# RATES IN THE CLT FOR VECTOR-VALUED RANDOM FIELDS 

ALEXANDER BULINSKI* and ALEXEY SHASHKIN<br>Dept. of Mathematics and Mechanics, Moscow State University, Moscow 119992, Russia<br>* E-mail: bulinski@mech.math.msu.su abulinsk@u-paris10.fr

The Lindeberg function or the Lyapunov fraction are used to establish convergence rates in the CLT for vector-valued random fields possessing dependence structure more general than positive or negative association. Thus a generalization of the classical Newman CLT is obtained. The Stein method and the Bernstein block techniques are employed. An application to kernel estimates for the density of a stationary random field is provided.

Keywords: dependence conditions, kernel estimates of density, random fields, rates in the CLT.

## 1. Introduction

The aim of this paper is to establish convergence rates in the CLT for sums of dependent multiindexed random vectors with values in $\mathbb{R}^{s}$. We develop an approach to description of the dependence structure proposed by Doukhan and Louhichi (1999) for stochastic processes and by Bulinski and Suquet (2001) for random fields.

Let $X=\left\{X_{t}, t \in T\right\}$ be a random field defined on a probability space $(\Omega, \mathcal{F}, \mathrm{P})$ such that $X_{t}$ takes values in a metric space $(\mathrm{M}, \varkappa)$ for each $t \in T$. The key idea is to measure, for finite disjoint sets $I, J \subset T$ (with cardinalities $|I|,|J|)$, the dependence between collections of random variables $X_{I}=\left(X_{t}, t \in\right.$ $I)$ and $X_{J}=\left(X_{t}, t \in J\right)$ in terms of a functional

$$
\begin{equation*}
F(f, g ; I, J)=\left|\operatorname{cov}\left(f\left(X_{I}\right), g\left(X_{J}\right)\right)\right| \tag{1}
\end{equation*}
$$

where $f: \mathrm{M}^{I} \rightarrow \mathbb{R}, g: \mathrm{M}^{J} \rightarrow \mathbb{R}$ belong to specified classes of "test functions" (whenever the covariance exists). Some restrictions can be imposed on $I$ and $J$ as well, for instance, $|I|=1$.

We use here classes of bounded Lipschitz functions $f, g$. Recall that $G$ : $\mathrm{K} \rightarrow \mathrm{L}$ (where $(\mathrm{K}, \tau)$ and $(\mathrm{L}, \nu)$ are some metric spaces) is a Lipschitz function
if

$$
\begin{equation*}
\operatorname{Lip}(G)=\sup _{x \neq y} \frac{\nu(G(x), G(y))}{\tau(x, y)}<\infty . \tag{2}
\end{equation*}
$$

When $\mathrm{K}=\mathrm{M}^{I}$ we take $\tau(x, y)=\sum_{t \in I} \varkappa\left(x_{t}, y_{t}\right)$ for $x=\left(x_{t}, t \in I\right), y=\left(y_{t}, t \in\right.$ $I)$ in formula (2). Let $B L(\mathrm{~K})$ denote the class of bounded Lipschitz functions $G: \mathrm{K} \rightarrow \mathbb{R}$ (in $\mathbb{R}$ we use the Euclidean distance).

Often it is natural to suppose, when $T$ is endowed with a metric $\rho$, that the dependence between $X_{I}$ and $X_{J}$ is "rather small" if the distance $\rho(I, J)=$ $\inf \{\rho(t, v): t \in I, v \in J\}$ is "large enough". At the same time the dependence can increase if the distance $\rho(I, J)$ is fixed but $I$ and $J$ are growing in a sense.

To give an exact formulation consider $T=\mathbb{Z}^{d}$ and introduce a set $\Theta$ consisting of functions $\theta(I, J)$ depending on finite disjoint sets $I, J \subset \mathbb{Z}^{d}$ such that

$$
\begin{equation*}
\theta\left(\tau_{n} I, \tau_{m} J\right) \rightarrow 0 \text { as }|n-m| \rightarrow \infty \quad\left(n, m \in \mathbb{Z}^{d}\right) \tag{3}
\end{equation*}
$$

where $\tau_{n} I=\{t+n: t \in I\}$ is a shift of $I,|n|=\max _{i=1, \ldots, d}\left|n_{i}\right|$.
For example, $\theta \in \Theta$ if

$$
\begin{equation*}
\theta(I, J) \leq a(|I|,|J|) u(\rho(I, J)) \tag{4}
\end{equation*}
$$

where a function $a \geq 0$ is nondecreasing in each variable and $u(r) \searrow 0$ as $r \rightarrow \infty$.

Definition 1 (Bulinski and Suquet (2001)). A random field $X=\left\{X_{j}, j \in\right.$ $\left.\mathbb{Z}^{d}\right\}$ with values in a metric space $(\mathrm{M}, \varkappa)$ is called $(B L, \theta)$-dependent if there is a function $\theta \in \Theta$ such that

$$
\begin{equation*}
F(f, g ; I, J) \leq \operatorname{Lip}(f) \operatorname{Lip}(g) \theta(I, J) \tag{5}
\end{equation*}
$$

for all finite disjoint sets $I, J \subset \mathbb{Z}^{d}$ and any $f \in B L\left(\mathrm{M}^{I}\right), g \in B L\left(\mathrm{M}^{J}\right)$.
The appearance of Lipschitz constants in the right-hand side of (5) is clear since covariance is a bilinear function and $\operatorname{Lip}(c f)=|c| \operatorname{Lip}(f)$ for every $c \in \mathbb{R}$.

The motivation for the concept of $(B L, \theta)$-dependence is the following. There are a number of interesting models described by means of families of random variables possessing properties of positive or negative association or their modifications. For definitions and examples we refer to the pioneering papers by Harris (1960), Lehmann (1966), Esary et al. (1967), Fortuin et al. (1971), Joag-Dev and Proschan (1983); see also, e.g., Pitt (1982), Newman (1984), Lindqvist (1988), Evans (1990), Lee et al. (1990), Rachev and Xin
(1996), Ebrahimi (2002). Due to Bulinski and Shabanovich (1998) for a positively or negatively associated real-valued random field $X=\left\{X_{j}, j \in \mathbb{Z}^{d}\right\}$ having finite second moments the inequality (5) holds with

$$
\begin{equation*}
\theta(I, J)=\sum_{i \in I} \sum_{j \in J}\left|\operatorname{cov}\left(X_{i}, X_{j}\right)\right| . \tag{6}
\end{equation*}
$$

So, for this function $\theta$ the bound (4) is valid with $a(|I|,|J|)=\min \{|I|,|J|\}$ and with an analogue of the Cox-Grimmett coefficient

$$
u(r)=\sup _{j \in \mathbb{Z}^{d}} \sum_{q:|q-j| \geq r}\left|\operatorname{cov}\left(X_{j}, X_{q}\right)\right|, \quad r \geq 1 .
$$

Thus, $\theta$ appearing in (6) satisfies (3) if $u(1)<\infty$ and $u(r) \rightarrow 0$ as $r \rightarrow \infty$.
In other words Definition 1 provides a unified approach to studying both families of positively or negatively dependent random variables.

A variant of inequality (5) for smooth functions $f$ and $g$ in associated real-valued random variables was firstly established by Birkel (1988) (related results were proved by Newman (1984), Roussas (1994), Peligrad and Shao (1995), Bulinski (1996)). Some modifications of association for vector-valued processes and random fields leading to (5) were used by Burton et al. (1986), Bulinski (2000), Shashkin (2002).

Note that the choice of indicator functions $f$ and $g$ as the "test-functions" in (1) would lead to the Rosenblatt-type mixing coefficient (see, e.g., Doukhan (1994) and the references therein showing that the calculation or estimation of mixing coefficients is in general a difficult problem whereas using of a covariance function is much more simple). The choice of certain power-type functions $f$ and $g$ in (1) was applied by Bakhtin and Bulinski (1997) to get bounds for absolute moments of partial sums of multiindexed dependent random variables. Linear functions $f$ and $g$ and a correlation coefficient instead of covariance in (1) was recently used by Bradley (2002) to study the boundedness properties for spectral density of weakly stationary random field. The choice of "complex exponential" functions is discussed by Jakubowski (1993), Doukhan and Louhichi (1999).

Remark 1. Following Doukhan and Louhichi (1999) we can define the dependence conditions for a field $X=\left\{X_{t}, t \in T\right\}$ by means of specified "test functions" $f$ and $g$ and inequalities

$$
F(f, g ; I, J) \leq c(f, g ;|I|,|J|) v(\rho(I, J))
$$

where $I$ and $J$ are finite disjoint subsets of $T, c$ is a nonnegative function (nondecreasing in $|I|$ and $|J|)$ and $v(r) \rightarrow 0$ as $r \rightarrow \infty$.

Remark 2. In many problems we need not consider the whole random field $X$ on $\mathbb{Z}^{d}$ but only "a part" $X_{U}=\left(X_{j}, j \in U\right), U \subset \mathbb{Z}^{d},|U|<\infty$. Then it is sufficient to use $I, J \subset U, I \cap J=\varnothing$ in (5). Moreover, we can introduce

$$
\begin{equation*}
\theta_{1}=\theta_{1}\left(X_{U}\right)=\sup F(f, g ;\{j\}, U \backslash\{j\}) \tag{7}
\end{equation*}
$$

where the supremum is taken over all $j \in U$ and all $f \in B L(\mathrm{M}), g \in$ $B L\left(\mathrm{M}^{U \backslash\{j\}}\right)$ with $\operatorname{Lip}(f) \leq 1, \operatorname{Lip}(g) \leq 1$. Note that in (7) a set $U$ need not be a subset of $\mathbb{Z}^{d}$, that is we can use any finite collection of random variables $X_{t}, t \in U$, with values in some metric space $(\mathrm{M}, \varkappa)$.

Further on let $T=\mathbb{Z}^{d}$ and $\mathrm{M}=\mathbb{R}^{s}$, that is we study a random field $X=\left\{X_{j}, j \in \mathbb{Z}^{d}\right\}$ with values in $\mathbb{R}^{s}$. As usual $E Y$ and $\operatorname{Var}(Y)$ denote respectively the mean and covariance matrix of a random vector $Y$ defined on a probability space $(\Omega, \mathcal{F}, \mathrm{P})$.

Let $U$ be a finite subset of $\mathbb{Z}^{d}$. Assume that

$$
\begin{equation*}
\mathrm{E} X_{j}=0 \in \mathbb{R}^{s}, \quad \mathrm{E}\left\|X_{j}\right\|^{2}<\infty \text { for all } j \in U \tag{8}
\end{equation*}
$$

where $\|\cdot\|$ is the Euclidean norm in $\mathbb{R}^{s}$. We shall use $\|z\|_{1}=\sum_{i=1}^{k}\left|z_{i}\right|$ for $z \in \mathbb{R}^{k}$ as well. Note that $|z| \leq\|z\| \leq\|z\|_{1}$ for all $z \in \mathbb{R}^{k}$. These norms coinside if $k=1$. Moreover, all norms are equivalent in our finitedimensional case.

Set

$$
S=\sum_{j \in U} X_{j}, \quad V^{2}=\sum_{j \in U} \operatorname{Var}\left(X_{j}\right) .
$$

Suppose $\operatorname{det} V^{2}>0$ and define

$$
\begin{equation*}
Y_{j}=V^{-1} X_{j}, \quad W=\left(W_{1}, \ldots, W_{s}\right)=\sum_{j \in U} Y_{j}, \quad R=\left\|V^{-1}\right\|_{1}^{2}|U| \theta_{1} \tag{9}
\end{equation*}
$$

where $V^{-1}$ is the inverse matrix to the square root of $V^{2},\|A\|_{1}$ is the matrix norm corresponding to the vector norm $\|z\|_{1}$.

Evidently, $V^{2}, W, R$ are functions of $X_{j}, j \in U$, and we use also notation $V^{2}\left(X_{U}\right), W\left(X_{U}\right), R\left(X_{U}\right)$.

Consider a function $h: \mathbb{R}^{s} \rightarrow \mathbb{R}$ such that for some positive constants $M_{0}, M_{1}, M_{2}$ and for all $x, x^{\prime} \in \mathbb{R}^{s}, k=1, \ldots, s$, one has

$$
\begin{equation*}
|h(x)| \leq M_{0}, \quad\left|\frac{\partial h(x)}{\partial x_{k}}\right| \leq M_{1},\left|\frac{\partial h(x)}{\partial x_{k}}-\frac{\partial h\left(x^{\prime}\right)}{\partial x_{k}}\right| \leq M_{2}\left\|x-x^{\prime}\right\| . \tag{10}
\end{equation*}
$$

Using the Stein method (see Stein $(1972,1986)$ ) we provide, for a $(B L, \theta)$ dependent random field $X=\left\{X_{j}, j \in \mathbb{Z}^{d}\right\}$ with values in $\mathbb{R}^{s}$, upper estimates of a functional

$$
\begin{equation*}
\Delta\left(h, X_{U}\right)=|\mathrm{E} h(W)-\mathrm{E} h(Z)| \tag{11}
\end{equation*}
$$

where $h$ is a function satisfying conditions (10), $Z$ is a standard normal vector in $\mathbb{R}^{s}$ and $U$ is a finite subset of $\mathbb{Z}^{d}$.

It is worth remarking that there are various generalizations of the Stein method. We refer to Chen (1975), Tikhomirov (1980, 1983), Barbour (1990), Götze(1991). The approach based on diffusion approximation for positively or negatively dependent random field was used in Bulinski and Shabanovich (1998). Interesting applications of the Stein techniques (with semigroup approach) in the framework of statistical models are discussed in Baddeley (2000).

Our main result (Theorem 1) gives an estimate for $\Delta\left(h, X_{U}\right)$ in terms of the Lindeberg function

$$
\begin{equation*}
\mathcal{L}_{\varepsilon}=\mathcal{L}_{\varepsilon}\left(X_{U}\right)=\sum_{j \in U} E\left\|Y_{j}\right\|^{2} \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\}, \quad \varepsilon>0 \tag{12}
\end{equation*}
$$

of the function $R$ appearing in (9); here random vectors $Y_{j}$ are defined in (9) and $\mathbf{1}\{\cdot\}$ is an indicator function.

If, moreover, for some $\delta \in(0,1]$

$$
\begin{equation*}
\mathrm{E}\left\|X_{j}\right\|^{2+\delta}<\infty, \quad j \in U \tag{13}
\end{equation*}
$$

then Theorem 2 gives an estimate of $\Delta\left(h, X_{U}\right)$ in terms of the Lyapunov fraction instead of $\mathcal{L}_{\varepsilon}$.

Using the smoothing techniques we establish (Theorem 3) the upper bound for

$$
\begin{equation*}
\Delta\left(B, X_{U}\right)=|\mathrm{P}(W \in B)-\mathrm{P}(Z \in B)| \tag{14}
\end{equation*}
$$

where $B$ is an arbitrary convex set in $\mathbb{R}^{s}$.
It is also shown that the Bernstein block techniques is useful in combination with the above mentioned theorems (see Theorem 4).

An application to kernel estimates of unknown density of a vector-valued stationary random field is provided as well. Theorem 5 extends some results obtained by Bosq et al. (1999), Roussas (2000, 2001), Veretennikov (2000), Bulinski and Millionshchikov (2002).

## 2. Results and proofs

Here we keep the notation used in the Introduction.
Theorem 1. Let $X=\left\{X_{j}, j \in \mathbb{Z}^{d}\right\}$ be a random field with values in $\mathbb{R}^{s}$ satisfying condition (8) where $U$ is a finite subset of $\mathbb{Z}^{d}$. Assume that for a function $h$ condition (10) holds. Then for every $\varepsilon>0$

$$
\begin{align*}
& \Delta\left(h, X_{U}\right) \leq s\left(D_{1}+\varepsilon(s+1) D_{2}\right) R+2 \varepsilon c(s) D_{2} \\
& +\left(2 \varepsilon^{-1} D_{0}+(6 s+1) D_{1}+(\varepsilon / 2) s(s+1) D_{2}\right) \mathcal{L}_{\varepsilon} \tag{15}
\end{align*}
$$

where $c(s)=\sum_{k=1}^{s} k^{3 / 2} \leq s^{5 / 2}$,

$$
\begin{align*}
& D_{0}=\sqrt{2 \pi} M_{0}, \quad D_{1}=\max \left\{4 M_{0}, \sqrt{2 \pi} M_{1}\right\} \\
D_{2}= & \sqrt{2} \max \left\{\sqrt{2 \pi} M_{0}+2 M_{1}, 4 \sqrt{s} M_{1}, \sqrt{2 \pi} M_{2}\right\} \tag{16}
\end{align*}
$$

and the constants $M_{0}, M_{1}, M_{2}$ appear in (10).
Proof. For $i=1, \ldots, s(s \geq 1)$ and $x_{i}, \ldots, x_{s} \in \mathbb{R}$ introduce functions

$$
\begin{equation*}
H_{i}\left(x_{i}, \ldots, x_{s}\right)=\mathrm{E}\left(h\left(Z_{1}, \ldots, Z_{i-1}, x_{i}, \ldots, x_{s}\right)-h\left(Z_{1}, \ldots, Z_{i}, x_{i+1}, \ldots x_{s}\right)\right), \tag{17}
\end{equation*}
$$

here $Z=\left(Z_{1}, \ldots, Z_{s}\right)$ is a standard normal vector in $\mathbb{R}^{s}$ (as usual if $s=1$ one has $H_{1}\left(x_{1}\right)=\mathrm{E}\left(h\left(x_{1}\right)-h\left(Z_{1}\right)\right)$, if $s \geq 2$ then $H_{1}\left(x_{1}, \ldots, x_{s}\right)=\mathrm{E}\left(h\left(x_{1}, \ldots, x_{s}\right)-\right.$ $\left.h\left(Z_{1}, x_{2}, \ldots, x_{s}\right)\right)$ and $\left.H_{s}\left(x_{s}\right)=\mathrm{E}\left(h\left(Z_{1}, \ldots, Z_{s-1}, x_{s}\right)-h\left(Z_{1}, \ldots, Z_{s}\right)\right)\right)$.

For $i=1, \ldots, s$ consider a differential equation

$$
\begin{equation*}
\frac{\partial f_{i}}{\partial x_{i}}-x_{i} f_{i}=H_{i} \tag{18}
\end{equation*}
$$

where functions $H_{i}$ are defined in (17).
Below we employ the solution of this equation given by the formula

$$
f_{i}=f_{i}\left(x_{i}, \ldots, x_{s}\right)=e^{x_{i}^{2} / 2} \int_{-\infty}^{x_{i}} H_{i}\left(u, x_{i+1}, \ldots, x_{s}\right) e^{-u^{2} / 2} d u
$$

(for $i=s$ one has $f_{s}=f_{s}\left(x_{s}\right)=e^{x_{s}^{2} / 2} \int_{-\infty}^{x_{s}} H_{s}(u) e^{-u^{2} / 2} d u$ ).
Lemma 1. For all $x=\left(x_{i}, \ldots, x_{s}\right), x^{\prime}=\left(x_{i}^{\prime}, \ldots, x_{s}^{\prime}\right) \in \mathbb{R}^{s-i+1}$ and any $i=1, \ldots, s, k=i, \ldots, s$ the following inequalities are valid

$$
\begin{equation*}
\left|f_{i}(x)\right| \leq D_{0}, \quad\left|\partial f_{i}(x) / \partial x_{k}\right| \leq D_{1} \tag{19}
\end{equation*}
$$

$$
\begin{equation*}
\left|\partial f_{i}(x) / \partial x_{k}-\partial f_{i}\left(x^{\prime}\right) / \partial x_{k}\right| \leq D_{2}\left\|x-x^{\prime}\right\| \tag{20}
\end{equation*}
$$

where $D_{0}, D_{1}, D_{2}$ are indicated in (16).
Proof. Note that $\mathrm{E} h\left(Z_{1}, \ldots, Z_{i}, x_{i+1}, \ldots, x_{s}\right)$ for a Borel function $h: \mathbb{R}^{s} \rightarrow$ $\mathbb{R}$ is given by the expression

$$
(2 \pi)^{-i / 2} \int_{\mathbb{R}^{i}} e^{-\frac{u_{1}^{2}+\ldots+u_{i}^{2}}{2}} h\left(u_{1}, \ldots, u_{i}, x_{i+1}, \ldots, x_{s}\right) d u_{1} \ldots d u_{i}
$$

and, for $h$ having bounded partial derivatives in $x_{i+1}, \ldots, x_{s}$,

$$
\frac{\partial}{\partial x_{k}} \mathrm{E} h\left(Z_{1}, \ldots, Z_{i}, x_{i+1}, \ldots, x_{s}\right)=\mathrm{E} \frac{\partial}{\partial x_{k}} h\left(Z_{1}, \ldots, Z_{i}, x_{i+1}, \ldots, x_{s}\right)
$$

for every $k=i+1, \ldots, s$. A simple calculation shows that, if $K: \mathbb{R} \rightarrow \mathbb{R}$ is a bounded Borel function, $|K(x)| \leq K_{0}, x \in \mathbb{R}$, and $\int_{\mathbb{R}} K(u) e^{-u^{2} / 2} d u=0$, then for all $x \in \mathbb{R}$

$$
\begin{equation*}
e^{x^{2} / 2}\left|\int_{-\infty}^{x} K(u) e^{-u^{2} / 2} d u\right| \leq K_{0} \sqrt{\pi / 2}, \quad|x| e^{x^{2} / 2}\left|\int_{-\infty}^{x} K(u) e^{-u^{2} / 2} d u\right| \leq K_{0} \tag{21}
\end{equation*}
$$

Thus to establish (19) we use the upper bounds for the absolute values of functions appearing under the signs of integrals in representations for $f_{i}\left(x_{i}, \ldots, x_{s}\right)$ and $\partial f_{i}\left(x_{i}, \ldots, x_{s}\right) / \partial x_{k}$.

To obtain (20) consider separately the cases $k>i$ and $k=i$. Moreover, each time we have to consider whether $x_{i}=x_{i}^{\prime}$ or $x_{k}=x_{k}^{\prime}, k=i+1, \ldots, s$. As for the case of $x_{k}=x_{k}^{\prime}, k=i+1, \ldots, s$, let us note the existence of the second partial derivatives $\partial^{2} f_{i} / \partial x_{i} \partial x_{k}, k=i, \ldots, s$. Clearly from the equation (18) we have

$$
\begin{gathered}
\partial^{2} f_{i} / \partial x_{i} \partial x_{k}=\partial H_{i} / \partial x_{k}+x_{i} \partial f_{i} / \partial x_{k}, k=i+1, \ldots, s ; \\
\partial^{2} f_{i} / \partial x_{i}^{2}=\partial H_{i} / \partial x_{i}+\left(1+x_{i}^{2}\right) f_{i}+x_{i} H_{i} .
\end{gathered}
$$

The estimate for the second partial derivative in $x_{i}$ and $x_{k}, k>i$, follows now from (21). To estimate $\partial^{2} f_{i} / \partial x_{i}^{2}$ it suffices to integrate by parts in the integral representation for $x_{i}^{2} f_{i}$. In the case $x_{i}=x_{i}^{\prime}$ we use again the representations for $f_{i}$ and $\partial f_{i} / \partial x_{k}, k=i, \ldots, s$. The Lemma is proved.

Continuing the proof of Theorem 1, observe that (17) and (18) imply

$$
\begin{equation*}
\sum_{i=1}^{s} \mathrm{E}\left(\frac{\partial}{\partial x_{i}}-W_{i}\right) f_{i}\left(W_{i}, \ldots, W_{s}\right)=\mathrm{E} h(W)-\mathrm{E} h(Z) \tag{22}
\end{equation*}
$$

where the vector $W$ is defined in (9).
For each fixed $i \in\{1, \ldots, s\}$ we estimate the summand

$$
\mathrm{E}\left(\frac{\partial}{\partial x_{i}}-W_{i}\right) f_{i}\left(W_{i}, \ldots, W_{s}\right)
$$

in the left-hand side of (22). Analogously to Bulinski and Suquet (2001) introduce for a given $\varepsilon>0$ auxiliary random vectors

$$
T_{j}=\left(T_{j 1}, \ldots, T_{j s}\right)=\left(b\left(Y_{j 1}\right), \ldots, b\left(Y_{j s}\right)\right), \quad V_{j}=\left(V_{j 1}, \ldots, V_{j s}\right)=Y_{j}-T_{j}
$$

where $b(y)=\operatorname{sign}(y) \min \{|y|, \varepsilon\}, y \in \mathbb{R}$. For the sake of brevity we write $\mathbb{W}=\left(W_{i}, \ldots, W_{s}\right)$ and $\mathbb{T}_{j}=\left(T_{j i}, \ldots, T_{j s}\right)$. Set

$$
\mathbb{W}^{(j)}=\mathbb{W}-\left(Y_{j i}, \ldots, Y_{j s}\right), \quad j \in U .
$$

It can be seen that

$$
\mathrm{E} W_{i} f_{i}(\mathbb{W})=\sum_{q=1}^{4} R_{i q}
$$

where

$$
\begin{gathered}
R_{i 1}=\sum_{j \in U} \mathrm{E} Y_{j i} f_{i}\left(\mathbb{W}^{(j)}\right), \\
R_{i 2}=\sum_{j \in U} \mathrm{E} V_{j i}\left(f_{i}(\mathbb{W})-f_{i}\left(\mathbb{W}^{(j)}\right)\right), \\
R_{i 3}=\sum_{j \in U} \mathrm{E} T_{j i}\left(f_{i}(\mathbb{W})-f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)\right), \\
R_{i 4}=\sum_{j \in U} \mathrm{E} T_{j i}\left(f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)-f_{i}\left(\mathbb{W}^{(j)}\right)\right) .
\end{gathered}
$$

Note that for a Lipschitz function $G: \mathbb{R}^{m} \rightarrow \mathbb{R}$ and a linear map $A$ : $\mathbb{R}^{n} \rightarrow \mathbb{R}^{m}(m, n \in \mathbb{N})$ the composition $G(A(\cdot))$ is a Lipschitz function with $\operatorname{Lip}(G A) \leq \operatorname{Lip}(G)\|A\|_{1}$. Using this fact, definitions (7), (12) and inequalities (19) we derive the following estimates

$$
\begin{gathered}
\left|R_{i 1}\right| \leq \sum_{j \in U}\left|\operatorname{cov}\left(Y_{j i}, f_{i}\left(\mathbb{W}^{(j)}\right)\right)\right| \leq \sum_{j \in U}\left|\operatorname{cov}\left(T_{j i}, f_{i}\left(\mathbb{W}^{(j)}\right)\right)\right|+2 D_{0} \sum_{j \in U} \mathrm{E}\left|V_{j i}\right| \\
\leq D_{1}\left\|V^{-1}\right\|_{1}^{2}|U| \theta_{1}+2 D_{0} \varepsilon^{-1} \sum_{j \in U} \mathrm{E} Y_{j i}^{2} \boldsymbol{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\}
\end{gathered}
$$

$$
\begin{gathered}
=D_{1} R+2 D_{0} \varepsilon^{-1} \sum_{j \in U} \mathrm{E} Y_{j i}^{2} \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\} \\
\left|R_{i 2}\right| \leq \sum_{j \in U}\left|\mathrm{E} V_{j i}\left(f_{i}(\mathbb{W})-f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)\right)\right|+\left|\mathrm{E} V_{j i}\left(f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)-f_{i}\left(\mathbb{W}^{(j)}\right)\right)\right| \\
\leq D_{1} \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E}\left(\left|V_{j i} V_{j k}\right|+\left|V_{j i} T_{j k}\right|\right) \\
\leq D_{1} \sum_{j \in U} \sum_{k=i}^{s}\left(\frac{1}{2} \mathrm{E}\left(Y_{j i}^{2}+Y_{j k}^{2}\right) \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\}+\varepsilon \mathrm{E}\left|Y_{j i}\right| \mathbf{1}\left\{\left|Y_{j i}\right|>\varepsilon\right\}\right) \\
\leq \frac{3}{2} D_{1}(s-i+1) \sum_{j \in U} \mathrm{E} Y_{j i}^{2} \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\}+\frac{1}{2} D_{1} \mathcal{L}_{\varepsilon} .
\end{gathered}
$$

In a similar way

$$
\left|R_{i 3}\right| \leq D_{1} \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E}\left|T_{j i} V_{j k}\right| \leq D_{1} \sum_{j \in U} \sum_{k=i}^{s} \varepsilon \mathrm{E}\left|Y_{j k}\right| \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\} \leq D_{1} \mathcal{L}_{\varepsilon}
$$

Due to differentiability of $f_{i}$ one has

$$
\begin{gathered}
\left|f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)-f_{i}\left(\mathbb{W}^{(j)}\right)-\sum_{k=i}^{s} \frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}} T_{j k}\right| \\
\leq \sum_{k=i}^{s}\left|\frac{\partial f_{i}\left(\mathbb{W}^{(j)}+\tau_{0} \mathbb{T}_{j}\right)}{\partial x_{k}}-\frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}}\right|\left|T_{j k}\right| \leq D_{2}(s-i+1)^{1 / 2} \sum_{k=i}^{s} \max _{i \leq p \leq s}\left|T_{j p}\right|\left|T_{j k}\right|
\end{gathered}
$$

where $\tau_{0} \in[0,1]$. Taking into account the relation $\sum_{j \in U} \operatorname{Var}\left(Y_{j}\right)=I$ where $I$ is a unit matrix of order $s$ we conclude that

$$
R_{i 4}=\sum_{j \in U} \sum_{k=i}^{s} \mathrm{E} T_{j i} T_{j k} \frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}}+\triangle_{i 1}
$$

One has

$$
\begin{gathered}
\left|\triangle_{i 1}\right| \leq \varepsilon D_{2}(s-i+1)^{1 / 2} \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E}\left|T_{j i} T_{j k}\right| \\
\leq \varepsilon D_{2}(s+i-1)^{1 / 2} \sum_{k=i}^{s} \sum_{j \in U} \mathrm{E}\left(Y_{j i}^{2}+Y_{j k}^{2}\right) / 2 \leq \varepsilon D_{2}(s-i+1)^{3 / 2}
\end{gathered}
$$

and

$$
\begin{aligned}
& \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E} T_{j i} T_{j k} \frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}}=\sum_{j \in U} \sum_{k=i}^{s} \operatorname{cov}\left(T_{j i} T_{j k}, \frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}}\right) \\
&+ \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E} T_{j i} T_{j k} \mathrm{E}\left(\frac{\partial f_{i}(\mathbb{W}(j)}{\partial x_{k}}-\frac{\partial f_{i}(\mathbb{W})}{\partial x_{k}}\right) \\
&+\sum_{j \in U} \sum_{k=i}^{s} \mathrm{E} T_{j i} T_{j k} \mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{k}}=\sum_{q=1}^{3} C_{i q}
\end{aligned}
$$

For any $i, k=1, \ldots, s$ the function $b_{i k}(x)=b\left(x_{i}\right) b\left(x_{k}\right), x=\left(x_{1}, \ldots, x_{s}\right) \in \mathbb{R}^{s}$, is a Lipschitz one with $\operatorname{Lip}\left(b_{i k}\right) \leq 2 \varepsilon$. Consequently

$$
\begin{gathered}
\left|C_{i 1}\right| \leq \sum_{j \in U} \sum_{k=i}^{s}\left|\operatorname{cov}\left(T_{j i} T_{j k}, \frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}}\right)\right| \leq 2 D_{2} \varepsilon\left\|V^{-1}\right\|_{1}^{2}(s-i+1)|U| \theta_{1} \\
\leq 2 D_{2}(s-i+1) \varepsilon R . \\
\left|C_{i 2}\right| \leq\left|\sum_{j \in U} \sum_{k=i}^{s} \mathrm{E} T_{j i} T_{j k} \mathrm{E}\left(\frac{\partial f_{i}\left(\mathbb{W}^{(j)}\right)}{\partial x_{k}}-\frac{\partial f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)}{\partial x_{k}}\right)\right| \\
+\left|\sum_{j \in U} \sum_{k=i}^{s} \mathrm{E} T_{j i} T_{j k} \mathrm{E}\left(\frac{\partial f_{i}\left(\mathbb{W}^{(j)}+\mathbb{T}_{j}\right)}{\partial x_{k}}-\frac{\partial f_{i}(\mathbb{W})}{\partial x_{k}}\right)\right| \\
\leq D_{2}(s-i+1)^{1 / 2} \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E}\left|T_{j i} T_{j k}\right| \mathrm{E} \max _{p=i, \ldots, s}\left|T_{j p}\right|+D_{2} \sum_{j \in U} \sum_{k=i}^{s} \sum_{p=i}^{s} \mathrm{E}\left|T_{j i} T_{j k}\right| \mathrm{E}\left|V_{j p}\right| \\
\leq \varepsilon D_{2}(s-i+1)^{3 / 2}+\varepsilon D_{2}(s-i+1) \mathcal{L}_{\varepsilon} .
\end{gathered}
$$

To estimate $C_{i 3}$ consider at first the case $k=i$. Then

$$
\begin{aligned}
& \sum_{j \in U} \mathrm{E} T_{j i}^{2} \mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{i}}=\mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{i}} \sum_{j \in U} \mathrm{E} Y_{j i}^{2}+\triangle_{i 2}=\mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{i}}+\triangle_{i 2}, \\
& \left|\triangle_{i 2}\right|=\left|\mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{i}} \sum_{j \in U}\left(\mathrm{E} T_{j i}^{2}-\mathrm{E} Y_{j i}^{2}\right)\right| \leq D_{1} \sum_{j \in U} \mathrm{E} Y_{j i}^{2} \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\} .
\end{aligned}
$$

The bounds for $R_{i 2}$ and $R_{i 3}$ imply that

$$
\begin{gathered}
\left|\triangle_{i 3}\right|=\left|\sum_{j \in U} \sum_{k=i+1}^{s} \mathrm{E} T_{j i} T_{j k} \mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{k}}\right| \\
=\left|\sum_{k=i+1}^{s} \mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{k}} \sum_{j \in U} \mathrm{E}\left(Y_{j i} Y_{j k}-T_{j i} V_{j k}-V_{j i} T_{j k}-V_{j i} V_{j k}\right)\right| \\
\leq \sum_{j \in U} \sum_{k=i+1}^{s}\left|\mathrm{E} \frac{\partial f_{i}(\mathbb{W})}{\partial x_{k}}\right|\left(\mathrm{E}\left|T_{j i} V_{j k}\right|+\mathrm{E}\left|V_{j i} T_{j k}\right|+\mathrm{E}\left|V_{j i} V_{j k}\right|\right) \\
\leq D_{1}\left((3 / 2) \mathcal{L}_{\varepsilon}+(3 s / 2) \sum_{j \in U} \mathrm{E} Y_{j i}^{2} \mathbf{1}\left\{\left\|Y_{j}\right\|>\varepsilon\right\}\right)
\end{gathered}
$$

in view of the relation $\sum_{j \in U} \operatorname{Var}\left(Y_{j}\right)=I$, that is $\sum_{j \in U} \mathrm{E} Y_{j i} Y_{j k}=0, k \neq i$.
Finally we have

$$
\begin{gathered}
\Delta\left(h, X_{U}\right) \leq \sum_{i=1}^{s}\left|\mathrm{E}\left(\frac{\partial}{\partial x_{i}}-W_{i}\right) f_{i}\left(W_{i}, \ldots, W_{s}\right)\right| \\
\leq \sum_{i=1}^{s}\left(\sum_{q=1}^{3}\left(\left|R_{i q}\right|+\left|\triangle_{i q}\right|\right)+\sum_{q=1}^{2}\left|C_{i q}\right|\right) .
\end{gathered}
$$

Observing that $\sum_{i=1}^{s}(s-i+1)^{3 / 2}=c(s)$ and $\sum_{i=1}^{s}(s-i+1)=s(s+1) / 2$ we come to (15). The proof of Theorem 1 is complete.

Corollary 1. For a family of centered random fields $X^{(n)}=\left\{X_{j}^{(n)}, j \in \mathbb{Z}^{d}\right\}$ $(n \in \mathbb{N})$ with values in $\mathbb{R}^{s}$ and a family of finite subsets $U_{n}$ of $\mathbb{Z}^{d}$, the CLT holds, that is

$$
W\left(X_{U_{n}}^{(n)}\right) \xrightarrow{\text { Law }} Z \text { as } n \rightarrow \infty,
$$

whenever, for every $\varepsilon>0$,

$$
\mathcal{L}_{\varepsilon}\left(X_{U_{n}}^{(n)}\right) \rightarrow 0 \text { and } R\left(X_{U_{n}}^{(n)}\right) \rightarrow 0 \text { as } n \rightarrow \infty
$$

If $X_{U_{n}}^{(n)}$ consists of independent random vectors then $R\left(X_{U_{n}}^{(n)}\right)=0$. Thus Theorem 1 comprises the multidimensional Lindeberg theorem for independent summands. Analogously to Bulinski and Suquet (2001) one can also obtain from (15) the generalization of the classical Newman CLT for associated random fields.

Theorem 2. Assume that conditions of Theorem 1 are satisfied and, moreover, (13) holds. Then

$$
\begin{array}{r}
\Delta\left(h, X_{U}\right) \leq s\left(D_{1}+D_{2}(s+1)\right) R  \tag{23}\\
+\left(2 D_{0}+(6 s+1) D_{1}+(2 c(s)+s(s+1) / 2) D_{2}\right) L_{2+\delta}
\end{array}
$$

where the Lyapunov fraction

$$
L_{2+\delta}=L_{2+\delta}\left(X_{U}\right)=\sum_{j \in U} \mathrm{E}\left\|Y_{j}\right\|^{2+\delta} \leq\left\|V^{-1}\right\|_{1}^{2+\delta} \sum_{j \in U} \mathrm{E}\left\|X_{j}\right\|_{1}^{2+\delta},
$$

random vectors $Y_{j}$ are defined in (9) and $c(s)$ appears in (15).
Proof. Observe that for $\varepsilon=1$ and $\delta \in(0,1]$ one has $\mathcal{L}_{\varepsilon} \leq L_{2+\delta}$. To estimate $\Delta_{i 1}$ (and analogously $C_{i 2}$ ) we use inequalities $|a b c| \leq\left(|a|^{3}+|b|^{3}+\right.$ $\left.|c|^{3}\right) / 3$ for all $a, b, c \in \mathbb{R}$ and $|a|^{3} \leq|a|^{2+\delta}$ when $|a| \leq 1$. Thus

$$
\begin{gathered}
D_{2}(s-i+1)^{1 / 2} \sum_{j \in U} \sum_{k=i}^{s} \mathrm{E}\left|T_{j i} T_{j k} \max _{p=i, \ldots, s} T_{j p}\right| \\
\leq \frac{1}{3} D_{2}(s+i-1)^{1 / 2} \sum_{j \in U} \sum_{k=i}^{s}\left(\mathrm{E}\left|T_{j i}\right|^{3}+\mathrm{E}\left|T_{j k}\right|^{3}+\mathrm{E} \max _{p=i, \ldots, s}\left|T_{j p}\right|^{3}\right) \\
\leq D_{2}(s-i+1)^{3 / 2} \sum_{j \in U} \mathrm{E}\left\|Y_{j}\right\|^{2+\delta}=D_{2}(s-i+1)^{3 / 2} L_{2+\delta} .
\end{gathered}
$$

Theorem is proved.
Now we need some new notation. Let $B^{(\gamma)}$ be a $\gamma$-neighborhood of a set $B \subset \mathbb{R}^{s}$ with respect to Euclidean distance (that is $B^{(\gamma)}=\left\{x \in \mathbb{R}^{s}\right.$ : $\left.\left.\inf _{y \in B}\|x-y\|<\gamma\right\}\right), \partial B$ being the boundary of $B$.

Remark 3. From Theorem 1.4 by Goldstein and Rinott (1996) the estimate for $\Delta\left(h, X_{U}\right)$ (in our notation) can be established when $h \in C_{b}^{3}\left(\mathbb{R}^{s}\right)$ and $\mathrm{E} X_{j}=0, \mathrm{E}\left\|X_{j}\right\|^{4}<\infty, j \in U$. Using a function $h \in C_{b}^{3}$ approximating the indicator function of a convex set $B \subset \mathbb{R}^{s}$ (more precisely, for a given $\gamma \in(0,1)$ let $h(x)=1$ for $x \in B, h(x)=0$ for $x \notin B^{(\gamma)}$ and $0 \leq h(x) \leq 1$ for all $x \in \mathbb{R}^{s}$ ) one can derive from the mentioned estimate that

$$
\begin{equation*}
\Delta\left(B, X_{U}\right) \leq \mathrm{P}\left(Z \in(\partial B)^{(\gamma)}\right)+G_{\gamma}\left(X_{U}\right) \tag{24}
\end{equation*}
$$

where $\Delta\left(B, X_{U}\right)$ is defined in (15), $G_{\gamma}(\cdot)$ is a certain (nonrandom) functional in $X_{U}$. In Theorems 1 and 2 of our paper the estimates of $\Delta\left(h, X_{U}\right)$ are
obtained in other terms under lower moment assumptions and for a wider class of functions $h$ satisfying the conditions (10). We have

$$
\begin{equation*}
\Delta\left(B, X_{U}\right) \leq \mathrm{P}\left(Z \in(\partial B)^{(\gamma)}\right)+H_{\gamma}\left(X_{U}\right) \tag{25}
\end{equation*}
$$

where $H_{\gamma}(\cdot)$ is a specified nonrandom functional in $X_{U}$, as the next theorem shows. For fixed $U, \varepsilon$ and $s$ one has $G_{\gamma}\left(X_{U}\right)=O\left(\gamma^{-3}\right)$ as $\gamma \rightarrow 0$ whereas $H_{\gamma}\left(X_{U}\right)=O\left(\gamma^{-2}\right)$ as $\gamma \rightarrow 0$.

Write

$$
\gamma_{0}(s)=\min \{1,3 \sqrt{\pi /(2 s)}\}
$$

Theorem 3. Let conditions of Theorem 1 be satisfied and $B$ be a convex set in $\mathbb{R}^{s}$. Then for any $\gamma \in\left(0, \gamma_{0}(s)\right]$ the estimate (25) holds with

$$
\begin{align*}
H_{\gamma}\left(X_{U}\right)= & \sqrt{2 \pi} s \gamma^{-1}\left(2+12 \sqrt{2}(s+1) \varepsilon \gamma^{-1}\right) R+48 \sqrt{\pi} c(s) \varepsilon \gamma^{-2}  \tag{26}\\
& +\sqrt{2 \pi}\left(\varepsilon^{-1}+(12 s+2) \gamma^{-1}+6 \sqrt{2} s(s+1) \varepsilon \gamma^{-2}\right) \mathcal{L}_{\varepsilon} .
\end{align*}
$$

If, moreover, conditions of Theorem 2 are satisfied then one can take

$$
\begin{array}{r}
H_{\gamma}\left(X_{U}\right)=\sqrt{2 \pi} s \gamma^{-1}\left(2+12 \sqrt{2}(s+1) \gamma^{-1}\right) R \\
+\sqrt{2 \pi}\left(1+(12 s+2) \gamma^{-1}+12 \sqrt{2}(2 c(s)+s(s+1) / 2) \gamma^{-2}\right) L_{2+\delta} . \tag{27}
\end{array}
$$

Proof. Introduce a function $\psi$ setting

$$
\psi(x)= \begin{cases}1, & x \leq 0,  \tag{28}\\ 1-\frac{16 x^{3}}{3 \gamma^{3}}, & x \in\left(0, \frac{\gamma}{4}\right], \\ \frac{3}{2}-\frac{2 x}{\gamma}-\frac{16}{3 \gamma^{3}}\left(\frac{\gamma}{2}-x\right)^{3}, & x \in\left(\frac{\gamma}{4}, \frac{3 \gamma}{4}\right], \\ \frac{16}{3 \gamma^{3}}(\gamma-x)^{3}, & x \in\left(\frac{3 \gamma}{4}, \gamma\right], \\ 0, & x \geq \gamma .\end{cases}
$$

It is easy to verify that the following statement is true.
Lemma 2. The function $\psi \in C^{2}(\mathbb{R})$ and for all $u \in \mathbb{R}$ one has

$$
0 \leq \psi(u) \leq 1, \quad\left|\psi^{\prime}(u)\right| \leq 2 \gamma^{-1}, \quad\left|\psi^{\prime \prime}(u)\right| \leq 8 \gamma^{-2}
$$

For a convex set $B \subset \mathbb{R}^{s}$ define the function

$$
\begin{equation*}
h(x)=\psi(\rho(x, B)), \quad x \in \mathbb{R}^{s} \tag{29}
\end{equation*}
$$

where $\psi$ is given by (28) and $\rho$ is the Euclidean distance in $\mathbb{R}^{s}$.
Obviously

$$
\begin{equation*}
0 \leq h(x) \leq 1, x \in \mathbb{R}^{s}, \quad h(x)=1, x \in B, \quad h(x)=0, x \notin B^{(\gamma)} . \tag{30}
\end{equation*}
$$

Lemma 3. For the nonnegative function $h \in C^{1}\left(\mathbb{R}^{s}\right)$ (see (29)) the condition (10) is satisfied with $M_{0}=1, M_{1}=2 \gamma^{-1}$ and $M_{2}=12 \gamma^{-2}$.

To prove this result one can use the properties of $\psi$ given in Lemma 2 and take into account that for all $x \notin[B]$ and any $i=1, \ldots, s$ there exists

$$
\frac{\partial}{\partial x_{i}} \rho(x, B)=-\cos \left(e_{i}, n\right) .
$$

Here $e_{i}$ is the $i$-th unit vector of the natural orthonormal basis of $\mathbb{R}^{s}, n=y-x$ where $y \in[B]$ and $\|x-y\|=\rho(x, B) ;[B]$ is a closure of $B$ in the Euclidean distance.

Lemma 4. If $\gamma$ satisfies the conditions of Theorem 3 and the function $h$ is given by (29), the statement of Lemma 1 is valid with

$$
\begin{equation*}
D_{0}=\sqrt{\pi / 2}, \quad D_{1}=2 \sqrt{2 \pi} \gamma^{-1}, \quad D_{2}=24 \sqrt{\pi} \gamma^{-2} \tag{31}
\end{equation*}
$$

Proof. The indicated values for $D_{0}, D_{1}, D_{2}$ can be easily obtained from the formula (19), analogously to the proof of Lemma 1, using the fact that $\left|H_{i}(x)\right| \leq 1$ since $0 \leq h(x) \leq 1, x \in \mathbb{R}^{s}$.

Now we proceed with the proof of Theorem 3. Due to Theorem 1 and Lemma 4 we come to the estimate (15) with $D_{0}, D_{1}, D_{2}$ indicated in (31). In view of (30)

$$
\begin{equation*}
\mathrm{P}(W \in B)-\mathrm{P}\left(Z \in B^{(\gamma)}\right) \leq \mathrm{E} h(W)-\mathrm{E} h(Z) \leq \Delta\left(h, X_{U}\right) \tag{32}
\end{equation*}
$$

In a similar way for the set $B_{(\gamma)}=B \backslash(\partial B)^{(\gamma)}$ one has

$$
\begin{equation*}
\mathrm{P}(W \in B)-\mathrm{P}\left(Z \in B_{(\gamma)}\right) \geq-\Delta\left(h, X_{U}\right) . \tag{33}
\end{equation*}
$$

Now (32) and (33) imply (25) with $H\left(X_{U}\right)$ given by (26). The second assertion of Theorem 3 follows in the same manner as Theorem 2 was obtained using Theorem 1. Theorem 3 is proved.

The next result gives the possibility to provide an estimate for $\Delta\left(B, X_{U}\right)$ which is uniform on the class $\mathcal{C}_{s}$ of convex sets of $\mathbb{R}^{s}$. It is easy to derive immediately the following bound from Corollary 3.2 of Bhattacharia and Ranga Rao (1976).

Lemma 5. For any $k \in \mathbb{N}$, all $\gamma>0$ and every convex set $B \subset \mathbb{R}^{s}$

$$
\mathrm{P}\left(Z \in(\partial B)^{\gamma}\right) \leq a(s) \gamma
$$

where $a(1)=2 \sqrt{2 / \pi}$ and $a(s)=2^{1 / 2}(s-1) \Gamma((s-1) / 2) / \Gamma(s / 2)$ for $s \geq 2$. Thus

$$
\begin{equation*}
a(s) \leq a_{0} s^{1 / 2}, \quad a_{0}=\text { const }, \quad s \in \mathbb{N} . \tag{34}
\end{equation*}
$$

Corollary 2. Let conditions of Theorem 2 be satisfied. Then

$$
\sup _{B \in \mathcal{C}_{s}} \Delta\left(B, X_{U}\right) \leq c\left(R\left(X_{U}\right)+L_{2+\delta}\left(X_{U}\right)\right)^{1 / 3}
$$

where a factor $c=c(s)$.
Remark 4. We are interested in asymptotical behaviour of random vectors $W=W\left(X_{U}\right)$ as $U \rightarrow \infty$ in a sense. In this regard note that in general $R=R\left(X_{U}\right)$ need not tend to zero for growing sets $U$ (if for dependent summands there are points of $U$ which are "rather close" to each other). So, it is natural to use the combination of the obtained results with the Bernstein block techniques. Our next two theorems provide examples of this approach.

Let $X=\left\{X_{j}, j \in \mathbb{Z}^{d}\right\}$ be a $(B L, \theta)$-dependent random field with values in $\mathbb{R}^{s}$ such that (4) holds with $a(I, J)=\min \{|I|,|J|\}$ and some function $u(r) \searrow 0$ as $r \rightarrow \infty$. Define for a Lipschitz function $\varphi: \mathbb{R}^{s} \rightarrow \mathbb{R}^{m}$ (in Euclidean spaces we use the norm $\|\cdot\|_{1}$ ) a random field $\widetilde{X}=\left\{\widetilde{X}_{j}, j \in \mathbb{Z}^{d}\right\}$ where $\widetilde{X}_{j}=\varphi\left(X_{j}\right), j \in \mathbb{Z}^{d}$.

Lemma 6. A random field $\widetilde{X}$ is $(B L, \widetilde{\theta})$-dependent where for any finite disjoints sets $I, J \subset \mathbb{Z}^{d}$

$$
\begin{equation*}
\widetilde{\theta}(I, J) \leq \min \{|I|,|J|\}(\operatorname{Lip}(\varphi))^{2} u(\rho(I, J)) \tag{35}
\end{equation*}
$$

Proof. Note that if $f_{k}\left(f_{k-1}\left(\ldots f_{1}\right)\right)$ is a composition of Lipschitz functions $f_{1}, \ldots, f_{k}$ then

$$
\operatorname{Lip}\left(f_{k}\left(f_{k-1}\left(\ldots f_{1}\right)\right)\right) \leq \operatorname{Lip}\left(f_{k}\right) \ldots \operatorname{Lip}\left(f_{1}\right)
$$

Now (35) is obvious due to ( $B L, \theta$ )-dependence of a field $X$.
Let $U$ be a finite subset of $\mathbb{Z}^{d}$ such that

$$
\begin{equation*}
U=\bigcup_{k=0}^{N} U^{(k)}, \quad N \in \mathbb{N} \tag{36}
\end{equation*}
$$

where $U^{(0)}, \ldots, U^{(N)}$ are disjoint sets and for some positive $b$ and $q$

$$
\begin{equation*}
\left|U^{(k)}\right| \leq b \text { and } \rho\left(U^{(k)}, U^{(l)}\right) \geq q, \quad k, l=1, \ldots, N, \quad l \neq k . \tag{37}
\end{equation*}
$$

Assume that

$$
\begin{equation*}
\mathrm{E} \widetilde{X}_{j}=0 \in \mathbb{R}^{m}, \quad \mathrm{E}\left\|\widetilde{X}_{j}\right\|^{2}<\infty, \quad j \in \mathbb{Z}^{d} \tag{38}
\end{equation*}
$$

Set for $k=0, \ldots, N$

$$
\begin{equation*}
\bar{X}_{k}=\sum_{j \in U^{(k)}} \tilde{X}_{j}, \quad V_{1}^{2}=\sum_{k=1}^{N} \operatorname{Var}\left(\bar{X}_{k}\right), \quad V_{0}^{2}=\operatorname{Var}\left(\bar{X}_{0}\right) . \tag{39}
\end{equation*}
$$

Using this notation we have $\widetilde{S}=\widetilde{S}(U)=\sum_{j \in U} \widetilde{X}_{j}=\sum_{k=1}^{N} \bar{X}_{k}$. Suppose that $\operatorname{det} V_{1}^{2}>0$ and introduce for random vectors $\bar{Y}_{k}=V_{1}^{-1} \bar{X}_{k}(k=1, \ldots, N)$ the Lindeberg function

$$
\begin{equation*}
\overline{\mathcal{L}}_{\varepsilon}=\sum_{k=1}^{N} \mathrm{E}\left\|\bar{Y}_{k}\right\|^{2} \mathbf{1}\left\{\left\|\bar{Y}_{k}\right\|>\varepsilon\right\}, \quad \varepsilon>0 \tag{40}
\end{equation*}
$$

Theorem 4. Let $\widetilde{X}$ be a random field satisfying all the above mentioned conditions and $U$ be a finite set appearing in (36). Then for any nonrandom matrix $A$ of order $m$ and all $\varepsilon>0, \gamma \in\left(0, \gamma_{0}(m)\right]$ one has

$$
\begin{gather*}
\Delta:=\sup _{B \in \mathcal{C}_{m}}|\mathrm{P}(A \widetilde{S} \in B)-\mathrm{P}(Z \in B)| \\
\leq 2 a(m) \gamma+\gamma^{-2}\left\{m \sqrt{2 \pi}(2+12 \sqrt{2}(m+1) \varepsilon) N b(\operatorname{Lip}(\varphi))^{2} u(q)\left\|V_{1}^{-1}\right\|_{1}^{2}\right. \\
\left.+48 c(m) \sqrt{\pi} \varepsilon+\sqrt{2 \pi}\left(\varepsilon^{-1}+12 m+2+6 \sqrt{2} m(m+1) \varepsilon\right) \overline{\mathcal{L}}_{\varepsilon}\right\}+\gamma^{-2} \bar{\Delta} \tag{41}
\end{gather*}
$$

where

$$
\begin{equation*}
\bar{\Delta}=2 m\left\{\left\|A V_{1}-I\right\|_{1}^{2}\left(m+\left\|V_{1}^{-1}\right\|_{1}^{2} N b(\operatorname{Lip}(\varphi))^{2} u(q)\right)+m\|A\|_{1}^{2}\left\|V_{0}^{2}\right\|_{1}\right\} \tag{42}
\end{equation*}
$$

$\mathcal{C}_{m}$ is a class of convex sets in $\mathbb{R}^{m}, a(m)$ appears in (34) and $Z$ is a Gaussian vector in $\mathbb{R}^{m}$. If, moreover, for some $\delta \in(0,1]$

$$
\begin{equation*}
\sup _{j \in \mathbb{Z}^{d}} \mathrm{E}\left\|\widetilde{X}_{j}\right\|^{2+\delta} \leq \widetilde{c}<\infty \tag{43}
\end{equation*}
$$

then

$$
\Delta \leq 2 a(m) \gamma+\gamma^{-2} \sqrt{2 \pi}\left\{a_{1}(m)\left\|V_{1}^{-1}\right\|_{1}^{2} N b(\operatorname{Lip}(\varphi))^{2} u(q)\right.
$$

$$
\begin{equation*}
\left.+\widetilde{c} b^{2+\delta} a_{2}(m) N\left\|V_{1}^{-1}\right\|_{1}^{2+\delta}\right\}+\gamma^{-2} \bar{\Delta} \tag{44}
\end{equation*}
$$

where $a_{1}(m)=m(2+12 \sqrt{2}(m+1))$ and $a_{2}(m)=m^{2+(3 \delta) / 2}(12 m+3+$ $12 \sqrt{2}(2 c(m)+m(m+1) / 2))$.

Proof. We need the following two elementary results.
Lemma 7. Let $\zeta_{0}, \zeta_{1}$ be random vectors with values in $\mathbb{R}^{m}$ such that $\mathrm{E}\left\|\zeta_{0}\right\|^{2}<\infty$. Then for any value $\gamma>0$

$$
\begin{aligned}
\sup _{B \in \mathcal{C}_{m}} \mid \mathrm{P}\left(\zeta_{0}+\zeta_{1} \in B\right) & -\mathrm{P}(Z \in B)\left|\leq \sup _{B \in \mathcal{C}_{m}}\right| \mathrm{P}\left(\zeta_{1} \in B\right)-\mathrm{P}(Z \in B) \mid \\
& +\gamma a(m)+\gamma^{-2} \mathrm{E}\left\|\zeta_{0}\right\|^{2} .
\end{aligned}
$$

Lemma 8. Let $\zeta$ be a centered random vector with values in $\mathbb{R}^{m}$ such that $\mathrm{E}\|\zeta\|^{2}<\infty$. Then for any nonrandom matrix $A$ of order $m$ one has

$$
\mathrm{E}\|A \zeta\|^{2} \leq m^{2}\|A\|_{1}^{2}\|\operatorname{Var}(\zeta)\|_{1}
$$

To prove Theorem 4 note that $A \widetilde{S}=\zeta_{0}+\zeta_{1}$ where

$$
\zeta_{1}=V_{1}^{-1} \sum_{k=1}^{N} \bar{X}_{k}, \quad \zeta_{0}=\left(A V_{1}-I\right) V_{1}^{-1} \sum_{k=1}^{N} \bar{X}_{k}+A \bar{X}_{0}
$$

here $I$ is a unit matrix of order $m$.
Using estimate (49) and Lemmas 5,7 we get

$$
\begin{equation*}
\Delta \leq \sup _{B \in \mathcal{C}_{m}}\left|\mathrm{P}\left(\zeta_{1} \in B\right)-\mathrm{P}(Z \in B)\right|+\gamma a(m)+\gamma^{-2} \mathrm{E}\left\|\zeta_{0}\right\|^{2} . \tag{45}
\end{equation*}
$$

Theorem 3 provides a bound

$$
\begin{gather*}
\sup _{B \in \mathcal{C}_{m}}\left|\mathrm{P}\left(\zeta_{1} \in B\right)-\mathrm{P}(Z \in B)\right| \leq a(m) \gamma \\
+\gamma^{-2}\{m \sqrt{2 \pi}(2+12 \sqrt{2}(m+1) \varepsilon) \bar{R}+48 \sqrt{\pi} c(m) \varepsilon \\
\left.+\sqrt{2 \pi}\left(\varepsilon^{-1}+12 m+2+6 \sqrt{2} m(m+1) \varepsilon\right) \overline{\mathcal{L}}_{\varepsilon}\right\} \tag{46}
\end{gather*}
$$

where $\overline{\mathcal{L}}_{\varepsilon}$ appears in (40) and $\bar{R}$ is defined for a collection of random vectors $\bar{X}_{1}, \ldots, \bar{X}_{N}$ in the same manner as $R$ in (9). Namely

$$
\bar{R}=\left\|V_{1}^{-1}\right\|_{1}^{2} N \bar{\theta}_{1}
$$

where $\bar{\theta}_{1}$ is given for $\bar{X}_{U}=\left\{\bar{X}_{k}, k=1, \ldots, N\right\}$ accordingly to (7). Hence

$$
\begin{equation*}
\bar{R} \leq N b(\operatorname{Lip}(\varphi))^{2} u(q)\left\|V_{1}^{-1}\right\|_{1}^{2} . \tag{47}
\end{equation*}
$$

Set $\xi_{k}=\left(A V_{1}-I\right) V_{1}^{-1} \bar{X}_{k}$. Note that $\sum_{k=1}^{N} \operatorname{Var}\left(\xi_{k}\right)=\left(A V_{1}-I\right)\left(A V_{1}-I\right)^{*}$ and consequently

$$
\begin{gather*}
\mathrm{E}\left\|\sum_{k=1}^{N} \xi_{k}\right\|^{2} \leq \sum_{k=1}^{N} \mathrm{E}\left\|\xi_{k}\right\|^{2}+\sum_{k=1}^{N} \sum_{r=1}^{m}\left|\operatorname{cov}\left(\xi_{k r}, \sum_{l \neq k} \xi_{l r}\right)\right| \\
\leq \operatorname{Tr}\left(A V_{1}-I\right)\left(A V_{1}-I\right)^{*}+m N\left\|A V_{1}-I\right\|_{1}^{2}\left\|V_{1}^{-1}\right\|_{1}^{2} \bar{\theta}_{1} \\
\leq m\left\|A V_{1}-I\right\|_{1}^{2}(m+\bar{R}), \tag{48}
\end{gather*}
$$

here $*$ and $\operatorname{Tr}$ stand respectively for conjugation and trace of a matrix. Due to Lemma 8 and (48)

$$
\begin{equation*}
\mathrm{E}\left\|\zeta_{0}\right\|^{2} \leq \bar{\Delta} \tag{49}
\end{equation*}
$$

where $\bar{\Delta}$ is defined in (42). Estimates (45) - (49) imply (41).
To prove (44) we use instead of (46) the relation (27). The proof is complete.

## 3. Application to the kernel estimate of a density

We consider a stationary random field $X=\left\{X_{j}, j \in \mathbb{Z}^{d}\right\}$ with values in $\mathbb{R}^{s}$. Assume that $X_{0}$ has a density $f=f(x), x \in \mathbb{R}^{s}$.

Recall that, given a probability density function $K$, the Parzen - Rosenblatt (kernel) estimators for a density $f$ are defined as follows

$$
\widehat{f}_{n}(x)=\frac{1}{\left|U_{n}\right| h_{n}^{s}} \sum_{j \in U_{n}} K\left(\frac{x-X_{j}}{h_{n}}\right), x \in \mathbb{R}^{s},
$$

where $U_{n}, n \in \mathbb{N}$ are finite subsets of $\mathbb{Z}^{d},\left\{h_{n}\right\}_{n \geq 1}$ is a sequence of positive numbers such that

$$
h_{n} \rightarrow 0 \text { and }\left|U_{n}\right| h_{n}^{s} \rightarrow 0 \text { as } n \rightarrow \infty .
$$

Introduce three types of conditions concerning the properties of a field $X$, a kernel $K$ and sets $U_{n}(n \in \mathbb{N})$. Further on $C$ denotes some positive factors (not necessary the same in different expressions) which do not depend on $n$.
$\left(\mathbf{A}_{\mathbf{1}}\right)$ Let $X$ be a $(B L, \theta)$-dependent random field such that (4) holds with $a(|I|,|J|)=\min \{|I|,|J|\}$ and for some $\lambda>d(s+2) / s$

$$
\begin{equation*}
u(r)=O\left(r^{-\lambda}\right) \text { as } r \rightarrow \infty . \tag{50}
\end{equation*}
$$

Assume that the density $f$ is a Lipschitz function and for all $j \in \mathbb{Z}^{d}$ there exists a joint density $f_{j}(x, y)$ for random vectors $X_{0}$ and $X_{j}$ such that

$$
\begin{equation*}
\sup _{x, y \in \mathbb{R}^{s}} f_{j}(x, y) \leq c_{0} \tag{51}
\end{equation*}
$$

where $c_{0}$ does not depend on $j$.
$\left(\mathbf{A}_{\mathbf{2}}\right)$ Let kernel $K$ be a Lipschitz function and

$$
\begin{equation*}
\int_{\mathbb{R}^{s}}\|x\| K(x) d x<\infty \tag{52}
\end{equation*}
$$

$\left(\mathbf{A}_{\mathbf{3}}\right)$ Let $U_{n}, n \geq 1$, be regularly growing parallelepipeds in $\mathbb{Z}^{d}$, that is

$$
\begin{equation*}
U_{n}=\left\{\left(a_{1}(n), a_{1}(n)+l_{1}(n)\right] \times \ldots \times\left(a_{d}(n), a_{d}(n)+l_{d}(n)\right]\right\} \cap \mathbb{Z}^{d} \tag{53}
\end{equation*}
$$

and for some $C_{1}>0$ and all $n \in \mathbb{N}, i, k=1, \ldots, d$ one has

$$
\begin{equation*}
a_{i}(n) \in \mathbb{Z}, l_{i}(n) \in \mathbb{N}, i=1, \ldots, d \text { and } l_{i}(n) / l_{k}(n) \leq C_{1} . \tag{54}
\end{equation*}
$$

For fixed $m \in \mathbb{N}$ and different points $x_{1}, \ldots, x_{m} \in \mathbb{R}^{s}$ such that $f\left(x_{i}\right)>0, i=$ $1, \ldots, m$, define

$$
\sigma_{i}^{2}=f\left(x_{i}\right) \int_{\mathbb{R}^{s}} K^{2}(x) d x
$$

Consider a random vector $L(n)=\left(L_{1}(n), \ldots, L_{m}(n)\right)$ with components

$$
L_{i}=\sigma_{i}^{-1} \sqrt{\left|U_{n}\right| h_{n}^{s}}\left(\widehat{f}_{n}\left(x_{i}\right)-\mathrm{E} \widehat{f}_{n}\left(x_{i}\right)\right), \quad i=1, \ldots, m
$$

Theorem 5. Let $X$ be a random field and $U_{n}$ be a sequence of subsets of $\mathbb{Z}^{d}$ satisfying all the conditions mentioned above. Then there exist $\beta>0$ such that for $h_{n}=\left|U_{n}\right|^{-\beta}, n \in \mathbb{N}$, and some $C_{0}, \mu>0$ (independent of $U_{n}$ ) the following inequality holds

$$
\begin{equation*}
\sup _{B \in \mathcal{C}_{m}}|\mathrm{P}(L(n) \in B)-\mathrm{P}(Z \in B)| \leq C_{0}\left|U_{n}\right|^{-\mu}, \quad n \in \mathbb{N} . \tag{55}
\end{equation*}
$$

If $s>2 d+1$ and $\lambda \geq \lambda_{0}=d(s+4) /(s-2 d-1)$, one can choose

$$
\beta=\frac{1}{2 \Lambda(s, d, \lambda)}, \quad \mu=\frac{1}{3 \Lambda(s, d, \lambda)}
$$

where $\Lambda(s, d, \lambda)=2+s / 2+2 d(2 \lambda+s+4) / \lambda$. If $s \leq 2 d+1$ or $\lambda<\lambda_{0}$ one can take

$$
\beta=\frac{1}{5 s}, \mu=\frac{2(\lambda s-(s+2) d)}{15 s(2 \lambda d+2 d+\lambda)} .
$$

Proof. Let $\left\{p_{n}\right\}_{n \geq 1}$ and $\left\{q_{n}\right\}_{n \geq 1}$ be nonrandom sequences with values in $\mathbb{N}$ such that

$$
q_{n} \rightarrow \infty, \quad q_{n} / p_{n} \rightarrow 0, \quad p_{n}^{d} h_{n}^{s} \rightarrow 0 \text { as } n \rightarrow \infty .
$$

Following the Bernstein method each edge ( $\left.a_{k}(n), a_{k}(n)+l_{k}(n)\right]$ can be represented as a union of disjoint "large" and "small" intervals (open from the left and closed from the right) having the respective lengths $p_{n}, q_{n}, p_{n}, \ldots, q_{n}, \widetilde{p_{n k}}$ where $p_{n} \leq \widetilde{p_{n k}} \leq 3 p_{n}(n \in \mathbb{N}, k=1, \ldots, d)$; see the details in Bulinski and Millionshchikov (2002). Through every end point of these intervals (for each $\mathrm{k}=1, \ldots, \mathrm{~d})$ draw the hyperplane which is perpendicular to $i$-th axis. Thus we have a partition of $\left|U_{n}\right|$ into blocks (that is parallelepipeds) of $2^{d}$ types. Namely, to each block there corresponds a vector $\alpha=\left(\alpha_{1}, \ldots, \alpha_{d}\right)$ determining the type such that $\alpha_{i}=0$ if the length of block's edge along the $i$-th axis is equal to $q_{n}$ and $\alpha_{i}=1$ otherwise. Denote by $\Gamma$ the set of all such vectors $\alpha$. Numerate as $1,2, \ldots, m_{k}$ the intervals of lengths $p_{n}, q_{n}, p_{n}, \ldots, q_{n}, \widetilde{p_{n k}}$ along the $i$-th edge of $U_{n}$. Then we get a partition of $U_{n}$ into blocks $\Pi_{j}$ where $j \in \Lambda_{n}=\left\{\left(j_{1}, \ldots, j_{d}\right): 1 \leq j_{k} \leq m_{k}, k=1, \ldots, d\right\}$, with $m_{k}=m_{k}\left(U_{n}\right)$. Let $j \in \Lambda_{n}^{\alpha}$ if the block $\Pi_{j}$ has type $\alpha=\left(\alpha_{1}, \ldots, \alpha_{d}\right)$. Below we write $\Pi_{j}^{\alpha}$ instead of $\Pi_{j}$ to indicate the type of a parallelepiped.

Keeping the notation of Theorem 4, let $\varphi(y)=\left(\varphi_{1}(y), \ldots, \varphi_{m}(y)\right)$, where

$$
\varphi_{i}(y)=\sigma_{i}^{-1} h_{n}^{-s / 2}\left(K\left(\frac{x_{i}-y}{h_{n}}\right)-\mathrm{E} K\left(\frac{x_{i}-X_{0}}{h_{n}}\right)\right), i=1, \ldots, m, \quad y \in \mathbb{R} .
$$

We have

$$
\begin{equation*}
\operatorname{Lip}\left(\varphi_{i}\right) \leq \sigma_{i}^{-1} \operatorname{Lip}(K) h_{n}^{-s / 2-1} . \tag{56}
\end{equation*}
$$

Let $\widetilde{X}_{j}=\varphi\left(X_{j}\right), j \in \mathbb{Z}^{d}$. Clearly, the field $\left\{\widetilde{X}_{j}, j \in U\right\}$ satisfies the condition (38). Enumerate the blocks $\Pi_{j}^{1}$ of type $\mathbf{1}$, where $\mathbf{1}=(1, \ldots, 1)$, as $U^{(1)}, \ldots, U^{(N)}, N=\left|\Lambda_{n}^{1}\right|$. Put also $U^{(0)}=U \backslash\left(\cup_{k=1}^{N} U^{(k)}\right)$. It is easy to see that condition (37) holds with $b=\left(3 p_{n}\right)^{d}, q=q_{n}$. For $\alpha \in \Gamma$ we introduce auxiliary random fields with values in $\mathbb{R}^{m}$

$$
\bar{X}^{\alpha}(n)=\left\{\sum_{i \in \Pi_{j}} \widetilde{X}_{i}(n), j \in \Lambda_{n}^{\alpha}\right\}, \quad n \in \mathbb{N} .
$$

So $\bar{X}^{\mathbf{1}}(n)=\bar{X}(n)$.
Lemma 9. For each $\alpha \in \Gamma$

$$
\theta_{1}\left(\bar{X}^{\alpha}(n)\right) \leq 3^{d} p_{n}^{\|\alpha\|_{1}} q_{n}^{d-\|\alpha\|_{1}} h_{n}^{-s-2} u\left(q_{n}\right)\left(\sum_{i=1}^{s} \sigma_{i}^{-1}\right)^{2}(\operatorname{Lip}(K))^{2} .
$$

Proof. Evidently for every $\Pi_{j}^{\alpha}$ one has

$$
\rho_{0}\left(\Pi_{j}^{\alpha}, \bigcup_{i \neq j} \Pi_{i}^{\alpha}\right) \geq q_{n}
$$

where $\rho_{0}(B, D)=\min _{s \in B, t \in D}\left\{\max _{k=1}^{d}\left|s_{k}-t_{k}\right|\right\}, B, D \subset \mathbb{Z}^{d}$, and $\left|\Pi_{j}^{\alpha}\right| \leq$ $3^{d} p_{n}^{\|\alpha\|_{1}} q_{n}^{d-\|\alpha\|_{1}}$. The Lemma now follows from (35) and (56).

The next Lemma describes some properties of the field $\left\{\widetilde{X}_{j}, j \in U\right\}$ determined by the distributions of order not greater than two.

Lemma 10. For all $i, j \in U_{n}, i \neq j$ and $r, t=1, \ldots, m$

$$
\begin{gather*}
\left|\operatorname{Var}\left(\widetilde{X}_{i r}(n)\right)-1\right| \leq C h_{n}  \tag{57}\\
\left|\operatorname{cov}\left(\widetilde{X}_{i r}(n), \widetilde{X}_{j t}(n)\right)\right| \leq C h_{n}^{s} . \tag{58}
\end{gather*}
$$

If $r \neq t$ one has

$$
\begin{equation*}
\left|\operatorname{cov}\left(\widetilde{X}_{i r}(n), \widetilde{X}_{i t}(n)\right)\right| \leq C h_{n} \tag{59}
\end{equation*}
$$

Also for all $j \in U_{n}, r=1, \ldots, m$

$$
\begin{equation*}
\mathrm{E}\left|\widetilde{X}_{j r}(n)\right|^{3} \leq C h_{n}^{-s / 2} \tag{60}
\end{equation*}
$$

Proof. Note that $f$ is a bounded Lipschitz function. Due to (51) and $\left(\mathbf{A}_{\mathbf{2}}\right)$ simple estimates of the corresponding integrals lead to (57)-(60).

Now we turn to the properties of matrices $V_{1}^{2}, V_{0}^{2}$ defined in (39), denoting by $I$ the unit matrix of order $m$.

Lemma 11. The following estimate holds

$$
\left\|V_{1}^{2}-\sum_{k=1}^{N}\left|U_{n}^{(k)}\right| I\right\|_{1} \leq C \sum_{k=1}^{N}\left|U_{n}^{(k)}\right|\left(h_{n}+p_{n}^{d} h_{n}^{s}\right)
$$

Proof. By definition of $V_{1}^{2}$,(57) and (58) for any $r=1, \ldots, m$ we have

$$
\begin{aligned}
& \left|\left(V_{1}^{2}-\left(\sum_{k=1}^{N}\left|U_{n}^{(k)}\right|\right) I\right)_{r r}\right|=\mid \sum_{k=1}^{N} \sum_{j \in U^{(k)}} \operatorname{Var}\left(\widetilde{X}_{j r}(n)\right) \\
& \quad+\sum_{k=1}^{N} \sum_{\substack{j, q \in U^{(k)} \\
j \neq q}} \operatorname{cov}\left(\widetilde{X}_{j r}(n), \widetilde{X}_{q r}(n)\right)-\sum_{k=1}^{N}\left|U_{n}^{(k)}\right| \mid
\end{aligned}
$$

$$
\begin{gathered}
=\left|\sum_{k=1}^{N} \sum_{j \in U^{(k)}}\left(\operatorname{Var}\left(\widetilde{X}_{j r}(n)\right)-1\right)+\sum_{k=1}^{N} \sum_{\substack{j, q \in U^{(k)} \\
j \neq q}} \operatorname{cov}\left(\widetilde{X}_{j r}(n), \widetilde{X}_{q r}(n)\right)\right| \\
\leq C\left(h_{n}+p_{n}^{d} h_{n}^{s}\right) \sum_{k=1}^{N}\left|U_{n}^{(k)}\right|
\end{gathered}
$$

If $t=1, \ldots, m, t \neq r$ then, making use of (58) and (59),

$$
\begin{gathered}
\left|\left(V_{1}^{2}-\left(\sum_{k=1}^{N}\left|U_{n}^{(k)}\right|\right) I\right)_{r t}\right|=\left|\sum_{k=1}^{N} \sum_{j, q \in U^{(k)}} \operatorname{cov}\left(\widetilde{X}_{j r}(n), \widetilde{X}_{q t}(n)\right)\right| \\
\leq C\left(h_{n}+p_{n}^{d} h_{n}^{s}\right) \sum_{k=1}^{N}\left|U_{n}^{(k)}\right| .
\end{gathered}
$$

Lemma 11 now follows from standard norm estimates.
Thus we can write

$$
\begin{equation*}
V_{1}^{2}=\left(I+M_{n}\right) \sum_{k=1}^{N}\left|U_{n}^{(k)}\right| \tag{61}
\end{equation*}
$$

where $M_{n}$ is a matrix of order $m$, and $\left\|M_{n}\right\|_{1} \leq C\left(h_{n}+p_{n}^{d} h_{n}^{s}\right)$. Clearly, if $\left\|M_{n}\right\|_{1} \leq 1 / 2$ then $\operatorname{det} V_{1}^{2}>0$. To apply Theorem 4 we introduce $A=$ $\left|U_{n}\right|^{-1 / 2} I$. According to (60) we can use the second proposition of Theorem 4 with $\delta=1$. The next Lemma gives the estimates of matrix expressions occurring in (42).

Lemma 12. Assume that for $M_{n}$ in (61)

$$
\begin{equation*}
\left\|M_{n}\right\|_{1} \leq 1 / 2 \text { and } q_{n} / p_{n} \leq 1 / 2 \tag{62}
\end{equation*}
$$

Then

$$
\begin{gather*}
\left\|V_{1}^{-1}\right\|_{1} \leq C\left|U_{n}\right|^{-1 / 2}  \tag{63}\\
\left\|A V_{1}-I\right\|_{1}^{2} \leq C\left(q_{n} / p_{n}+h_{n}+p_{n}^{d} h_{n}^{s}\right)^{2}  \tag{64}\\
\left\|V_{0}^{2}\right\|_{1} \leq C\left|U_{n}\right|\left(q_{n} / p_{n}+p_{n}^{d-1} q_{n} h_{n}^{s}+\left(q_{n} / p_{n}\right) h_{n}^{-s-2} u\left(q_{n}\right)\right) . \tag{65}
\end{gather*}
$$

Proof. To establish (63), (64) we need the simple fact that for every real square matrix $T$ with $\|T\|_{1} \leq 1 / 2$

$$
\begin{align*}
& \left\|(I+T)^{-1 / 2}-I\right\|_{1} \leq 2^{1 / 2}\|T\|_{1}, \\
& \left\|(I+T)^{1 / 2}-I\right\|_{1} \leq 2^{-1 / 2}\|T\|_{1} . \tag{66}
\end{align*}
$$

We shall also use (61) and the estimate $\left|U_{n}^{(0)}\right| \leq N p_{n}^{d-1} q_{n} \leq\left|U_{n}\right| q_{n} / p_{n}$. Clearly

$$
\begin{gather*}
V_{1}^{-1}=\left(\left|U_{n}\right|-\left|U_{n}^{(0)}\right|\right)^{-1 / 2}\left(I+M_{n}\right)^{-1 / 2}=\left|U_{n}\right|^{-1 / 2}\left(1-\frac{\left|U_{n}^{(0)}\right|}{\left|U_{n}\right|}\right)^{-1 / 2}\left(I+M_{n}\right)^{-1 / 2} \\
A V_{1}-I=\left(1-\frac{\left|U_{n}^{(0)}\right|}{\left|U_{n}\right|}\right)^{1 / 2}\left(I+M_{n}\right)^{1 / 2}-I \tag{67}
\end{gather*}
$$

Thus in view of (66), we derive (63) and (64) from (67).
Consider now $V_{0}^{2}$. At first, we have

$$
\left\|V_{0}^{2}\right\|_{1}=\left\|\operatorname{Var}\left(\sum_{\alpha \neq 1} \sum_{j \in \Lambda_{n}^{\alpha}} \bar{X}^{(\alpha)}(n)\right)\right\|_{1} \leq\left(2^{d}-1\right) m^{3 / 2} \sum_{\alpha \neq 1}\left\|\operatorname{Var}\left(\sum_{j \in \Lambda_{n}^{\alpha}} \bar{X}^{(\alpha)}(n)\right)\right\|_{1}
$$

For fixed $\alpha \in \Gamma, \alpha \neq \mathbf{1}$,

$$
\begin{aligned}
& \left\|\operatorname{Var}\left(\sum_{j \in \Lambda_{n}^{\alpha}} \bar{X}_{j}^{(\alpha)}(n)\right)\right\|_{1} \leq\left\|\sum_{j \in \Lambda_{n}^{\alpha}} \operatorname{Var}\left(\bar{X}_{j}^{(\alpha)}(n)\right)\right\|_{1} \\
& \quad+\sum_{\substack{j, q \in \Lambda_{n}^{\alpha} \\
j \neq q}} \sum_{r, t=1}^{m}\left|\operatorname{cov}\left(\bar{X}_{j r}^{(\alpha)}(n), \bar{X}_{q t}^{(\alpha)}(n)\right)\right| .
\end{aligned}
$$

Since $\left|\Pi_{j}^{\alpha}\right| \leq p_{n}^{d-1} q_{n}$, Lemma 10 provides (analogously to the proof of Lemma 11) the bound

$$
\left\|\operatorname{Var}\left(\bar{X}_{j}^{(\alpha)}(n)\right)\right\|_{1} \leq p_{n}^{d-1} q_{n}\left(1+C\left(h_{n}+p_{n}^{d-1} q_{n} h_{n}^{s}\right)\right)
$$

and $p_{n}^{d-1} q_{n}\left|\Lambda_{n}^{\alpha}\right| \leq\left|U_{n}\right| q_{n} / p_{n}$ follows from the fact that $\left|\Lambda_{n}^{\alpha}\right| \leq N$. Consequently,

$$
\left\|\sum_{j \in \Lambda_{n}^{\alpha}} \operatorname{Var}\left(\bar{X}_{j}^{(\alpha)}(n)\right)\right\|_{1} \leq\left|U_{n}\right|\left(q_{n} / p_{n}\right)\left(1+C\left(h_{n}+p_{n}^{d-1} q_{n} h_{n}^{s}\right)\right)
$$

Applying Lemma 9 we conclude that for $j \in \Lambda_{n}^{\alpha}, r, t=1, \ldots, m$

$$
\sum_{q \in \Lambda_{n}^{\alpha}, j \neq q}\left|\operatorname{cov}\left(\bar{X}_{j r}^{(\alpha)}(n), \bar{X}_{q t}^{(\alpha)}(n)\right)\right| \leq \theta_{1}\left(\bar{X}^{\alpha}(n), \Lambda_{n}^{\alpha}\right)
$$

$$
\leq C p_{n}^{d-1} q_{n} h_{n}^{-s-2} u\left(q_{n}\right) \leq C p_{n}^{d}\left(q_{n} / p_{n}\right) h_{n}^{-s-2} u\left(q_{n}\right)
$$

The assertion of Lemma 12 is now evident.
Let us return now to the proof of Theorem 5. From (44), using (50), (60), (63) - (65), we find that, if (62) holds, then for every $\gamma \in\left(0, \gamma_{0}(m)\right]$

$$
\begin{align*}
& \sup _{B \in \mathcal{C}_{m}}|\mathrm{P}(L(n) \in B)-\mathrm{P}(Z \in B)| \leq C\left(\gamma+\gamma^{-2}\left(h_{n}^{-s-2} q_{n}^{-\lambda}\right.\right.  \tag{68}\\
& \left.\left.+\left|U_{n}\right|^{-1 / 2} h_{n}^{-s / 2} p_{n}^{2 d}+q_{n} / p_{n}+h_{n}^{2}+p_{n}^{2 d} h_{n}^{2 s}\right)\right) .
\end{align*}
$$

Let $\tau, \beta, \eta, \mu>0$ be such positive numbers that $\lambda \tau-(s+2) \beta-\mu>0, \beta-\mu>$ $0,1 / 2-s \beta / 2-2 d \eta-2 \mu>0, \eta-\tau-\mu>0, s \beta-d \eta-\mu>0$. For $n \in \mathbb{N}$ we take $\gamma=\gamma_{n}=\gamma_{0}(m)\left|U_{n}\right|^{-\mu}, h_{n}=\left|U_{n}\right|^{-\beta}, p_{n}=\left[\left|U_{n}\right|^{\eta}\right], q_{n}=\left[\left|U_{n}\right|^{\tau}\right] / 2$, [ $\left.\cdot\right]$ meaning the integer part of a number. Optimization of $(\tau, \beta, \eta, \mu)$ in (68) leads to (55) for all $n$ such that $\left|U_{n}\right|$ is large enough to guarantee the condition (62). For other $n$ the estimate (55) is obvious. The proof of Theorem 5 is complete.

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