ON PROPERTIES OF POLYNOMIALS IN RANDOM ELEMENTS

ULYANOV V. V.†

Abstract. The paper deals with different properties of polynomials in random elements: bounds for characteristics functionals of polynomials, stochastic generalization of the Vinogradov mean value theorem, characterization problem, bounds for probabilities to hit the balls. These results cover the cases when the random elements take values in finite as well as infinite dimensional Hilbert spaces.

Key words. polynomials in random elements, bounds for characteristic functions, the Vinogradov mean value theorem, characterization problem, stability problem, distribution tails, distribution of quadratic forms, singular distributions, the Cantor distribution.

1. Introduction. Let $X$ be a random element with values in a real separable Hilbert space $\mathbb{H}$. We consider a sequence $X_1, X_2, \ldots$ of independent copies of $X$. Put

$$S_n = n^{-1/2}(X_1 + \cdots + X_n), \quad g_n(t) = \mathbb{E}\exp\{itf(S_n)\},$$

where $f$ is a measurable map from $\mathbb{H}$ to the real line $\mathbb{R}$.

In order to find estimates for the rate of convergence in functional limit theorems and to construct asymptotic expansions, we apply inequalities for $|g_n(t)|$ of the type

$$|g_n(t)| \leq c(n^{-A} + |t|^{-D}),$$

where $c, A, D$ are some constants. One can derive estimate (1) using the so-called symmetrization inequality (see [1] and [2], in particular, Corollary 3.26 and the remark after it in [2] concerning the paper of Weyl [3]; see also [4]). For example, when $\mathbb{H} = \mathbb{R}$ this approach leads to the following inequality (cf. [5, Lemmas 3.1, 4.1, 4.2]).

Set

$$f(x) = x^m + \alpha_{m-1}x^{m-1} + \cdots + \alpha_0, \quad m \geq 2.$$

Then we have (1) with $D = (m \cdot 2^m)^{-1}$ and an arbitrary $A > 0$ for $t$, satisfying the condition

$$|t| \leq c_1 n^{(m-1)/2 - \varepsilon},$$

where $\varepsilon > 0$.

It is well known that if a random vector $X$ has a discrete distribution, then $g_n(t, a)$ is an almost periodic function. Therefore,

$$\limsup_{t \to \infty} |g_n(t, a)| = 1.$$

Thus, we have (1) provided $t$ is restricted to a finite interval depending on $n$ and $f$.

In Section 2 we give the bounds for characteristic functions $g_n(t, a)$, which improve (1) so that $D$ increases essentially and becomes proportional to $1/m$. It follows from Theorem 4 with two-sided bounds for $g_n(t, a)$, that this order of $D$ is right one.

*This research was supported by RSCF 14-11-00196.
†Lomonosov Moscow State University, Faculty of Computational Mathematics and Cybernetics; National Research University Higher School of Economics (HSE), Moscow, Russia; e-mail: vulyanov@cs.msu.ru

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2. Bounds for characteristic functions of polynomials in random elements. As noted in the introduction, (1) could be improved so that $D$ becomes larger (see [6], [7]). The improvement is obtained in the following cases:

a) the distribution of $X$ has a non-degenerate discrete component;

b) for some $n_0$ the distribution of $S_{n_0}$ has an absolutely continuous component;

c) the characteristic function of $X$ satisfies some conditions “in average”, in particular, the conditions are met for a class of singular distributions.

Recall that if $P, P_1, P_2$ denote three probability measures in $H$ and $\gamma$ is a positive number, $0 < \gamma \leq 1$, then $P = \gamma P_1 + (1 - \gamma) P_2$.

In case a) the following theorem is true (see [6, Theorem 2]).

**Theorem 1.** Let $X_1, X_2, \ldots$ be independent identically distributed (i.i.d.) random variables in $R$ with a distribution $F$, admitting a non-degenerate discrete component. Then for any $\varepsilon > 0$ and integer $m \geq 2$ there exist absolute constants $c_3, c_4$ and a constant $c_2$, depending on $F, \varepsilon$ and $m$, such that, for all $n \geq 1$ and any $t$, in the range (3), one has

$$
\sup_{a \in R} |E \exp\{itf(S_n + a)\}| \leq c_2 \left( n^{-c_3/(m^2 \ln m)} + \left( \frac{\sqrt{n}}{|t|} \right)^{c_4/m} \right),
$$

where $f$ is the polynomial defined in (2).

Now we consider the case of continuous components. We assume that a random vector $Y$ in $R^k$ has a distribution $P_Y$ which admits an absolutely continuous component of weight $\lambda$. Then $P_Y$ has an absolutely continuous component of weight $\gamma$, $0 < \gamma < \lambda$, such that the Lebesgue density of this component is bounded. Take $\gamma = \lambda/2$. Then (see [14, p. 4] or [15, \S 16]) this implies that $P_Y \ast P_Y$ has a component of weight $(\lambda/2)^2$ with uniformly continuous density function on $R^k$, say, $r(x)$. If $r(a_0) = r_0 > 0$ for some $a_0$, then there exists $\eta$ such that for all $x$, with $|x-a_0| \leq \eta$, we
have $r(x) \geq r_0/2$. Thus, $P_Y * P_Y$ admits a component which is uniformly distributed on some ball.

In case b) the following theorem is true (see [6, Theorem 3]).

**Theorem 2.** Set for a vector $x = (x^{(1)}, \ldots, x^{(k)})$

$$f_0(x) = (x^{(1)})^m + \cdots + (x^{(k)})^m, \quad m \geq 2.$$  

Let $X_1, X_2, \ldots$ denote a sequence of i.i.d. random vectors in $\mathbb{R}^k$. Suppose that for some $n_0$ the distribution of $X_1 + \cdots + X_{n_0}$ has an absolutely continuous component. Let $\eta$ denote a real number such that the distribution of $X_1 + \cdots + X_{2n_0}$ has a component of weight $\delta$, uniformly distributed on the ball $\{x: |x-a_0| \leq \eta\}$. Then there exist the constants $c_5, c_6$, depending on $k, n_0, \delta, \eta$ and such that for all $n \geq 6n_0/\delta$ and $t \neq 0$ we have

$$\sup_{a \in \mathbb{R}^k} |\mathbb{E} \exp\{itf_0(S_n + a)\}| \leq c_5(\exp\{-c_6n\} + |t|^{-k/m}).$$

The expressions for $c_5, c_6$ see in [6, Theorem 3]. For components with differentiable densities Theorem 2 could be refined (see [6, Theorem 4]).

We now consider the general polynomial

$$f(x) = \sum_{m_1=0}^{m} \cdots \sum_{m_k=0}^{m} \alpha(m_1, \ldots, m_k) x_1^{m_1} \cdots x_k^{m_k},$$

where $\alpha(0, \ldots, 0) = 0$. Let $M$ stand for the degree of $f(x)$, that is,

$$M = \max\{m_1 + \cdots + m_k: \alpha(m_1, \ldots, m_k) \neq 0\}.$$

Clearly, $km \geq M$. Put

$$\alpha_* = \max|\alpha(m_1, \ldots, m_k)| > 0,$$

where the maximum is taken over all $m_1, \ldots, m_k$ such that

$$0 \leq m_1, \ldots, m_k \leq m, \quad m_1 + \cdots + m_k = M.$$

The following theorem is true (see [6, Theorem 5]).

**Theorem 3.** Let $Z$ be a random vector in $\mathbb{R}^k$ with independent standard normal coordinates $Z^{(i)}$, $i = 1, \ldots, k$. If $f(x)$, is defined by (5), with $M \geq 2$, then

$$\sup_{a \in \mathbb{R}^k} |\mathbb{E} \exp\{itf(Z + a)\}| \leq c_7|\alpha_* t|^{-1/m} \ln^{c_8}(2 + |\alpha_* t|)$$

for some constants $c_7, c_8$, depending on $k, M, m$.

The expressions for $c_7, c_8$ see in [6, Theorem 5].

Notice that until now we gave the upper bounds for characteristic functionals. In some cases it is possible to prove two-sided bounds. Moreover, the orders of the upper and lower bounds with respect to $t$ coincide. Set

$$f_1(x) = x^{(1)} \cdots x^{(k)}, \quad k \geq 2.$$

Let $Z$, as above, be a random vector in $\mathbb{R}^k$ with independent standard normal coordinates. Applying Theorem 3, we obtain

$$|\mathbb{E} \exp\{itf_1(Z)\}| \leq c(k)|t|^{-1} \ln^{k-1}(2 + |t|)$$
only. In fact, one can get not only better upper bound but prove a lower bound of the same order in $t$ as well (see [6, Theorem 6]).

**Theorem 4.** For any $k \geq 2$ and all $|t| \geq 1$ we have

$$l_k |t|^{-1} \ln k - 2 |t| \leq |\mathbb{E} \exp(itf_1(Z))| \leq L_k |t|^{-1} \ln k - 2 |t|$$

with some $l_k, L_k$, depending on $k$ only.

The expressions for $l_k, L_k$ see in [6, Theorem 6].

Now we consider the case when the characteristic function of $X$ satisfies some “averaged” conditions. Let $g_X(t) = \mathbb{E} \exp(itX)$ be the characteristic function of $X$. Put for $T > 0$

$$\phi_X(T) = \int_{-T}^{T} |g_X(t)| \, dt.$$

The following theorem is true (see [7, Theorem 1]).

**Theorem 5.** Let $X$ be a random variable such that $|X| \leq 1$. Assume that there exists a nondecreasing positive function $\phi(t)$ on $(0, \infty)$ and a constant $\varepsilon : 0 \leq \varepsilon < 1/m$, such that for all $t \geq 1$ and $b \geq 1$

$$\phi_X(bt) \leq b^{\varepsilon} \phi(t).$$

Then for $|t| \geq 1$ we have

$$|\mathbb{E} \exp(itf(X))| \leq c\phi(|t|)|t|^{-1/m},$$

where $f(x)$ is defined by (2) and the constant $c$ does not depend on $t$.

The expressions for $c$ see [7, relations (2.5) and (2.10)].

**Remark 1.** If the characteristic function $g_X(t)$ satisfies

$$\int_{-\infty}^{\infty} |g_X(t)| \, dt = A < \infty,$$

we may choose, in Theorem 5, $\varepsilon = 0$ and $\phi(t) = A$. Thus, assuming condition (8) we get

$$|\mathbb{E} \exp(itf(X))| \leq cA|t|^{-1/m}.$$

Moreover, if (8) holds we drop the assumption $|X| \leq 1$ (see [6, Lemma 7 and bounds (22), (23)]).

**Remark 2.** One can show (see [7, Corollary 3]), the distribution function of $X$, which satisfies the conditions of Theorem 5 is necessarily a Lipschitz function of order $1 - \varepsilon$.

Condition (6) is satisfied by many absolutely continuous distributions (cf. [16, Ch. VI, §3]) as well as by some singular distributions (see [17]–[19] and [16, Ch. XII, §§10, 11]).

We give an example of singular distribution of $X$, such that its convolutions of any order are singular as well. However, the behavior of the characteristic function of polynomial mapping $f(X)$ is very similar to the case when $X$ is absolutely continuous (see Theorem 2) or even Gaussian (see Theorem 3).
We consider the Cantor distribution with characteristic function

\begin{equation}
L(t) = \prod_{j=1}^{\infty} \cos(2\pi \cdot 3^{-j}t).
\end{equation}

The Cantor distribution satisfies the Cramèr condition (see [17]), i.e.,

\begin{equation}
\limsup_{|t| \to \infty} |L(t)| < 1.
\end{equation}

In [7, Theorem 2] it was proved a sharper bound than (11), namely, for all $|t| \geq 8.5$ we have

$$|L(t)| \leq e^{-0.027}.$$ 

Moreover, in fact, not only is the Cantor distribution singular but its $k$-fold convolutions of any order $k$ are singular as well (see [7, Theorem 2]). The convolutions have characteristic functions $L^k(t)$. At the same time the behavior of the characteristic function of the Cantor distribution after polynomial mapping is similar to the case of absolutely continuous distributions.

The following theorem is true (see [7, Corollary 5]).

**Theorem 6.** Let $f(x)$ be a polynomial defined in (2) with $m \geq 2$, and let $\varepsilon : 0 < \varepsilon < 1/m$. Let $X$ denote a random variable with singular distribution and characteristic function $L^k(t)$, where $k$ is any integer such that $k \geq 0.027^{-1}\ln(2/(3\varepsilon - 1))$. Then

$$E \exp\{itf(X)\} = O(|t|^\varepsilon^{-1/m}) \quad \text{as } |t| \to \infty.$$ 

The proofs of Theorems 1–4 are based on inequalities for trigonometric sums and integrals. If the distribution of $X$ has a non-degenerate discrete component we use estimates of the Vinogradov-type for trigonometric sums (see [20, Statement 4.2 in §4, Ch. VIII]). If for some $n_0$ the distribution of $S_{n_0}$ has an absolutely continuous component, we apply the Vinogradov inequality for trigonometric integrals (see [21, Lemma 4, Ch. 2]) and its generalizations to the multidimensional case (see [22, Theorem 5 in §3, Ch. 1]). Our improvements of estimate (1) obtained in [6] are similar in form with the improvements which gave the Vinogradov method for trigonometric sums comparing the Weyl results (see Introduction in [21]).

### 3. Stochastic generalization of the Vinogradov mean value theorem.

A desire to remove additional conditions on the type of the distribution $X$ leads to the necessity of stochastic generalization of the famous Vinogradov theorem on the mean. Theorems 1–4 are proved provided additional conditions on the type of the distribution $X$ are met. A desire to remove the conditions leads to the necessity of stochastic generalization of the famous Vinogradov mean value theorem (see [21]).

The stochastic generalization was obtained in [8]. Introduce necessary notation.

Let $J_k(P)$ denote the number of the simultaneous diophantine equations

$$\sum_{i=1}^{k} (x_i^j - y_i^j) = 0, \quad j = 1, \ldots, m,$$

where $1 \leq x_i, y_i \leq P$ for $i = 1, \ldots, k$. 

Since for any integer $x$ we have
$$\int_0^1 \exp\{2\pi i \alpha x\} \, d\alpha = \begin{cases} 1, & \text{if } x = 0, \\ 0, & \text{if } x \neq 0, \end{cases}$$
we get
$$J_k(P) = \int_0^1 \ldots \int_0^1 \left| \sum_{x=1}^P \exp\{2\pi i f(x)\} \right|^{2k} \, d\alpha_m \cdots d\alpha_1,$$
where $f(x) = \alpha_m x^m + \cdots + \alpha_1 x$.

Estimates of the magnitude order of the growth of $J_k(P)$ with increasing $P$ have important applications in analytic number theory. In 1934 I.M. Vinogradov proved the following theorem for $J_k(P)$.

**Theorem 7.** Let $\tau > 0$, $m > 2$ be integers with $k \geq m\tau$, and let $P \geq 1$. Then
$$J_k(P) \leq c_{\tau} \cdot P^{2k-\Delta(\tau)},$$
where
$$\Delta(\tau) = 0.5m(m+1)(1 - (1 - m^{-1})^\tau), \quad c_{\tau} = (m\tau)^{6m}(2m)^{4m(m+1)\tau}.$$

The bounds for $J_k(P)$ provide information about the size of the trigonometric Weyl sums
$$F = \sum_{x=1}^P \exp\{2\pi i f(x)\}$$
and trigonometric sums involving functions to which there are reasonable polynomial approximations. Since there are numerous applications for such bounds Vinogradov was able to use his method with great success in many different problems of number theory. For later refinements of the Vinogradov mean value theorem see, e.g., [23]–[25], and the discussion and references therein.

It is obvious that if a random variable $S$ takes values $1, 2, \ldots, P$ with the equal probability $1/P$, then
$$E \exp\{2\pi i f(S)\} = \frac{F}{P}.$$

In Section 2 we showed that using the bounds for $|F|$, known in the number theory, one can improve the upper bounds for $|E \exp\{2\pi i f(S_n)\}|$, where $S_n$ is a normalized sum of i.i.d. random variables $X_1, \ldots, X_n$ such that the distribution of $X_1$ has non-degenerate component (see [6]).

Stochastic generalization of the Vinogradov mean value theorem which is of independent interest, could be applied to find the upper bounds for $|E \exp\{2\pi i f(S_n)\}|$ without any additional conditions for the type of distribution of $X_1$.

The following theorem is true (see [8, Theorem 2]).

**Theorem 8.** Let $S$ be a random variable and $P \geq 1$, $m \geq 2$, $\tau \geq 1$ and $k$ denote the positive integers. Set $f(x) = \alpha_m x^m + \cdots + \alpha_1 x$ and

$$I_k(P) = P^{-(m-1)m/2} \times \int_{-P}^{P} \ldots \int_{-P}^{P} \left| E \exp\{2\pi i f(S)\} \mathbf{1}_{\{S \leq P\}} \right|^{2k} \, d\alpha_m \cdots d\alpha_1,$$

(12)
where \(1_A\) denote the indicator function of \(A\). Then for all \(k = m\tau\) there exists a constant \(D_\tau = D(k, \tau)\), depending on \(m\) and \(\tau\), such that one has

\[
I_k(P) \leq D_\tau P^{2k - \Delta(\tau)} \left( \sup_I P(S \in I) \right)^{2k},
\]

where

\[
\Delta(\tau) = 0.5m(m + 1)(1 - (1 - m^{-1})^\tau)
\]

and supremum is taken over all intervals \(I = \{x: a < x \leq a + 1\}\) with arbitrary real \(a\).

\textbf{Remark 3.} If \(S\) takes values \(1, 2, \ldots, P\) with equal probabilities then \(\mathbb{E} \exp\{2\pi i f(S)\}\) is a periodic function of period 1 with respect to each coefficient \(\alpha_1, \ldots, \alpha_k\). Thus, \(I_k(P) = 2^m P^{-2k} J_k(P)\), and our result coincides with the original Vinogradov estimate with respect to \(P\).

\textbf{Theorem 9.} Let \(m\) and \(b\) be positive integers and \(b \geq (m + 1)/2\). Let \(P\) be a real number, \(P \geq 32\), and \(S\) be a random variable with uniform distribution on \([-P, P]\). Then for \(k = bm\) we have

\[
I_k(P) \leq 2^{5m^2 + m} P^{-m^2} \frac{1}{1 - m/(2b)}.
\]

\textbf{Remark 4.} The proof of Theorem 8 is based on the original proof of the Vinogradov mean value theorem (see [21]). It is more suitable for generalizations to arbitrary random variables with non-lattice support than \(p\)-adic proof due to Linnik, Karatsuba and Wooley (see, e.g., [22], [24], [25]). In particular, Theorem 8 is proved by induction with respect to two parameters \(P\) and \(\tau\) simultaneously. Theorem 10 is the first step of induction when \(\tau = 1\) and \(P\) is arbitrary. This step is trivially clear for lattice case in the original Vinogradov mean value theorem.

\textbf{Theorem 10.} Suppose the conditions of Theorem 8 are met. Then for all \(k, 1 \leq k \leq m\), we have

\[
I_k(P) \leq M_k P^{(3k - 1)/2} \left( \sup_I P(S \in I) \right)^{2k},
\]

where \(M_k\) is some constant depending on \(k\) and \(m\) only.

\textbf{4. Characterization and stability properties of finite quadratic forms.} Now consider the particular case of functionals – quadratic forms. The case is very important both in theory and in practice. We start with characterization properties of finite quadratic forms.

Let \(Z_1, \ldots, Z_n\) be i.i.d. standard normal random variables and \(a_1, \ldots, a_n\) be real numbers with \(a_1 + \cdots + a_n \neq 0\) and \(a_1^2 + \cdots + a_n^2 \neq 0\). Suppose that \(X_1, \ldots, X_n\) are i.i.d. random variables such that \(a_1 Z_1 + \cdots + a_n Z_n \overset{d}{=} a_1 X_1 + \cdots + a_n X_n\), where \(\overset{d}{=}\) denotes the equality in distribution. Then, by Cramér’s decomposition theorem for the normal law (see [27, Theorem 3.1.4]) the random variables \(X_i\) are standard normal as well.

Lukacs and Laha (see [28, Theorem 9.1.1]) considered more general problem. Namely, let \(X_1, \ldots, X_n\) be i.i.d. random variables such that their linear combination
\[ L = a_1 X_1 + \cdots + a_n X_n \] has analytic characteristic function and \( a_1^s + \cdots + a_n^s \neq 0 \) for all \( s = 1, 2, \ldots \). Then the distribution of \( X_1 \) is uniquely determined by that of \( L \).

The aim of this Section is to obtain similar characterization property for quadratic forms in i.i.d. random variables \( Z_1, \ldots, Z_n \). Furthermore, we state stability property of such quadratic forms.

Consider symmetric matrix \( A = (a_{ij})_{i,j=1}^n \). Let

\[ Q(x_1, \ldots, x_n) = \sum_{i,j=1}^n a_{ij} x_i x_j \]

be quadratic form in variables \( x_1, \ldots, x_n \). Assume that \( Q \) is non-degenerate in the sense that \( A \) is not zero matrix. Suppose \( Z_1, \ldots, Z_n \) are i.i.d. random variables with symmetric distribution \( F \).

We say that pair \( (Q, F) \) has characterization property (CP), iff for sequence of i.i.d. symmetric random variables \( X_1, \ldots, X_n \) the equality

\begin{equation}
Q(Z_1, \ldots, Z_n) \overset{d}{=} Q(X_1, \ldots, X_n)
\end{equation}

implies \( Z_1 \overset{d}{=} X_1 \).

**Remark 6.** We require in the definition of CP that the random variables \( X_1, \ldots, X_n \) are symmetric. Otherwise the problem does not have solution even in the case \( n = 1 \) and \( Q(x_1) = x_1^2 \). Equation (14) holds for \( X_1 = Z_1 \), as well as for \( X_1 = |Z_1| \).

**Remark 7.** With symmetric distribution \( F \) an answer is trivial in the one dimensional case, i.e. any pair \( (Q, F) \) has CP. Therefore we assume that \( n \geq 2 \) everywhere below.

We are interested to find the sufficient conditions in order the pair \( (Q, F) \) has CP. The solution of the problem depends also on the coefficients of the matrix \( A \), where the following possibilities exist:

1. \( a_{ii} = 0 \) for all \( i = 1, \ldots, n \).
2. \( a_{ii} \neq 0 \) for some \( i = 1, \ldots, n \).
   2.1. \( a_{11}^{2k+1} + a_{22}^{2k+1} + \cdots + a_{nn}^{2k+1} \neq 0 \) for all \( k = 0, 1, 2, \ldots \).
   2.2. \( a_{11} + a_{22} + \cdots + a_{nn} = 0 \).
      2.2.1. \( a_{ij} = 0 \) for all \( i \neq j \).
      2.2.2. \( a_{ij} \neq 0 \) for some \( i \neq j \).
2.3. \( a_{11}^{2k+1} + a_{22}^{2k+1} + \cdots + a_{nn}^{2k+1} = 0 \) for some \( k = 1, 2, \ldots \).

So far the answers are known for the cases 1, 2.1 and 2.2.1 only.

Define now class \( \mathcal{F} \) of probability distributions so that \( F \in \mathcal{F} \) iff the following two conditions are satisfied:

(C1) \( F \) has moments \( \alpha_k = \int_{-\infty}^{\infty} x^k \, dF(x) \) of all orders \( k \);
(C2) \( F \) is uniquely specified \( \alpha_1, \alpha_2, \ldots \).

The following examples demonstrate when probability distribution \( F \) belongs to \( \mathcal{F} \).

**Example 1.** If \( F \) has an analytic characteristic function, then \( F \in \mathcal{F} \).

Recall (see, e.g., [29, §7.2]), that characteristic function is analytic iff

(i) condition (C1) is met and
(ii) \( \limsup_{n \to \infty} a_{2n}^{1/(2n)} /(2n) < \infty \).

The latter condition leads to (C2) (see, e.g., [28, Ch. 9]).
We say that probability distribution $F$ satisfies Cramér’s condition CC, iff

$$\int_{-\infty}^{\infty} \exp\{h|x|\} \ dF(x) < \infty \quad \text{for some } h > 0.$$  

Let $F$ satisfies CC, then $F \in \mathcal{F}$.

It follows from the fact that $F$ satisfies CC iff its characteristic function is analytic (see, e.g., [29, §7.2]).

Example 2. If the moments $\{\alpha_k\}$ of $F$ satisfy the Carleman condition, i.e.

\[
\sum_{n=1}^{\infty} \alpha_{2n}^{-1/(2n)} = \infty,
\]

then $F \in \mathcal{F}$.

In fact, the condition (15) yields the uniqueness of the moment problem for $\mathcal{F}$ (see, e.g., [30, Theorem 1.10]).

Note that the Carleman condition is weaker than CC. Other examples of probability distributions belonging to $\mathcal{F}$, as well as detailed discussion concerning the moment problem and other related topics see, e.g., [31], [32, Sec. VII.3] and [33, Sec. 8.12 and 11].

The following theorem is true (see [10, Theorem 3.2.1]).

Theorem 11. Let $F \in \mathcal{F}$ and the matrix $A$ be such that $a_{ii} = 0$ for all $i = 1, \ldots, n$. Then $(Q,F)$ has CP.

Example 3. Let $Z_1, Z_2, Z_3$ be i.i.d. standard normal random variables and $X_1, X_2, X_3$ be i.i.d. symmetric random variables such that $Z_1Z_2 - Z_2Z_3 \overset{d}{=} X_1X_2 - X_2X_3$, then by Theorem 12 the random variables $X_1, X_2, X_3$ are standard normal.

The following theorem is true (see [10, Theorem 3.2.2]).

Theorem 12. Let $F \in \mathcal{F}$ and the matrix $A$ be such that $a_{11}^{2k+1} + a_{22}^{2k+1} + \cdots + a_{nn}^{2k+1} \neq 0$ for all $k = 0, 1, 2, \ldots$. Then $(Q,F)$ has CP.

Example 4. Let $Z_1, Z_2$ be i.i.d. random variables with distribution $F$ and density function

\[
p(x) = \frac{1}{4} \exp\{-|x|^{1/2}\}, \quad x \in (-\infty, \infty).
\]

Then $F \in \mathcal{F}$ (see [33, Ch. 11]).

Let $X_1, X_2$ be i.i.d. symmetric random variables such that $2Z_1^2 + 4Z_1Z_2 - Z_2^2 \overset{d}{=} 2X_1^2 + 4X_1X_2 - X_2^2$. Then by Theorem 2 the random variables $X_1$ and $X_2$ have the density function defined in (16) as well.

The following theorem is true (see [10, Theorem 3.2.3]).

Theorem 13. Let $a_{ii} \neq 0$ for some $i \in \{1, \ldots, n\}$, but $a_{11} + \cdots + a_{nn} = 0$ and $a_{ij} = 0$ for all $j \neq i$. Then for any $F$ the pair $(Q,F)$ does not have CP.

Example 5. Let $Z$ be random variable with symmetric distribution $F$ independent of the random variable $\zeta$ with $P(\zeta = 1) = P(\zeta = -1) = 1/2$, and let $c > 0$ be real constant. Put $X = \zeta(Z^2 + c)^{1/2}$. Suppose now that both $Z, Z_1, Z_2, \ldots, Z_n$ and $X, X_1, X_2, \ldots, X_n$ are i.i.d. too. Under the conditions of Theorem 13, varying the constant $c$, we find family of symmetric distributions of $X_1$ such that (14) holds. In particular, if $Z_1^2 - Z_2^2 \overset{d}{=} X_1^2 - X_2^2$, then the distributions of $X_1$ and $Z_1$ may differ.

Example 5 proves Theorem 13. The proofs of Theorems 11 and 12 are based on the following facts:
a) if $F \in \mathcal{F}$, then $X_1$ has moments $\mathbf{E} X^k$ of all orders $k$;
b) under the given conditions, we have $\mathbf{E} X^k = \mathbf{E} Z^k$ for all $k = 1, 2, \ldots$.

The following stability theorem is true as well (see [10, Theorem 3.2.4]).

**Theorem 14.** Suppose that the pair $(Q, F)$ has CP. Let $X_{N,1}, \ldots, X_{N,n}$ for $N = 1, 2, \ldots$, be series f.i.d. symmetric random variables and

$$Q(X_{N,1}, \ldots, X_{N,n}) \overset{d}{\rightarrow} Q(Z_1, \ldots, Z_n) \quad \text{as } N \rightarrow \infty,$$

where $\overset{d}{\rightarrow}$ denotes the convergence in distribution. Then

$$X_{N,1} \overset{d}{\rightarrow} Z_1 \quad \text{as } N \rightarrow \infty.$$

We use the tightness of the converging sequences of quadratic forms while proving Theorem 14.

5. **Distributions of quadratic forms in Gaussian random variables.** Let $\mathbb{H}$ be a real separable Hilbert space with inner product $(\cdot, \cdot)$ and norm $|x| = (x, x)^{1/2}$, $x \in \mathbb{H}$. We denote by $Y$ an $\mathbb{H}$-valued Gaussian random element with zero mean and covariance operator $V$, i.e., the characteristic functional of $Y$ has the form

$$\mathbf{E} \exp\{i(x,Y)\} = \exp\left(-\frac{(Vx,x)}{2}\right), \quad x \in \mathbb{H},$$

where $V$ is defined for all $x, z \in \mathbb{H}$ by $(Vx,z) = \mathbf{E}(Y,x)(Y,z)$. Assume that the eigenvalues $\sigma_j^2$, $j = 1, 2, \ldots$, of the covariance operator $V$ are ordered so that

$$\sigma_1^2 = \sigma_2^2 = \cdots = \sigma_k^2 > \sigma_{k+1}^2 \geq \sigma_{k+2}^2 \geq \cdots$$

with some $k \geq 1$, which is the multiplicity of $\sigma_1^2$. It is well known that there exists an orthonormal system $\{e_i\}_{i=1}^\infty$ in $\mathbb{H}$ such that $V e_i = \sigma_i^2 e_i$, or $Y = \sum_{i=1}^\infty Y_i e_i$, where $Y_i = (Y,e_i)$ are independent real Gaussian random variables with zero mean and variances $\sigma_i^2$. The density function of the random variable $|Y - a|^2$ will be denoted by $p(u,a)$, $u \geq 0$, $a \in \mathbb{H}$. Let $f_k(u,a)$ be a density function of $\sum_{i=1}^k (Y_i - a_i)^2$. We recall (see, e.g., [34, vol. 6, p. 421]) that the exact expression for the density $f_k(u,a)$ is known.

Set $R = \sum_{j=k+1}^\infty (Y_j - a_j)^2$. Note that $\mathbf{E}R = \sum_{j=k+1}^\infty (\sigma_j^2 + a_j^2)$.

Let $u_0 = 2k(1 - \sigma_{k+1}^2/\sigma_1^2)^{-1/2} \mathbf{E} R$, $a_i = (a_1, \ldots, a_i)$.

The following theorem is true (see [11, Theorem 1]).

**Theorem 15.** Suppose $k \geq 3$.

a) For all $u \geq 0$ $a \in \mathbb{H}$ the upper bound

$$p(u,a) \leq f_k(u,a) \mathbf{E} \exp\left\{\frac{R}{2\sigma_1^2}\right\},$$

holds, where

$$\mathbf{E} \exp\left\{\frac{R}{2\sigma_1^2}\right\} = \prod_{j=k+1}^{\infty} \left(1 - \frac{\sigma_j^2}{\sigma_1^2}\right)^{-1/2} \exp\left\{\frac{a_j^2}{2(\sigma_j^2 - \sigma_1^2)}\right\}. \quad (18)$$

b) Let one of the following conditions i)–ii) be met:
i) $k = 3$, $u \geq u^* = 4.9u_0 + 16.94|\pi_3|^2\sigma_1^{-4}u_0^2$;

ii) $k \geq 4$, $u \geq u^{**} = 5.625(k-1)^2(k-3)^{-1}u_0 + 16|\pi_3|^2\sigma_1^{-4}u_0^2(k-3)^{-2}$. Then the lower bound

\begin{equation}
p(u, a) \geq 0.125f_k(u, a)\ E \exp \left\{ \frac{R}{2\sigma^2_i} \right\}.
\end{equation}

holds.

It follows from (17) and (18) that the properties of $p(u, a)$ with respect to $u$ and $\pi_k$ could be described via properties of $f_k(u, a)$. For example, it was proved in [35] (see (15) in [35]) the following inequality: for all $k \geq 1$, $u > 0$ and arbitrary $a \in \mathbb{R}$ one has

\[ f_k(u, a) \leq \left( 2\sigma_i^2 \Gamma \left( \frac{k}{2} \right) \right)^{-1} \left( \frac{u}{2\sigma_i^2} \right)^{k/2 - 1} \exp \left\{ -\frac{(u^{1/2} - |\pi_k|)^2}{2\sigma_i^2} \right\}. \]

In Lemma 2 in [36] another bound was obtained: if $k \geq 1$ and $|\pi_k| > 0$, then for all $u > 0$ one has

\[ f_k(u, a) \leq c_1(k)\sigma_i^{-1}u^{(k-3)/4}|\pi_k|^{-(k-1)/2} \exp \left\{ -\frac{(u^{1/2} - |\pi_k|)^2}{2\sigma_i^2} \right\}, \]

where $c_1(k) = \pi^{-1/2} + ((k-1)/2)^{(k-1)/2}/\Gamma(k/2)$ and $\Gamma(\cdot)$ is the Gamma function.

Furthermore, since

\[ P(|Y - a| > r) = \int_{r^2}^{\infty} p(u, a) \, du, \]

the two-sided bounds for $p(u, a)$ imply two-sided bounds for the probability $P(|Y - a| > r)$. In particular, the following theorem is true (see [11, Corollary 4.i]).

**Theorem 16.** Let $k \geq 4$ and $r, a$ be such that $|\pi_k| > 0$ and $r > \sigma^2_i / |\pi_k| + 2|\pi_k| + (u^{**})^{1/2}$. Then there exist constants $c_9(k), c_{10}(k)$, depending on $k$, only such that one has

\[ c_9(k) \ E \exp \left\{ \frac{R}{2\sigma_i^2} \right\} \leq P(|Y - a| > r) \exp \left\{ \frac{(r - |\pi_k|)^2}{2\sigma_i^2} \right\} \sigma_i r^{(k-3)/2}|\pi_k|^{(k-1)/2} \]

\[ \leq c_{10}(k) \ E \exp \left\{ \frac{R}{2\sigma_i^2} \right\}. \]

Another type of bounds for $P(|Y - a| > r)$ see in [11] and [37], and in the papers from bibliographies in [11] and [37].

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