from \( A \) into \( \mathcal{R}^n \). If \( f \) is bounded and \( f = g \), then \( g \) is bounded.
(17) For every partial function \( f \) from \( \mathbb{R} \) to \( \mathcal{R}^n \) and for every function \( g \) from
\( A \) into \( \mathcal{R}^n \) such that \( f = g \) holds \(|f| = |g|\).

(18) If \( A \subseteq \text{dom} \, h \), then \(|h| \, A| = |h| \, A|.

(19) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( h \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \). If \( A \subseteq \text{dom} \, h \) and \( h \, A \) is bounded, then \( |h| \, A \) is bounded.

(20) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( h \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \). Suppose \( A \subseteq \text{dom} \, h \) and \( h \, A \) is bounded and \( h \) is integrable on \( A \) and \(|h| \) is integrable on \( A \). Then \( \int_A h(x) \, dx \leq \int_A |h|(x) \, dx \).

(21) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( h \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \). Suppose \( a \leq b \) and \( [a, b] \subseteq \text{dom} \, h \) and \(|h| \) is integrable on \( [a, b] \) and \(|h| \, [a, b] \) is bounded. Then \( \int_a^b h(x) \, dx \leq \int_a^b |h|(x) \, dx \).

(22) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( f \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \). Suppose that \( a \leq b \) and \( f \) is integrable on \([a, b] \) and \(|f| \) is integrable on \([a, b] \) and \(|f|[a, b] \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \( c, d \in [a, b] \). Then \(|f| \) is integrable on \([\min(c, d), \max(c, d)] \) and \(|f|[\min(c, d), \max(c, d)] \) is bounded and \( \int_c^d f(x) \, dx \leq \int_c^d |f|(x) \, dx \).

(23) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( f \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \). Suppose that \( a \leq b \) and \( c \leq d \) and \( f \) is integrable on \([a, b] \) and \(|f| \) is integrable on \([a, b] \) and \(|f|[a, b] \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \( c, d \in [a, b] \). Then \(|f| \) is integrable on \([c, d] \) and \(|f|[c, d] \) is bounded and \( \int_c^d f(x) \, dx \leq \int_c^d |f|(x) \, dx \) and \( \int_c^d f(x) \, dx \leq \int_c^d |f|(x) \, dx \).

(24) Let \( n \) be a non empty element of \( \mathbb{N} \) and \( f \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \). Suppose that \( a \leq b \) and \( c \leq d \) and \( f \) is integrable on \([a, b] \) and \(|f| \) is integrable on \([a, b] \) and \(|f|[a, b] \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \( c, d \in [a, b] \) and for every real number \( x \) such that \( x \in [c, d] \) holds \(|f_x| \leq e \).

Then \( \int_c^d f(x) \, dx \leq e \cdot (d - c) \) and \( \int_c^d f(x) \, dx \leq e \cdot (d - c) \).

(25) If \( a \leq b \) and \( f \) is integrable on \([a, b] \) and \(|f|[a, b] \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \( c, d \in [a, b] \), then \( \int_c^d (r \cdot f) \, dx = r \cdot \int_c^d f \, dx \).

(26) If \( a \leq b \) and \( f \) is integrable on \([a, b] \) and \(|f|[a, b] \) is bounded and \([a, b] \subseteq \text{dom} \, f \) and \( c, d \in [a, b] \), then \( \int_c^d (-f) \, dx = -\int_c^d f \, dx \).
(27) Suppose that \( a \leq b \) and \( f \) is integrable on \([a, b]\) and \( g \) is integrable on \([a, b]\) and \( f \mid [a, b] \) is bounded and \( g \mid [a, b] \) is bounded and \([a, b] \subseteq \text{dom} \ f\) and \([a, b] \subseteq \text{dom} \ g\) and \( c, d \in [a, b] \). Then \( \int_{c}^{d} (f + g)(x) \, dx = \int_{c}^{d} f(x) \, dx + \int_{c}^{d} g(x) \, dx \).

(28) Suppose that \( a \leq b \) and \( f \) is integrable on \([a, b]\) and \( g \) is integrable on \([a, b]\) and \( f \mid [a, b] \) is bounded and \( g \mid [a, b] \) is bounded and \([a, b] \subseteq \text{dom} \ f\) and \([a, b] \subseteq \text{dom} \ g\) and \( c, d \in [a, b] \). Then \( \int_{c}^{d} (f - g)(x) \, dx = \int_{c}^{d} f(x) \, dx - \int_{c}^{d} g(x) \, dx \).

(29) Suppose \( a \leq b \) and \([a, b] \subseteq \text{dom} \ f\) and for every real number \( x \) such that \( x \in [a, b] \) holds \( f(x) = E \). Then \( f \) is integrable on \([a, b]\) and \( f \mid [a, b] \) is bounded and \( \int_{a}^{b} f(x) \, dx = (b - a) \cdot E \).

(30) Suppose \( a \leq b \) and for every real number \( x \) such that \( x \in [a, b] \) holds \( f(x) = E \) and \([a, b] \subseteq \text{dom} \ f\) and \( c, d \in [a, b] \). Then \( \int_{a}^{d} f(x) \, dx = (d - c) \cdot E \).

(31) If \( a \leq b \) and \( f \) is integrable on \([a, b]\) and \( f \mid [a, b] \) is bounded and \([a, b] \subseteq \text{dom} \ f\) and \( c, d \in [a, b] \), then \( \int_{a}^{d} f(x) \, dx = \int_{a}^{c} f(x) \, dx + \int_{c}^{d} f(x) \, dx \).

(32) Suppose that \( a \leq b \) and \( f \) is integrable on \([a, b]\) and \( f \mid [a, b] \) is bounded and \([a, b] \subseteq \text{dom} \ f\) and \( c, d \in [a, b] \) and for every real number \( x \) such that \( x \in [\min(c, d), \max(c, d)] \) holds \(|f_x| \leq e\). Then \( \int_{c}^{d} |f(x) - f(c)| \, dx \leq n \cdot e \cdot |d - c| \).

(33) \( \int_{a}^{b} f(x) \, dx = -\int_{b}^{a} f(x) \, dx \).

2. ON THE FUNCTIONS FROM \( \mathbb{R} \) INTO \( n \)-DIMENSIONAL REAL NORMED SPACE

Let \( R \) be a real normed space, let \( X \) be a non-empty set, and let \( g \) be a partial function from \( X \) to \( R \). We say that \( g \) is bounded if and only if:

(Def. 1) There exists a real number \( r \) such that for every set \( y \) such that \( y \in \text{dom} \ g \) holds \( \|g_y\| < r \).
Next we state a number of propositions:

(34) Let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \) and \( g \) be a partial function from \( \mathbb{R} \) to \( \langle \mathcal{E}^n, \| \cdot \| \rangle \). If \( f = g \), then \( f \) is bounded iff \( g \) is bounded.

(35) Let \( X, Y \) be sets and \( f_1, f_2 \) be partial functions from \( \mathbb{R} \) to \( \langle \mathcal{E}^n, \| \cdot \| \rangle \). Suppose \( f_1 \mid X \) is bounded and \( f_2 \mid Y \) is bounded. Then \( (f_1 + f_2) \mid (X \cap Y) \) is bounded and \( (f_1 - f_2) \mid (X \cap Y) \) is bounded.

(36) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \), \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \), \( D \) be a division of \( A \), \( p \) be a finite sequence of elements of \( \mathcal{R}^n \), and \( q \) be a finite sequence of elements of \( \langle \mathcal{E}^n, \| \cdot \| \rangle \). Suppose \( f = g \) and \( p = q \). Then \( p \) is a middle volume of \( f \) and \( D \) if and only if \( q \) is a middle volume of \( g \) and \( D \).

(37) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \), \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \), \( D \) be a division of \( A \), \( p \) be a middle volume of \( f \) and \( D \), and \( q \) be a middle volume of \( g \) and \( D \). If \( f = g \) and \( p = q \), then middle sum\((f,p) = middle sum(g,q)\).

(38) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \), \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \), \( T \) be a division sequence of \( A \), \( p \) be a function from \( N \) into \( (\mathcal{R}^n)^* \), and \( q \) be a function from \( N \) into \( (the\ carrier of \langle \mathcal{E}^n, \| \cdot \| \rangle)^* \). Suppose \( f = g \) and \( p = q \). Then \( p \) is a middle volume sequence of \( f \) and \( T \) if and only if \( q \) is a middle volume sequence of \( g \) and \( T \).

(39) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \), \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \), \( T \) be a division sequence of \( A \), \( S \) be a middle volume sequence of \( f \) and \( T \), and \( U \) be a middle volume sequence of \( g \) and \( T \). If \( f = g \) and \( S = U \), then middle sum\((f,S) = middle sum(g,U)\).

(40) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \), \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \), \( I \) be an element of \( \mathcal{R}^n \), and \( J \) be a point of \( \langle \mathcal{E}^n, \| \cdot \| \rangle \). Suppose \( f = g \) and \( I = J \). Then the following statements are equivalent:

(i) for every division sequence \( T \) of \( A \) and for every middle volume sequence \( S \) of \( f \) and \( T \) such that \( \delta_T \) is convergent and \( \lim(\delta_T) = 0 \) holds middle sum\((f,S) = \text{convergent and \(lim middle sum(f,S) = I\)}\),

(ii) for every division sequence \( T \) of \( A \) and for every middle volume sequence \( S \) of \( g \) and \( T \) such that \( \delta_T \) is convergent and \( \lim(\delta_T) = 0 \) holds middle sum\((g,S) = \text{convergent and \(lim middle sum(g,S) = J\)}\).

(41) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \) and \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \). Suppose \( f = g \) and \( f \) is bounded. Then \( f \) is integrable if and only if \( g \) is integrable.

(42) Let \( f \) be a function from \( A \) into \( \mathcal{R}^n \) and \( g \) be a function from \( A \) into \( \langle \mathcal{E}^n, \| \cdot \| \rangle \). Suppose \( f = g \) and \( f \) is bounded and integrable. Then \( g \) is integrable and integral \( f = \text{integral} g \).

(43) Let \( f \) be a partial function from \( \mathbb{R} \) to \( \mathcal{R}^n \) and \( g \) be a partial function
from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose $f = g$ and $f \upharpoonright A$ is bounded and $A \subseteq \text{dom } f$. Then $f$ is integrable on $A$ if and only if $g$ is integrable on $A$.

(44) Let $f$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$ and $g$ be a partial function from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose $f = g$ and $f \upharpoonright A$ is bounded and $A \subseteq \text{dom } f$ and $f$ is integrable on $A$. Then $g$ is integrable on $A$ and $\int_A f(x)dx = \int_A g(x)dx$.

(45) Let $f$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$ and $g$ be a partial function from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose $f = g$ and $a \leq b$ and $f \upharpoonright [a, b]$ is bounded and $[a, b] \subseteq \text{dom } f$ and $f$ is integrable on $[a, b]$. Then $\int_a^b f(x)dx = \int_a^b g(x)dx$.

(46) Let $f$, $g$ be partial functions from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose $a \leq b$ and $f$ is integrable on $[a, b]$ and $g$ is integrable on $[a, b]$ and $[a, b] \subseteq \text{dom } f$ and $[a, b] \subseteq \text{dom } g$. Then $\int_a^b (f + g)(x)dx = \int_a^b f(x)dx + \int_a^b g(x)dx$ and $\int_a^b (f - g)(x)dx = \int_a^b f(x)dx - \int_a^b g(x)dx$.

(47) For every partial function $f$ from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$ such that $a \leq b$ and $[a, b] \subseteq \text{dom } f$ holds $\int_b^a f(x)dx = -\int_a^b f(x)dx$.

(48) Let $f$ be a partial function from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$ and $g$ be a partial function from $\mathbb{R}$ to $\mathbb{R}^n$. Suppose $f = g$ and $a \leq b$ and $[a, b] \subseteq \text{dom } f$ and $f \upharpoonright [a, b]$ is bounded and $f$ is integrable on $[a, b]$ and $c, d \in [a, b]$. Then $\int_c^d f(x)dx = \int_c^d g(x)dx$.

(49) Let $f$, $g$ be partial functions from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose that $a \leq b$ and $f$ is integrable on $[a, b]$ and $g$ is integrable on $[a, b]$ and $f \upharpoonright [a, b]$ is bounded and $g \upharpoonright [a, b]$ is bounded and $[a, b] \subseteq \text{dom } f$ and $[a, b] \subseteq \text{dom } g$ and $c, d \in [a, b]$. Then $\int_c^d (f + g)(x)dx = \int_c^d f(x)dx + \int_c^d g(x)dx$.

(50) Let $f$, $g$ be partial functions from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose that $a \leq b$ and $f$ is integrable on $[a, b]$ and $g$ is integrable on $[a, b]$ and $f \upharpoonright [a, b]$ is bounded and $g \upharpoonright [a, b]$ is bounded and $[a, b] \subseteq \text{dom } f$ and $[a, b] \subseteq \text{dom } g$ and $c, d \in [a, b]$. Then $\int_c^d (f - g)(x)dx = \int_c^d f(x)dx - \int_c^d g(x)dx$.

(51) Let $E$ be a point of $\langle E^n, \| \cdot \| \rangle$ and $f$ be a partial function from $\mathbb{R}$ to $\langle E^n, \| \cdot \| \rangle$. Suppose $a \leq b$ and $[a, b] \subseteq \text{dom } f$ and for every real number
$x$ such that $x \in [a, b]$ holds $f(x) = E$. Then $f$ is integrable on $[a, b]$ and $f|_{[a, b]}$ is bounded and $\int_a^b f(x)dx = (b - a) \cdot E$.

(52) Let $E$ be a point of $\langle \mathcal{E}^n, \| \cdot \| \rangle$ and $f$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $a \leq b$ and $[a, b] \subseteq \text{dom } f$ and for every real number $x$ such that $x \in [a, b]$ holds $f(x) = E$ and $c, d \in [a, b]$. Then $\int_c^d f(x)dx = (d - c) \cdot E$.

(53) Let $f$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $a \leq b$ and $f$ is integrable on $[a, b]$ and $f|_{[a, b]}$ is bounded and $[a, b] \subseteq \text{dom } f$ and $c, d \in [a, b]$. Then $\int_a^d f(x)dx = \int_a^c f(x)dx + \int_c^d f(x)dx$.

(54) Let $f$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose that $a \leq b$ and $f$ is integrable on $[a, b]$ and $f|_{[a, b]}$ is bounded and $[a, b] \subseteq \text{dom } f$ and $c, d \in [a, b]$ and for every real number $x$ such that $x \in [\min(c,d), \max(c,d)]$ holds $\|f_x\| \leq e$. Then $\| \int_c^d f(x)dx \| \leq n \cdot e \cdot |d - c|$.

3. **Fundamental Theorem of Calculus**

The following two propositions are true:

(55) Let $n$ be a non empty element of $\mathbb{N}$ and $F$, $f$ be partial functions from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose that $a \leq b$ and $f$ is integrable on $[a, b]$ and $f|_{[a, b]}$ is bounded and $[a, b] \subseteq \text{dom } f$ and $]a, b[ \subseteq \text{dom } F$ and for every real number $x$ such that $x \in ]a, b[$ holds $F(x) = \int_a^x f(x)dx$ and $x_0 \in ]a, b[$ and $f$ is continuous in $x_0$. Then $F$ is differentiable in $x_0$ and $F'(x_0) = f_{x_0}$.

(56) Let $n$ be a non empty element of $\mathbb{N}$ and $f$ be a partial function from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$. Suppose $a \leq b$ and $f$ is integrable on $[a, b]$ and $f|_{[a, b]}$ is bounded and $[a, b] \subseteq \text{dom } f$ and $x_0 \in ]a, b[$ and $f$ is continuous in $x_0$. Then there exists a partial function $F$ from $\mathbb{R}$ to $\langle \mathcal{E}^n, \| \cdot \| \rangle$ such that $]a, b[ \subseteq \text{dom } F$ and for every real number $x$ such that $x \in ]a, b[$ holds $F(x) = \int_a^x f(x)dx$ and $F$ is differentiable in $x_0$ and $F'(x_0) = f_{x_0}$.

\footnote{Fundamental Theorem of Calculus (for $\mathbb{R}^n$)}
References


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Operations of Points on Elliptic Curve in Projective Coordinates

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Summary. In this article, we formalize operations of points on an elliptic curve over GF(p). Elliptic curve cryptography [7], whose security is based on a difficulty of discrete logarithm problem of elliptic curves, is important for information security. We prove that the two operations of points: compell_{ProjCo} and addell_{ProjCo} are unary and binary operations of a point over the elliptic curve.

The terminology and notation used here are introduced in the following papers: [5], [17], [3], [1], [13], [4], [2], [12], [14], [10], [9], [16], [15], [8], [11], and [6].

1. Arithmetic in GF(p)

For simplicity, we adopt the following convention: i, j denote integers, n denotes a natural number, K denotes a field, and a1, a2, a3, a4, a5, a6 denote elements of K.

One can prove the following propositions:

1. If a1 = −a2, then a1^2 = a2^2.
2. (1K)^{-1} = 1K.

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For all elements \(g_1, g_2, g_3, a\) of \(\text{GF}(p)\) such that \(g_1 = \text{i mod } p\) and \(g_2 = \text{j mod } p\) and \(g_3 = (i + j) \text{ mod } p\) holds \(g_1 \cdot a + g_2 \cdot a = g_3 \cdot a\).

For all elements \(g_1, g_2, a\) of \(\text{GF}(p)\) such that \(g_1 = \text{i mod } p\) and \(g_2 = \text{j mod } p\) and \(g_3 = (i - j) \text{ mod } p\) holds \(g_1 \cdot a - g_2 \cdot a = g_3 \cdot a\).

For all elements \(g_1, g_2, a\) of \(\text{GF}(p)\) such that \(g_1 = \text{i mod } p\) and \(g_2 = \text{j mod } p\) and \(i = j + 1\) holds \(g_1 \cdot a - g_2 \cdot a = a\).
Let $p$ be a 5 or greater prime number. The parameters of elliptic curve $p$ yielding a subset of $(\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p))$ is defined as follows:

(Def. 2) The parameters of elliptic curve $p = \{(a, b); a \text{ ranges over elements of } GF(p), b \text{ ranges over elements of } GF(p): \text{Disc}(a) \neq 0_{GF(p)}\}$.

Let $p$ be a 5 or greater prime number. Observe that the parameters of elliptic curve $p$ is non empty.

Let $p$ be a 5 or greater prime number and let $z$ be an element of the parameters of elliptic curve $p$. Then $z_1$ is an element of $GF(p)$. Then $z_2$ is an element of $GF(p)$.

The following proposition is true

(30) Let $p$ be a 5 or greater prime number and $z$ be an element of the parameters of elliptic curve $p$. Then $p > 3$ and $\text{Disc}(z_1) \neq 0_{GF(p)}$.

For simplicity, we adopt the following rules: $p_1, p_2, p_3$ denote sets, $P_1, P_2, P_3$ denote elements of $GF(p)$, $P$ denotes an element of $\text{ProjCo}(GF(p))$, and $O$ denotes an element of $\text{EC}_{\text{SetProjCo}}(a)$.

Let $p$ be a prime number, let $a, b$ be elements of $GF(p)$, and let $P$ be an element of $\text{EC}_{\text{SetProjCo}}(a)$. The functor $P_1$ yields an element of $GF(p)$ and is defined as follows:

(Def. 3) If $P = \{p_1, p_2, p_3\}$, then $P_1 = p_1$.

The functor $P_2$ yielding an element of $GF(p)$ is defined as follows:

(Def. 4) If $P = \{p_1, p_2, p_3\}$, then $P_2 = p_2$.

The functor $P_3$ yielding an element of $GF(p)$ is defined by:

(Def. 5) If $P = \{p_1, p_2, p_3\}$, then $P_3 = p_3$.

We now state three propositions:
(31) For every prime number \( p \) and for all elements \( a, b \) of \( GF(p) \) and for every element \( P \) of \( EC_{SetProjCo}(a) \) holds \( P = \langle P_1, P_2, P_3 \rangle \).

(32) Let \( p \) be a prime number, \( a, b \) be elements of \( GF(p) \), \( P \) be an element of \( EC_{SetProjCo}(a) \), and \( Q \) be an element of \( ProjCo(GF(p)) \). Then \( P = Q \) if and only if the following conditions are satisfied:
   (i) \( P_1 = Q_1 \).
   (ii) \( P_2 = Q_2 \), and
   (iii) \( P_3 = Q_3 \).

(33) Let \( p \) be a prime number, \( a, b, P_1, P_2, P_3 \) be elements of \( GF(p) \), and \( P \) be an element of \( EC_{SetProjCo}(a) \). If \( P = \langle P_1, P_2, P_3 \rangle \), then \( P_1 = P_1 \) and \( P_2 = P_2 \) and \( P_3 = P_3 \).

Let \( p \) be a prime number, let \( P \) be an element of \( ProjCo(GF(p)) \), and let \( C_1 \) be a function from (the carrier of \( GF(p) \)) \( \times \) (the carrier of \( GF(p) \)) \( \times \) (the carrier of \( GF(p) \)) into \( GF(p) \). We say that \( P \) is on curve defined by an equation \( C_1 \) if and only if:

(Def. 6) \( C_1(P) = 0_{GF(p)} \).

The following two propositions are true:

(34) \( P \) is on curve defined by an equation \( EC_{WEdProjCo}(a) \) iff \( P \) is an element of \( EC_{SetProjCo}(a) \).

(35) Let \( p \) be a prime number, \( a, b \) be elements of \( GF(p) \), and \( P \) be an element of \( EC_{SetProjCo}(a) \). Then \((P_2)^2 \cdot P_3 - ((P_1)^3 + a \cdot P_1 \cdot (P_3)^2 + b \cdot (P_3)^3) = 0_{GF(p)} \).

Let \( p \) be a prime number and let \( P \) be an element of \( ProjCo(GF(p)) \). The represent point of \( P \) yields an element of \( ProjCo(GF(p)) \) and is defined by:

(Def. 7)(i) The represent point of \( P = \langle P_1, (P_3)^{-1}, P_2 \cdot (P_3)^{-1}, 1 \rangle \) if \( P_3 \neq 0 \),
   (ii) the represent point of \( P = \langle 0, 1, 0 \rangle \) if \( P_3 = 0 \),
   (iii) \( P_3 = 0 \), otherwise.

The following propositions are true:

(36) Let \( p \) be a 5 or greater prime number, \( z \) be an element of the parameters of elliptic curve \( p \), and \( P \) be an element of \( EC_{SetProjCo}(z_1) \). Then the represent point of \( P \equiv P \) and the represent point of \( P \in EC_{SetProjCo}(z_1) \).

(37) Let \( p \) be a prime number, \( a, b \) be elements of \( GF(p) \), and \( P \) be an element of \( ProjCo(GF(p)) \). Suppose (the represent point of \( P \)) \( _3 = 0 \). Then the represent point of \( P = \langle 0, 1, 0 \rangle \) and \( P_3 = 0 \).

(38) Let \( p \) be a prime number, \( a, b \) be elements of \( GF(p) \), and \( P \) be an element of \( ProjCo(GF(p)) \). Suppose (the represent point of \( P \)) \( _3 \neq 0 \). Then the represent point of \( P = \langle P_1 \cdot (P_3)^{-1}, P_2 \cdot (P_3)^{-1}, 1 \rangle \) and \( P_3 \neq 0 \).

(39) Let \( p \) be a 5 or greater prime number, \( z \) be an element of the parameters of elliptic curve \( p \), and \( P, Q \) be elements of \( EC_{SetProjCo}(z_1) \). Then \( P \equiv Q \) if and only if the represent point of \( P = \) the represent point of \( Q \).
Let $p$ be a 5 or greater prime number and let $z$ be an element of the parameters of elliptic curve $p$. The functor $\text{compell}_{\text{ProjCo}}(z, p)$ yields a function from $\text{EC}_{\text{SetProjCo}}(z_1)$ into $\text{EC}_{\text{SetProjCo}}(z_1)$ and is defined as follows:

(Def. 8) For every element $P$ of $\text{EC}_{\text{SetProjCo}}(z_1)$ holds $(\text{compell}_{\text{ProjCo}}(z, p))(P) = \langle P_1, -P_2, P_3 \rangle$.

Let $p$ be a 5 or greater prime number, let $z$ be an element of the parameters of elliptic curve $p$, let $F$ be a function from $\text{EC}_{\text{SetProjCo}}(z_1)$ into $\text{EC}_{\text{SetProjCo}}(z_1)$, and let $P$ be an element of $\text{EC}_{\text{SetProjCo}}(z_1)$. Then $F(P)$ is an element of $\text{EC}_{\text{SetProjCo}}(z_1)$.

We now state a number of propositions:

(40) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $O$ be an element of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $O = \langle 0, 1, 0 \rangle$, then $(\text{compell}_{\text{ProjCo}}(z, p))(O) \equiv O$.

(41) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$ be an element of $\text{EC}_{\text{SetProjCo}}(z_1)$. Then $(\text{compell}_{\text{ProjCo}}(z, p))((\text{compell}_{\text{ProjCo}}(z, p))(P)) = P$.

(42) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$ be an element of $\text{EC}_{\text{SetProjCo}}(z_1)$. Suppose $P_3 \neq 0$. Then the represent point of $(\text{compell}_{\text{ProjCo}}(z, p))(P) = (\text{compell}_{\text{ProjCo}}(z, p))(\text{the represent point of } P)$.

(43) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. Then $P = Q$ if and only if $(\text{compell}_{\text{ProjCo}}(z, p))(P) = (\text{compell}_{\text{ProjCo}}(z, p))(Q)$.

(44) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$ be an element of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $P_3 \neq 0$, then $P \equiv (\text{compell}_{\text{ProjCo}}(z, p))(P)$ iff $P_2 = 0$.

(45) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $P_3 \neq 0$, then $P_1 = Q_1$ and $P_3 = Q_3$ iff $P = Q$ or $P = (\text{compell}_{\text{ProjCo}}(z, p))(Q)$.

(46) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. Then $P \equiv Q$ if and only if $(\text{compell}_{\text{ProjCo}}(z, p))(P) \equiv (\text{compell}_{\text{ProjCo}}(z, p))(Q)$.

(47) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. Then $P \equiv (\text{compell}_{\text{ProjCo}}(z, p))(Q)$ if and only if $(\text{compell}_{\text{ProjCo}}(z, p))(P) \equiv Q$.

(48) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. Suppose $P_3 \neq 0$ and $Q_3 \neq 0$. Then the represent point of $P = (\text{compell}_{\text{ProjCo}}(z, p))(\text{the}}
represent point of Q) if and only if $P \equiv (\text{compell}_{\text{ProjCo}}(z,p))(Q)$.

(49) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $P \equiv Q$, then $P_2 \cdot Q_3 = Q_2 \cdot P_3$.

(50) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. Suppose $P_3 \neq 0$ and $Q_3 \neq 0$. Then $P \equiv Q$ or $P \equiv (\text{compell}_{\text{ProjCo}}(z,p))(Q)$ if and only if $P_1 \cdot Q_3 = Q_1 \cdot P_3$.

(51) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $P_3 \neq 0$ and $Q_3 \neq 0$ and $P_2 \neq 0$, then if $P \equiv (\text{compell}_{\text{ProjCo}}(z,p))(Q)$, then $P_2 \cdot Q_3 \neq Q_2 \cdot P_3$.

(52) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, and $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $P \equiv (\text{compell}_{\text{ProjCo}}(z,p))(Q)$, then $P_2 \cdot Q_3 \neq Q_2 \cdot P_3$.

(53) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_5$ be an element of $\text{GF}(p)$, and $P$ be an element of $\text{EC}_{\text{SetProjCo}}(z_1)$. If $g_5 = 3 \mod p$ and $P_2 = 0$ and $P_3 \neq 0$, then $z_1 \cdot (P_3)^2 + g_5 \cdot (P_1)^2 \neq 0$.

(54) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_4$, $g_6$, $g_7$, $g_8$ be elements of $\text{GF}(p)$, $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$, and $R$ be an element of (the carrier of $\text{GF}(p)$) $\times$ (the carrier of $\text{GF}(p)$). Suppose that

(i) $g_4 = 2 \mod p$,
(ii) $g_6 = Q_2 \cdot P_3 - P_2 \cdot Q_3$,
(iii) $g_7 = Q_1 \cdot P_3 - P_1 \cdot Q_3$,
(iv) $g_8 = g_6^2 \cdot P_3 \cdot Q_3 - g_7^3 - g_4 \cdot g_7^2 \cdot P_1 \cdot Q_3$, and
(v) $R = (g_7 \cdot g_8, g_6 \cdot (g_7^2 \cdot P_1 \cdot Q_3 - g_8) - g_7^3 \cdot P_2 \cdot Q_3, g_7^3 \cdot P_3 \cdot Q_3)$. Then $g_7 \cdot P_3 \cdot R_2 = -(g_6 \cdot (R_1 \cdot P_3 - P_1 \cdot R_3) + g_7 \cdot P_2 \cdot R_3)$.

(55) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_4$, $g_6$, $g_7$, $g_8$ be elements of $\text{GF}(p)$, $P$, $Q$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$, and $R$ be an element of (the carrier of $\text{GF}(p)$) $\times$ (the carrier of $\text{GF}(p)$). Suppose that

(i) $g_4 = 2 \mod p$,
(ii) $g_6 = Q_2 \cdot P_3 - P_2 \cdot Q_3$,
(iii) $g_7 = Q_1 \cdot P_3 - P_1 \cdot Q_3$,
(iv) $g_8 = g_6^2 \cdot P_3 \cdot Q_3 - g_7^3 - g_4 \cdot g_7^2 \cdot P_1 \cdot Q_3$, and
(v) $R = (g_7 \cdot g_8, g_6 \cdot (g_7^2 \cdot P_1 \cdot Q_3 - g_8) - g_7^3 \cdot P_2 \cdot Q_3, g_7^3 \cdot P_3 \cdot Q_3)$. Then $-g_7^2 \cdot (P_3 \cdot Q_3 \cdot R_1 + P_3 \cdot Q_1 \cdot R_3 + P_1 \cdot Q_3 \cdot R_3) + P_3 \cdot Q_3 \cdot R_3 \cdot g_6^2 = 0_{\text{GF}(p)}$. 

(56) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_4$, $g_6$, $g_7$, $g_8$ be elements of $GF(p)$, $P$, $Q$ be elements of $EC_{SetProjCo}(z_1)$, and $R$ be an element of $(the\ carrier\ of\ GF(p)) \times (the\ carrier\ of\ GF(p))$. Suppose that

(i) $g_4 = 2 \mod p,$
(ii) $g_6 = Q_2 \cdot P_3 - P_2 \cdot Q_3,$
(iii) $g_7 = Q_1 \cdot P_3 - P_1 \cdot Q_3,$
(iv) $g_8 = g_6^2 \cdot P_3 \cdot Q_3 - g_3^2 - g_4 \cdot g_7^2 \cdot P_1 \cdot Q_3,$ and
(v) $R = \langle g_7 \cdot g_8, g_6 \cdot (g_7^2 \cdot P_1 \cdot Q_3 - g_8) - g_7^3 \cdot P_2 \cdot Q_3, g_7^3 \cdot P_3 \cdot Q_3 \rangle.$

Then $z_2 \cdot g_7^2 \cdot (P_3)^2 \cdot Q_3 \cdot R_3 = -g_7^2 \cdot P_3 \cdot P_1 \cdot Q_1 \cdot R_1 + (g_7 \cdot P_2 - g_6 \cdot P_1)^2.$

(57) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_4$, $g_6$, $g_7$, $g_8$ be elements of $GF(p)$, $P$, $Q$ be elements of $EC_{SetProjCo}(z_1)$, and $R$ be an element of $(the\ carrier\ of\ GF(p)) \times (the\ carrier\ of\ GF(p))$. Suppose that

(i) $g_4 = 2 \mod p,$
(ii) $g_6 = Q_2 \cdot P_3 - P_2 \cdot Q_3,$
(iii) $g_7 = Q_1 \cdot P_3 - P_1 \cdot Q_3,$
(iv) $g_8 = g_6^2 \cdot P_3 \cdot Q_3 - g_3^2 - g_4 \cdot g_7^2 \cdot P_1 \cdot Q_3,$ and
(v) $R = \langle g_7 \cdot g_8, g_6 \cdot (g_7^2 \cdot P_1 \cdot Q_3 - g_8) - g_7^3 \cdot P_2 \cdot Q_3, g_7^3 \cdot P_3 \cdot Q_3 \rangle.$

Then $z_1 \cdot g_7^2 \cdot P_3 \cdot Q_3 \cdot R_3 = g_7^2 \cdot (P_1 \cdot Q_1 \cdot R_3 + P_3 \cdot Q_1 \cdot R_1 + P_1 \cdot Q_3 \cdot R_1) + g_4 \cdot g_6 \cdot Q_3 \cdot R_3 \cdot (g_7 \cdot P_2 - g_6 \cdot P_1).$

(58) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_4$, $g_6$, $g_7$, $g_8$ be elements of $GF(p)$, $P$, $Q$ be elements of $EC_{SetProjCo}(z_1)$, and $R$ be an element of $(the\ carrier\ of\ GF(p)) \times (the\ carrier\ of\ GF(p))$. Suppose that

(i) $g_4 = 2 \mod p,$
(ii) $g_6 = Q_2 \cdot P_3 - P_2 \cdot Q_3,$
(iii) $g_7 = Q_1 \cdot P_3 - P_1 \cdot Q_3,$
(iv) $g_8 = g_6^2 \cdot P_3 \cdot Q_3 - g_3^2 - g_4 \cdot g_7^2 \cdot P_1 \cdot Q_3,$ and
(v) $R = \langle g_7 \cdot g_8, g_6 \cdot (g_7^2 \cdot P_1 \cdot Q_3 - g_8) - g_7^3 \cdot P_2 \cdot Q_3, g_7^3 \cdot P_3 \cdot Q_3 \rangle.$

Then $g_7^2 \cdot (P_3)^2 \cdot Q_3 \cdot ((R_2)^2 \cdot R_3 - ((R_1)^3 + z_1 \cdot R_1 \cdot (R_3)^2 + z_2 \cdot (R_3)^3)) = 0_{GF(p)}.$

(59) Let $p$ be a 5 or greater prime number, $z$ be an element of the parameters of elliptic curve $p$, $g_4$, $g_5$, $g_9$, $g_6$, $g_7$, $g_8$, $g_{10}$ be elements of $GF(p)$, $P$ be an element of $EC_{SetProjCo}(z_1)$, and $R$ be an element of $(the\ carrier\ of\ GF(p)) \times (the\ carrier\ of\ GF(p)) \times (the\ carrier\ of\ GF(p))$. Suppose that $g_4 = 2 \mod p$ and $g_5 = 3 \mod p$ and $g_{11} = 4 \mod p$ and $g_9 = 8 \mod p$ and $g_6 = z_1 \cdot (P_3)^2 + g_5 \cdot (P_1)^2$ and $g_7 = P_2 \cdot P_3$ and $g_8 = P_1 \cdot g_7$ and $g_{10} = g_6^2 - g_9 \cdot g_8$ and $R = \langle g_4 \cdot g_{10} \cdot g_7, g_6 \cdot (g_{11} \cdot g_8 - g_{10}) - g_9 \cdot (P_2)^2 \cdot g_7^2, g_9 \cdot g_7^3 \rangle.$ Then $g_4 \cdot g_7 \cdot P_3 \cdot R_2 = -(g_6 \cdot (P_3 \cdot R_1 - P_1 \cdot R_3) + g_4 \cdot g_7 \cdot P_2 \cdot R_3).$
Let \( p \) be a 5 or greater prime number, \( z \) be an element of the parameters of elliptic curve \( p, g_4, g_5, g_{11}, g_9, g_6, g_7, g_8, g_{10} \) be elements of \( GF(p) \), \( P \) be an element of \( EC_{SetProjCo}(z, 1) \), and \( R \) be an element of \((\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p))\) (Def. 9). Suppose that \( g_4 = 2 \mod p \) and \( g_5 = 3 \mod p \) and \( g_{11} = 4 \mod p \) and \( g_9 = 8 \mod p \) and \( g_6 = z_1 \cdot (P_3)^2 + g_5 \cdot (P_1)^2 \) and \( g_7 = P_2 \cdot P_3 \) and \( g_8 = P_1 \cdot P_2 \cdot g_7 \) and \( g_{10} = g_6 \cdot g_9 \cdot g_8 \) and \( R = (g_4 \cdot g_{10} \cdot g_7, g_6 \cdot (g_{11} \cdot g_8 - g_{10}) - g_9 \cdot (P_2)^2 \cdot g_7^2, g_9 \cdot g_7^3). \) Then \( g_{11} \cdot g_7^2 \cdot P_3 \cdot R_1 = R_3 \cdot (g_6 \cdot P_3 - g_6 \cdot P_1)^2 = g_{11} \cdot g_7^2 \cdot (P_1)^2 \cdot R_1^2 \).

Let \( p \) be a 5 or greater prime number, \( z \) be an element of the parameters of elliptic curve \( p, g_4, g_5, g_{11}, g_9, g_6, g_7, g_8, g_{10} \) be elements of \( GF(p) \), \( P \) be an element of \( EC_{SetProjCo}(z, 1) \), and \( R \) be an element of \((\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p))\) (Def. 9). Suppose that \( g_4 = 2 \mod p \) and \( g_5 = 3 \mod p \) and \( g_{11} = 4 \mod p \) and \( g_9 = 8 \mod p \) and \( g_6 = z_1 \cdot (P_3)^2 + g_5 \cdot (P_1)^2 \) and \( g_7 = P_2 \cdot P_3 \) and \( g_8 = P_1 \cdot P_2 \cdot g_7 \) and \( g_{10} = g_6 \cdot g_9 \cdot g_8 \) and \( R = (g_4 \cdot g_{10} \cdot g_7, g_6 \cdot (g_{11} \cdot g_8 - g_{10}) - g_9 \cdot (P_2)^2 \cdot g_7^2, g_9 \cdot g_7^3). \) Then \( g_{11} \cdot g_7^2 \cdot (P_3)^2 \cdot (z_1 \cdot R_3) = R_3 \cdot (g_6 \cdot g_7 \cdot P_2 - g_6 \cdot P_1)^2 = g_{11} \cdot g_7^2 \cdot (P_1)^2 \cdot R_1^2 \).

Let \( p \) be a 5 or greater prime number, \( z \) be an element of the parameters of elliptic curve \( p, g_4, g_5, g_{11}, g_9, g_6, g_7, g_8, g_{10} \) be elements of \( GF(p) \), \( P \) be an element of \( EC_{SetProjCo}(z, 1) \), and \( R \) be an element of \((\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p)) \times (\text{the carrier of } GF(p))\) (Def. 9). Suppose that \( g_4 = 2 \mod p \) and \( g_5 = 3 \mod p \) and \( g_{11} = 4 \mod p \) and \( g_9 = 8 \mod p \) and \( g_6 = z_1 \cdot (P_3)^2 + g_5 \cdot (P_1)^2 \) and \( g_7 = P_2 \cdot P_3 \) and \( g_8 = P_1 \cdot P_2 \cdot g_7 \) and \( g_{10} = g_6 \cdot g_9 \cdot g_8 \) and \( R = (g_4 \cdot g_{10} \cdot g_7, g_6 \cdot (g_{11} \cdot g_8 - g_{10}) - g_9 \cdot (P_2)^2 \cdot g_7^2, g_9 \cdot g_7^3). \) Then \( g_{11} \cdot g_7^2 \cdot (P_3)^2 \cdot (z_1 \cdot R_3) = R_3 \cdot (g_6 \cdot g_7 \cdot P_2 - g_6 \cdot P_1) + g_7^2 \cdot (g_{11} \cdot P_1 \cdot P_3 \cdot R_1 + g_4 \cdot (P_1)^2 \cdot R_3). \)

Let \( p \) be a 5 or greater prime number and let \( z \) be an element of the parameters of elliptic curve \( p \). The functor \( addell_{ProjCo}(z, p) \) yields a function from \( EC_{SetProjCo}(z, 1) \times EC_{SetProjCo}(z, 1) \) into \( EC_{SetProjCo}(z, 1) \) and is defined by the condition (Def. 9).

(Def. 9) Let \( P, Q, O \) be elements of \( EC_{SetProjCo}(z, 1) \) such that \( O = (0, 1, 0) \). Then

(i) if \( P \equiv O \), then \( (addell_{ProjCo}(z, p))(P, Q) = Q \),
(ii) if \( Q \equiv O \) and \( P \not\equiv O \), then \( (addell_{ProjCo}(z, p))(P, Q) = P \),
(iii) if $P \not= O$ and $Q \not= O$ and $P \not= Q$, then for all elements $g_4$, $g_6$, $g_7$, $g_8$ of $GF(p)$ such that $g_4 = 2 \bmod p$ and $g_6 = Q_2 \cdot P_3 - P_2 \cdot Q_3$ and $g_7 = Q_1 \cdot P_3 - P_1 \cdot Q_3$ and $g_8 = g_6^2 \cdot P_3 \cdot Q_3 - g_7^3 - g_4 \cdot g_7^2 \cdot P_1 \cdot Q_3$ holds
\[(\text{addell}_{\text{ProjCo}}(z,p))(P,Q) = (g_7 \cdot g_8, g_6 \cdot (g_7^2 \cdot P_1 \cdot Q_3 - g_8) - g_9^3 \cdot P_2 \cdot Q_3, g_7^3 \cdot P_3 \cdot Q_3),\]
and
(iv) if $P \not= O$ and $Q \not= O$ and $P = Q$, then for all elements $g_4$, $g_5$, $g_{11}$, $g_9$, $g_6$, $g_7$, $g_8$, $g_{10}$ of $GF(p)$ such that $g_4 = 2 \bmod p$ and $g_5 = 3 \bmod p$ and $g_{11} = 4 \bmod p$ and $g_9 = 8 \bmod p$ and $g_6 = z_1 \cdot (P_3)^2 + g_5 \cdot (P_1)^2$ and $g_7 = P_2 \cdot P_3$ and $g_8 = P_1 \cdot P_2 \cdot g_7$ and $g_{10} = g_6^2 - g_9 \cdot g_8$ holds
\[(\text{addell}_{\text{ProjCo}}(z,p))(P,Q) = (g_4 \cdot g_{10} \cdot g_7, g_6 \cdot (g_{11} \cdot g_8 - g_{10}) - g_9 \cdot (P_2)^2 \cdot g_7^2, g_9 \cdot g_3^3).\]

Let $p$ be a 5 or greater prime number, let $z$ be an element of the parameters of elliptic curve $p$, let $F$ be a function from $\text{EC}_{\text{SetProjCo}}(z_1) \times \text{EC}_{\text{SetProjCo}}(z_1)$ into $\text{EC}_{\text{SetProjCo}}(z_1)$, and let $Q, R$ be elements of $\text{EC}_{\text{SetProjCo}}(z_1)$. Then $F(Q,R)$ is an element of $\text{EC}_{\text{SetProjCo}}(z_1)$.

**References**


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