# M-ESTIMATION, CONVEXITY AND QUANTILES ${ }^{1}$ 

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#### Abstract

The paper develops a class of extensions of the univariate quantile function to the multivariate case ( $M$-quantiles), related in a certain way to $M$-parameters of a probability distribution and their $M$-estimators. The spatial (geometric) quantiles, recently introduced by Koltchinskii and Dudley and by Chaudhuri as well as the regression quantiles of Koenker and Basset, are the examples of the $M$-quantile function discussed in the paper. We study the main properties of $M$-quantiles and develop the asymptotic theory of empirical $M$-quantiles. We use $M$-quantiles to extend $L$-parameters and $L$-estimators to the multivariate case; to introduce a bootstrap test for spherical symmetry of a multivariate distribution, and to extend the notion of regression quantiles to multiresponse linear regression models.


1. Introduction. The univariate quantile function has been extensively used in statistical inference. One of the reasons of such a wide applicability is that many important robust statistics (like, for instance, $L$ - and $R$-estimators) are defined as functionals of the empirical quantile function. Our goal is to extend the notion of quantiles to a broader context, specifically, to probability measures on finite-dimensional spaces. Since there is no really natural ordering of vectors in $\mathbb{R}^{d}$ for $d>1$, such an extension of quantiles is a hard problem in the statistics of multivariate data. A brief outline of some of the known approaches for the study of order statistics and quantiles in the multivariate case is as follows:
2. Tukey (1975), Barnett (1976), Eddy (1985) and Reiss (1989) suggested different ways of ordering multivariate data (by peeling the convex hulls of data points, using auxillary real-valued functions, etc.).
3. Pyke $(1975,1985)$ considered extensions of the empirical quantile process. Einmahl and Mason (1992) developed a class of extensions of the univariate quantile function to the multivariate case. Their quantile function $U_{6}$ is defined as

$$
U_{6}(t):=\inf \{\lambda(C): P(C) \geq t, C \in \mathscr{C}\} \quad \text { for } t \in(0,1),
$$

where $P$ is a probability measure on $\mathbb{R}^{d}, \mathscr{C} \subset \mathbb{R}^{d}$ is a class of Borel subsets and $\lambda$ is a real-valued function defined on $\ell$. They studied asymptotics of the corresponding empirical quantile processes.

[^0]3. Koltchinskii and Dudley (1996) ([K-D]) [see also Koltchinskii (1994a, b, 1996)], and Chaudhuri (1996) ([Ch]) suggested an extension of the distribution function $F_{P}$ (of a probability measure $P$ ) to the multivariate case such that the inverse function $F_{P}^{-1}$ can be viewed as an extension of the quantile function. These functions possess many properties known and used in the one-dimensional case. For instance, $F_{P}$ characterizes $P, F_{P}^{-1}$ is equivariant with respect to translations (as well as with respect to orthogonal transformations), and the asymptotic properties of their empirical versions are the same as in the univariate case [see, e.g., Shorack and Wellner (1986)]. The corresponding "median" $F_{P}^{-1}(0)$ is nothing but the well-known Haldane's spatial $L_{1}$-median [see Haldane (1948)].

Our approach in this paper is rather close to Koltchinskii (1996). It allows us, for example, to define $L$-parameters of the distribution $P$ as integral functionals of the form $L(P):=\int h\left(F_{P}^{-1}(t)\right) d m(t)$, and to define $L$-estimators as $L$-parameters of the empirical measure $P_{n}$, avoiding ordering of multivariate data [see Section 4 of this paper and Koltchinskii (1995a)]. A multivariate extension of $R$-estimators can be also defined in a similar way, as a certain functional of the empirical quantile function.

In the one-dimensional case, it is well known that the set of all medians of $P$ coincides with the set of all minimal points of the function $\int(|s-x|-|x|) P(d x)$. However, it is less known that quantiles of $P$ can be also characterized as minimal points of a certain function. For instance, let

$$
\begin{equation*}
f_{P}(s):=0.5 \int_{\mathbb{R}^{1}}(|s-x|-|x|+s) P(d x) \quad \text { for } s \in \mathbb{R}^{1} . \tag{1.1}
\end{equation*}
$$

Then for any $t \in(0,1)$, the set of all $t$ th quantiles of $P$ is exactly the set of all minimal points of the function $f_{P, t}(s):=f_{P}(s)-s t, s \in \mathbb{R}^{1}$. Note that the function $f_{P}$ is convex, the distribution function $F_{P}$ is a subgradient of $f_{P}$ and the quantile function $F_{P}^{-1}$ is a subgradient of the Young-Fenchel conjugate of $f_{P}$, defined as

$$
f_{P}^{*}(t):=\sup _{s \in \mathbb{R}^{1}}\left[s t-f_{P}(s)\right] \quad \text { for } t \in \mathbb{R}^{1}
$$

and known in smooth analysis as the Legendre transformation of $f_{P}$. In what follows we extend this construction to the multivariate case and define $M$-distribution and $M$-quantile functions of a probability measure $P$ on an arbitrary measurable space. These functions are related to the class of $M$-parameters of $P$, defined by convex minimization. The relationship is essentially the same as between the distribution and quantile functions and the median, so that the $M$-parameter can be viewed as a "median," whereas the corresponding $M$-estimator is a "sample median."

We use some notations, as introduced in the Appendix. Let $(\mathbf{X}, \mathscr{A}, P)$ be a probability space, and let $f: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{1}$. Suppose that for any $s \in \mathbb{R}^{d}, f(s, \cdot)$ is an integrable function. We define the integral transformation $(P, s) \mapsto f_{P}(s) \in$
$\mathbb{R}^{1}$ of the probability measure $P$ with the kernel $f$ :

$$
f_{P}(s):=\int_{\mathbf{X}} f(s, x) P(d x) \quad \text { for } s \in \mathbb{R}^{d} .
$$

If there exists a point $s_{0}$, which minimizes $f_{P}$ on $\mathbb{R}^{d}$, it is called an $M$ parameter of $P$ with respect to the kernel $f$. Under some smoothness, the $M$-parameters are the solutions of the equation $F_{P}(s)=0$, where $F_{P}(s):=$ $\int_{\mathbf{X}} \nabla f(s ; x) P(d x)$.

Given a sample ( $X_{1}, \ldots, X_{n}$ ) from the distribution $P, P_{n}$ denotes the empirical measure based on this sample; that is, $P_{n}:=n^{-1} \sum_{1}^{n} \delta_{X_{k}}$, where $\delta_{x}$ is the unit point mass at $x \in \mathbf{X}$. An $M$-parameter of $P_{n}$ is known as an $M$ estimator.

Definition 1.1. We call a minimal point of the functional

$$
f_{P, t}(s):=f_{P}(s)-\langle s, t\rangle, \quad s \in \mathbb{R}^{d}
$$

an ( $M, t$ )-parameter of $P$ with respect to $f$. An ( $M, t$ )-parameter of the empirical measure $P_{n}$ will be called an ( $M, t$ )-estimator (with respect to $f$ ) based on a sample ( $X_{1}, \ldots, X_{n}$ ).

Of course, $M$-parameters of $P$ coincide with ( $M, 0$ )-parameters, and $M$ estimators are just ( $M, 0$ )-estimators.

Huber (1967) developed the asymptotic theory of $M$-estimators. Pollard (1988), Haberman (1989), Niemiro (1992) and others have emphasized the role of convexity of the kernel $f(s, x)$ in $s$ in the asymptotic study of $M$-estimators (as well as some other estimators defined by minimization of statistical functionals). They showed that in this case the theory can be substantially simplified and the asymptotic normality of $M$-estimators can be established under minimal and very natural assumptions.

We use convex analysis rather extensively in this paper. The relevant facts can be found in Rockafellar (1970) ([Ro]). Some frequently used notions are given in the Appendix.

In what follows, we assume that $f(\cdot, x)$ is a convex function for $P$ almost all $x$. We call such a kernel $f P$-convex. If there exists $A \in \mathscr{A}$ such that $f(\cdot, x)$ is strictly convex for all $x \in A$ and $P(A)>0$, then $f$ will be called a $P$-strictly convex kernel.

Given a $P$-convex kernel $f$, we call the subdifferential map $\partial f_{P}$ of the convex function $f_{P}$ the $M$-distribution map of $P$ with respect to $f$. Its inverse, $\partial f_{P}^{-1}=\partial f_{P}^{*}$, is called the $M$-quantile map of $P$ with respect to $f$. For the convex function $f_{P}$ given by (1.1), $\partial f_{P}$ and $\partial f_{P}^{-1}$ are multivalued versions of the usual distribution and quantile functions of $P$ (see Example 2.2), which motivates our definitions. For any $t \in \mathbb{R}^{d},\left(\partial f_{P}\right)^{-1}(t)$ is the set of all $(M, t)-$ parameters of $P$. If, in addition, $f$ is a $P$-strictly convex kernel, then $\partial f_{P}^{-1}$ is a single-valued map, so that the ( $M, t$ )-parameter is unique (for all $t$ such that it exists). In this case we call a subgradient $F_{P}$ (i.e., a map $F_{P}$ from $\mathbb{R}^{d}$ into $\mathbb{R}^{d}$ such that $F_{P}(s) \in \partial f_{P}(s)$ for all $\left.s \in \mathbb{R}^{d}\right)$ the $M$-distribution function
of $P$ with respect to $f$, and we call its inverse $F_{P}^{-1}$ from $F_{P}\left(\mathbb{R}^{d}\right)$ onto $\mathbb{R}^{d}$, the $M$-quantile function. In particular, if $0 \in F_{P}\left(\mathbb{R}^{d}\right)$, then the $M$-parameter of $P$ is unique and it coincides with the "median" $F_{P}^{-1}(0)$ and $M$-estimators are "sample medians" $F_{P_{n}}^{-1}(0)$. Using some basic facts of convex analysis, it is easy to describe the range $F_{P}\left(\mathbb{R}^{d}\right)$ of the map $F_{P}$ (see Section 2). We also give some conditions when $F_{P}$ charaterizes $P$ (in the sense that $F_{P}=F_{Q}$ implies $P=Q$ ). In Section 3, we study equivariance of $M$-distribution and $M$-quantile maps with respect to the groups of translations and orthogonal transformations and obtain representations of these functions for spherically symmetric probability distributions.

Section 5 develops the asymptotic theory of empirical $M$-distribution and $M$-quantile functions. Specifically, we prove that the sequences $n^{1 / 2}\left(F_{P_{n}}-F_{P}\right)$ and $n^{1 / 2}\left(F_{P_{n}}^{-1}-F_{P}^{-1}\right)$ converge in distribution to Gaussian processes.

We apply these results to two main examples. First of all, we consider a class of extensions of quantiles to the multivariate case. To define such spatial quantiles we use the kernel $f(s, x):=\varphi(|s-x|)-\varphi\left(\left|s_{0}-x\right|\right), s, x \in \mathbb{R}^{d}$, where $s_{0} \in \mathbb{R}^{d}$ is a fixed point and $\varphi$ is a convex nondecreasing function on $[0,+\infty)$, differentiable in $(0,+\infty)$. In this case $F_{P}$ is defined by

$$
F_{P}(s):=\int_{\{x \neq s\}} \Phi(|s-x|) \frac{s-x}{|s-x|} P(d x) \quad \text { for } s \in \mathbb{R}^{d},
$$

where $\Phi:=\varphi^{\prime}$. If $\varphi$ is strictly increasing and $P$ is not concentrated in a line, then the range $V_{P}:=F_{P}\left(\mathbb{R}^{d}\right)$ is a ball (which, in particular, could coincide with the whole space $\mathbb{R}^{d}$ ) with a number of spherical "holes" in it (one "hole" for each atom of $P$ ). The map $F_{P}$ is one-to-one from $\mathbb{R}^{d}$ onto $V_{P}$, so that $F_{P}^{-1}$ is well defined (see Proposition 2.6). In particular, if $\varphi(u) \equiv u$ and $\Phi \equiv 1$, we get an extension of the distribution and quantile functions to the multivariate case as introduced by [K-D] and [Ch]. We study the main properties of these $M$-distribution and $M$-quantile functions (analytical properties, characterization of the distribution $P$, asymptotics of their empirical counterparts, etc.). We show also their equivariance with respect to translations and orthogonal transformations and obtain representations of $F_{P}$ and $F_{P}^{-1}$ for spherically symmetric distributions. These representations are used in Section 4 to develop tests for spherical symmetry of multivariate distributions.

Another example is related to the notion of regression quantiles, introduced by Koenker and Basset (1978). Consider a sample ( $X_{i}, Y_{i}$ ), $i=1, \ldots, n$ of i.i.d. observations, where $X_{i}, i \geq 1$ take values in $\mathbb{R}^{m}$ and $Y_{i}, i \geq 1$ take values in $\mathbb{R}^{1}$. Suppose that our observations satisfy the following linear regression model with random design:

$$
Y_{i}=\left\langle s_{0}, X_{i}\right\rangle+\xi_{i}, \quad 1 \leq i \leq n,
$$

where $s_{0} \in \mathbb{R}^{m}$ and $\xi_{i}, i \geq 1$ are i.i.d. errors, independent of $X_{i} i \geq 1$. Koenker and Basset (1978) define (empirical) regression $\alpha$-quantiles for $\alpha \in[0,1]$ as
minimal points of the function in $s$,

$$
\sum_{1}^{n}\left[\left|Y_{i}-\left\langle s, X_{i}\right\rangle\right|-(2 \alpha-1)\left\langle s, X_{i}\right\rangle\right]
$$

One of the reasons to introduce this notion was to develop a technique of robust estimation for regression models. Koenker and Portnoy (1987) studied $L$ estimators for regression, defined in terms of regression quantiles. Chaudhuri (1991) considered nonparametric estimates of regression quantiles. See also Gutenbrunner and Jurečková (1992) and Gutenbrunner, Jurečková, Koenker and Portnoy (1993) for applications of this notion.
[Ch] suggested extending the notion to the case of multiresponse regression, using the version of spatial quantiles introduced by him and by [K-D]. In this paper, we define regression quantiles for multiresponse regression models and show that they can be viewed and studied as a special case of the $M$-quantile function with respect to a certain convex kernel.
2. $M$-distribution and $M$-quantile functions with respect to convex kernels. Let $f$ be a $P$-convex kernel such that for all $s \in \mathbb{R}^{d} f(s, \cdot)$ is a $P$ integrable function. Then $f_{P}$ is a convex function from $\mathbb{R}^{d}$ into $\mathbb{R}^{1}$. Moreover, if $f$ is a $P$-strictly convex kernel, then $f_{P}$ is a strictly convex function. Let $d f(s ; v ; x)$ denote the directional derivative of the function $f(\cdot, x)$ at a point $s \in \mathbb{R}^{d}$ in the direction $v \in \mathbb{R}^{d}$,

$$
d f(s ; v ; x):=\inf _{\lambda \in(0,+\infty)} \Delta f(s ; v ; \lambda ; x)=\lim _{\lambda \downarrow 0} \Delta f(s ; v ; \lambda ; x)
$$

where $\Delta f(s ; v ; \lambda ; x):=\lambda^{-1}(f(s+\lambda v ; x)-f(s ; x))$. The directional derivative exists for any $x$ such that $f(\cdot, x)$ is convex; thus, it exists $P$ a.e. Since the function $\Delta f(s ; v ; \lambda ; x)$ is increasing in $\lambda>0$, one can use the monotone convergence theorem to show that for all $s \in \mathbb{R}^{d}$,

$$
d f_{P}(s ; v)=\lim _{\lambda \downarrow 0} \Delta f_{P}(s ; v ; \lambda)=\lim _{\lambda \downarrow 0} \int_{\mathbf{X}} \Delta f(s ; v ; \lambda ; x) P(d x)=\int_{\mathbf{X}} d f(s ; v ; x) P(d x)
$$

Clearly,

$$
\partial f_{P}(s)=\bigcap_{v \in \mathbb{R}^{d}, v \neq 0}\left\{t: \int_{\mathbf{x}} d f(s ; v ; x) P(d x) \geq\langle v, t\rangle\right\}
$$

Given $s \in \mathbb{R}^{d}, x \mapsto \partial f(s ; x)$ is a closed convex multivalued map. Given a multivalued map $x \mapsto C(x) \subset \mathbb{R}^{d}$, an integral $\int_{\mathbf{X}} C(x) P(d x)$ is defined as a set

$$
\left\{\int_{\mathbf{X}} c(x) P(d x): c: \mathbf{X} \mapsto \mathbb{R}^{d} \text { is } P \text {-integrable, } c(x) \in C(x) \text { for } P \text {-almost all } x\right\}
$$

The next statement is a special case of Theorem 8.3.4 of Ioffe and Tihomirov (1974). It shows that one can subdifferentiate under the integral sign and gives an integral representation of the multivalued map $\partial f_{P}$.

THEOREM 2.1. Suppose $f$ is a $P$-convex kernel, and for all $s \in \mathbb{R}^{d} f(s ; \cdot)$ is $P$-integrable. Then, for all $s \in \mathbb{R}^{d}$,

$$
\partial f_{P}(s)=\int_{\mathbf{X}} \partial f(s ; x) P(d x)
$$

A map $F_{P}: \mathbb{R}^{d} \mapsto \mathbb{R}^{d}$ is a subgradient of the convex function $f_{P}$ if and only if there exists a map $F: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{d}$, such that $F(s, x) \in \partial f(s, x) P$ a.e. [so, $F(s, x)$ is a subgradient of $f(s, x) P$ a.e.], $F(s, \cdot)$ is $P$-integrable and

$$
\begin{equation*}
F_{P}(s):=\int_{\mathbf{X}} F(s ; x) P(d x) \tag{2.1}
\end{equation*}
$$

EXAMPLE 2.2. Let $\mathbf{X}=\mathbb{R}^{1}$, and let $f(s, x):=0.5(|s-x|-|x|+s)$ for $s, x \in \mathbb{R}^{1}$. Then

$$
f_{P}(s):=0.5 \int_{\mathbb{R}^{1}}(|s-x|-|x|) P(d x)+0.5 s \quad \text { for } s \in \mathbb{R}^{1}
$$

In this case for any $x \in \mathbb{R}^{1}$, a function $F(\cdot ; x)$, defined by $F(s ; x):=\mathbf{I}_{(-\infty, s]}(x)$ for $x, s \in \mathbb{R}^{1}$, is a subgradient of the convex function $f(\cdot ; x)$. Therefore, a function $F_{P}$, defined by

$$
F_{P}(s):=\int_{\mathbb{R}^{1}} \mathbf{I}_{(-\infty, s]}(x) P(d x)=P((-\infty, s]) \quad \text { for } s \in \mathbb{R}^{1}
$$

is a subgradient of $f_{P}$. Clearly, $F_{P}$ is the distribution function of the probability measure $P$. The inverse $F_{P}^{-1}$ of the function $F_{P}$ is the quantile function of $P$. It's easy to show that the subdifferential of $f_{P}$ is a multivalued map $\partial f_{P}$, defined by $\partial f_{P}(s)=\left[F_{P}(s-), F_{P}(s)\right]$ for $s \in \mathbb{R}^{1}$. On the other hand, the inverse of the multivalued map $\partial f_{P}$,

$$
\partial f_{P}^{-1}(t)= \begin{cases}{\left[F_{P}^{-1}(t), F_{P}^{-1}(t+)\right],} & \text { for } t \in(0,1) \\ \left(-\infty, F_{P}^{-1}(0+)\right], & \text { for } t=0 \\ {\left[F_{P}^{-1}(1),+\infty\right),} & \text { for } t=1 \\ \varnothing, & \text { for } t \notin[0,1]\end{cases}
$$

is the subdifferential $\partial f_{P}^{*}$ of the conjugate $f_{P}^{*}$. In particular, $F_{P}^{-1}(t)$ is a subgradient of the convex function $f_{P}^{*}(t), t \in(0,1)$ [for $t=0$ and $t=1, F_{P}^{-1}(0+)$ and $F_{P}^{-1}(1)$ are also subgradients of $f_{P}^{*}$, if they are finite]. The set of $(M, t)-$ parameters of $P$ coincides with the set $\left[F_{P}^{-1}(t), F_{P}^{-1}(t+)\right]$ of $t$-quantiles of $P$. Clearly, $\partial f_{P}$ and $\partial f_{P}^{*}=\partial f_{P}^{-1}$ can be viewed as multivalued versions of the usual distribution and quantile functions, respectively.

Motivated by this example, we introduce the following definition in the general case.

DEFINITION 2.3. The subdifferential map $\partial f_{P}$ of the function $f_{P}$ will be called the $M$-distribution map of $P$ with respect to $f$. Respectively, the subdifferential map $\partial f_{P}^{*}=\partial f_{P}^{-1}$ of the conjugate map $f_{P}^{*}$ will be called the $M-$ quantile map of $P$ with respect to $f$. Any point of the set $\partial f_{P}^{-1}(t)$ is called an ( $M, t$ )-quantile of $P$ with respect to $f$. If $\partial f_{P}$ (respectively, $\partial f_{P}^{-1}$ ) is a single-
valued map, we call it the $M$-distribution (respectively, the $M$-quantile) function of $P$ with respect to $f$.

Define

$$
r_{f}(v ; x):=\sup _{s \in \mathbb{R}^{d}} d f(s ; v ; x) \quad \text { for } v \in \mathbb{R}^{d} \text { and } x \in \mathbf{X},
$$

so that $r_{f}(\cdot, x)$ is the recession function of $f(\cdot, x)$. It follows from the properties of the recession function (see [Ro]), that for any $s \in \mathbb{R}^{d}$ and for $P$ almost all $x$

$$
r_{f}(v ; x)=\lim _{\lambda \uparrow+\infty} \Delta f(s ; v ; \lambda ; x) .
$$

Therefore, by the monotone convergence theorem, we have the following representation for the recession function of $f_{P}$ :

$$
\begin{aligned}
r_{P}(v) & :=r_{f_{P}}(v):=\sup _{s \in \mathbb{R}^{d}} d f_{P}(s ; v)=\lim _{\lambda \uparrow+\infty} \Delta f_{P}(s ; v ; \lambda) \\
& =\lim _{\lambda \uparrow+\infty} \int_{\mathbf{x}} \Delta f(s ; v ; \lambda ; x) P(d x)=\int_{\mathbf{x}} r_{f}(v ; x) P(d x) \quad \text { for } v \in \mathbb{R}^{d} .
\end{aligned}
$$

We define a convex set

$$
B_{P}:=B_{f_{P}}:=\bigcap_{v \in \mathbb{R}^{d}, v \neq 0}\left\{t:\langle v, t\rangle<r_{P}(v)\right\} .
$$

The following properties of $M$-distribution and $M$-quantile maps easily follow from the general facts of convex analysis (see [Ro], especially Sections 2327). They are quite similar to the properties of the multivalued versions of the usual distribution and quantile functions [such as their monotonicity, semicontinuity, relationship with ( $M, t$ )-parameters (quantiles) and $M$-parameters (medians), etc.; see Example 2.2].
2.1. Properties of $M$-distribution and $M$-quantile maps. Let $f$ be a $P$ convex kernel. Then we have the following properties.

1. The $M$-distribution map $\partial f_{P}$ of $P$ with respect to $f$ is a monotone upper semicontinuous multivalued map. Moreover, if $f$ is $P$-strictly convex, then $\partial f_{P}$ is strictly monotone.
2. The $M$-quantile map of $P$ with respect to $f$ is the inverse of the $M$ distribution map: $\partial f_{P}^{*}=\left(\partial f_{P}\right)^{-1}$. It is a monotone upper semicontinuous multivalued map.

Let $D_{P}:=D_{f_{P}}$ be the set of all points $s \in \mathbb{R}^{d}$ such that $d f_{P}(s ; v)$ is linear in $v$.
3. $f_{P}$ is differentiable at points $s \in D_{P}$ and $\partial f_{P}(s)=\left\{\nabla f_{P}(s)\right\}, s \in D_{P} . D_{P}$ is a dense subset of $\mathbb{R}^{d}$ and its complement has Lebesgue measure 0 . The gradient $\nabla f_{P}$ is continuous on $D_{P}$.
4. For any $t \in \mathbb{R}^{d}$, the set of all ( $M, t$ )-parameters of $P$ with respect to $f$ coincides with $\partial f_{P}^{*}(t)=\partial f_{P}^{-1}(t)$ [the set of all ( $\left.M, t\right)$-quantiles of $P$ with respect to $f]$. Thus, an $(M, t)$-parameter of $P$ exists iff $f_{P}^{*}$ is subdifferentiable at $t$. The set of ( $M, t$ )-parameters is nonempty and bounded iff $t \in B_{P}$. The set of ( $M, t$ )-parameters consists of the unique point $s$ iff $f_{P}^{*}$ is differentiable at $t$ and $s=\nabla f_{P}^{*}(t)$.

Let now $f$ be a $P$-strictly convex kernel. Denote

$$
U_{P}:=U_{f_{P}}:=B_{P} \backslash \bigcup_{s \in \mathbb{R}^{d} \backslash D_{P}} \partial f_{P}(s)
$$

Then:
5. The convex function $f_{P}^{*}$ is continuously differentiable in $B_{P}$ with the gradient (and subgradient) $\nabla f_{P}^{*}$. Moreover, $\nabla f_{P}^{*}$ is the $M$-quantile function of $P$ with respect to $f$.
6. For any $t \in B_{P}$, the unique ( $M, t$ )-parameter of $P$ coincides with $\nabla f_{P}^{*}(t)$. The gradient $\nabla f_{P}$ of the function $f_{P}$ is a homeomorphism of $D_{P}$ and $U_{P}$. For any $t \in U_{P}, \nabla f_{P}^{*}(t)=\left(\nabla f_{P}\right)^{-1}(t)$.
To be more specific, assume that $\mathbf{X}=\mathbb{R}^{d}$, and that $P$ is a Borel probability on $\mathbb{R}^{d}$. Define a function

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}\left(h(s-x)-h\left(s_{0}-x\right)\right) P(d x) \quad \text { for } s \in \mathbb{R}^{d},
$$

where $h: \mathbb{R}^{d} \mapsto \mathbb{R}^{1}$ is convex, $s_{0} \in \mathbb{R}^{d}$ is a fixed point, and $h(s-\cdot)-h\left(s_{0}-\cdot\right)$ is $P$-integrable. Using Theorem 2.1, we get the following representation of the $M$-distribution function:

$$
\begin{equation*}
\partial f_{P}(s)=\int_{\mathbb{R}^{d}} \partial h(s-x) P(d x) \quad \text { for } s \in \mathbb{R}^{d} . \tag{2.2}
\end{equation*}
$$

For all $x \in \mathbb{R}^{d}$ and $v \in \mathbb{R}^{d}$,

$$
r_{f}(v ; x)=\lim _{\lambda \uparrow+\infty} \frac{h(\lambda v-x)-h(-x)}{\lambda}=r_{h}(v) .
$$

Therefore

$$
r_{P}(v):=r_{f_{P}}(v)=\int_{\mathbb{R}^{d}} r_{f}(v ; x) P(d x)=r_{h}(v) \quad \text { for } v \in \mathbb{R}^{d},
$$

so the recession function of $f_{P}$ does not depend on $P$. It follows that $B_{P}$ coincides with a convex set

$$
B_{h}:=\bigcap_{v \in \mathbb{R}^{d}, v \neq 0}\left\{t:\langle v, t\rangle<r_{h}(v)\right\} .
$$

Assume that $h$ is continuously differentiable in $\mathbb{R}^{d} \backslash\{0\}$. It means that $D_{h}=\mathbb{R}^{d} \backslash\{0\}$, unless $h$ is differentiable at 0 too, and then $D_{h}=\mathbb{R}^{d}$. Let $t_{0} \in \partial h(0)$. Then the map $H: \mathbb{R}^{d} \mapsto \mathbb{R}^{d}$, defined by

$$
H(s):= \begin{cases}\nabla h(s), & \text { for } s \neq 0, \\ t_{0}, & \text { for } s=0\end{cases}
$$

is a subgradient of the convex function $h$. Therefore the map $F_{P}$, defined by

$$
F_{P}(s):=\int_{\mathbb{R}^{d}} H(s-x) P(d x),
$$

is a subgradient of $f_{P}$.

Let $A_{P}$ be the set of all atoms of the probability measure $P$, and let, for $s \in A_{P}, p(s):=P(\{s\})$. Since

$$
\partial h(s-x)= \begin{cases}\partial h(0), & \text { for } x=s \\ \{H(s-x)\}, & \text { for } x \neq s\end{cases}
$$

we have

$$
\partial f_{P}(s)= \begin{cases}\left\{F_{P}(s)\right\}, & \text { for } s \notin A_{P} \\ F_{P}(s)+p(s)(\partial h(0)-H(0)), & \text { for } s \in A_{P}\end{cases}
$$

Thus $f_{P}$ is continuously differentiable in $\mathbb{R}^{d} \backslash A_{P}$. If $h$ is not differentiable at 0 , then $D_{P}=\mathbb{R}^{d} \backslash A_{P}$, otherwise $f_{P}$ is continuously differentiable in $\mathbb{R}^{d}$ and $D_{P}=\mathbb{R}^{d}$. We also define

$$
U_{P}:=B_{h} \backslash \bigcup_{s \in A_{P}}\left(F_{P}(s)+p(s)(\partial h(0)-H(0))\right)
$$

The following statement immediately follows from the main properties of $M$-distribution and $M$-quantile maps.

Proposition 2.4. (i) For all $P, F_{P}$ is a monotone map, continuous in $\mathbb{R}^{d} \backslash A_{P}$. If $f(s ; x)=h(s-x)-h\left(s_{0}-x\right)$ is a $P$-strictly convex kernel, then
(ii) $F_{P}$ is a strictly monotone one-to-one map from $\mathbb{R}^{d}$ onto $U_{P} \cup F_{P}\left(A_{P}\right)$ and it is a homeomorphism of $\mathbb{R}^{d} \backslash A_{P}$ and $U_{P}$. If $P$ is nonatomic, or if $h$ is differentiable at 0 , then $F_{P}$ is a homeomorphism of $\mathbb{R}^{d}$ and $B_{h}$.

We extend the map $F_{P}^{-1}$ (defined in a natural way on the set $U_{P}$ ) to the whole set $B_{h}$ assuming that for $t \in B_{h} \backslash U_{P} F_{P}^{-1}(t)$ coincides with the unique $s \in A_{P}$ such that $t \in\left(F_{P}(s)+p(s)\left(\partial h(0)-F_{P}(0)\right)\right)$. Then the following holds.
(iii) For all $t \in B_{h}$ the unique $(M, t)$-parameter of $P$ coincides with $F_{P}^{-1}(t)$.

Note that if $h$ is a strictly convex function, then $f(s ; \cdot)=h(s-\cdot)-h\left(s_{0}-\cdot\right)$ is a $P$-strictly convex kernel. It also holds, if $P$ is not concentrated in a line, $h$ is convex and for all $s_{1}, s_{2} \in \mathbb{R}^{d}$, such that $s_{1} \neq \alpha s_{2}$ for all $\alpha \geq 0$, we have

$$
h\left(\lambda s_{1}+(1-\lambda) s_{2}\right)<\lambda h\left(s_{1}\right)+(1-\lambda) h\left(s_{2}\right) \quad \text { for } \lambda \in(0,1)
$$

It is well known that in the one-dimensional case the distribution function $F_{P}$ characterizes $P$, that is, $F_{P} \equiv F_{Q}$ implies $P=Q$. Now we extend this importatant property to the class of $M$-distribution functions defined above [see Koldobskii (1990) and Mattner (1992) for more general conditions on the kernels of convolution integrals, characterizing probability measures]. Denote $\rho^{\prime}\left(\mathbb{R}^{d}\right)$ the space of tempered generalized functions (distributions). Given $g \in$ $\rho^{\prime}\left(\mathbb{R}^{d}\right)$, denote $\tilde{g}$ the Fourier transform of $g$ [see Gelfand and Shilov (1964) or Reed and Simon (1975)]. Since Borel probabilities on $\mathbb{R}^{d}$ can be also viewed as generalized functions, we use the same notation for their Fourier transforms (characteristic functions):

$$
\tilde{P}(\lambda):=\int_{\mathbb{R}^{d}} e^{i\langle\lambda, x\rangle} P(d x) \quad \text { for } \lambda \in \mathbb{R}^{d}
$$

Given $h$, let $\mathscr{P}_{h}$ be the set of all Borel probability measures $P$ on $\mathbb{R}^{d}$, such that the following hold:

1. $h(s-x)-h(-x)$ is a $P$-strictly convex kernel;
2. $h(s-\cdot)-h(-\cdot)$ is $P$-integrable for all $s \in \mathbb{R}^{d}$;
3. for some $\alpha>0 \int_{\mathbb{R}^{d}}|H(s-x)| P(d x)=O\left(|s|^{\alpha}\right)$ as $s \rightarrow \infty$.

ThEOREM 2.5. Suppose that the Fourier transform $\tilde{h}$ is given by a Borel function on $\mathbb{R}^{d}$ such that $\tilde{h} \neq 0$ a.e. in $\mathbb{R}^{d}$. Then, for any two probability measures $P, Q \in \mathscr{P}_{h}$, the equality $F_{P} \equiv F_{Q}$ implies that $P=Q$.

Let us consider a special case of the previous kernel. Namely, assume that $h(x):=\varphi(|x|)$ for $x \in \mathbb{R}^{d}$, where $\varphi$ is a convex nondecreasing function on $[0,+\infty)$. We also suppose that the functions $\varphi(|s-\cdot|)-\varphi\left(\left|s_{0}-\cdot\right|\right), s \in \mathbb{R}^{d}$, are $P$-integrable, and define

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}\left(\varphi(|s-x|)-\varphi\left(\left|s_{0}-x\right|\right)\right) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

Denote

$$
\rho_{\varphi}:=\lim _{\varepsilon \downarrow 0} \frac{\varphi(\varepsilon)-\varphi(0)}{\varepsilon} \quad \text { and } \quad R_{\varphi}:=\lim _{\lambda \uparrow+\infty} \frac{\varphi(\lambda)}{\lambda} .
$$

The use of Theorem 2.1 and a bit of subdifferential calculus allow us to get the following representation of the $M$-distribution map $\partial f_{P}$ for all $s \in \mathbb{R}^{d}$ :

$$
\begin{aligned}
\partial f_{P}(s) & =\int_{\mathbb{R}^{d}} \partial(\varphi(|s-x|)) P(d x) \\
& =\int_{\mathbb{R}^{d} \backslash\{s\}} \partial \varphi(|s-x|) \frac{s-x}{|s-x|} P(d x)+p(s) \bar{B}\left(0, \rho_{\varphi}\right)
\end{aligned}
$$

We assume for simplicity that $\varphi$ is differentiable in $(0,+\infty)$. Let $\Phi(\lambda):=$ $\varphi^{\prime}(\lambda), \lambda>0$. Then $\partial f_{P}(s)=F_{P}(s)+p(s) \bar{B}\left(0, \rho_{\varphi}\right)=\bar{B}\left(F_{P}(s), p(s) \rho_{\varphi}\right), s \in \mathbb{R}^{d}$, where

$$
F_{P}(s):=\int_{\mathbb{R}^{d} \backslash\{s\}} \Phi(|s-x|) \frac{s-x}{|s-x|} P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

It follows that $F_{P}$ is a subgradient of $f_{P}$. For $h(x):=\varphi(|x|)$ we have

$$
r_{h}(v):=\lim _{\lambda \uparrow+\infty} \frac{h(\lambda v)-h(0)}{\lambda}=\lim _{\lambda \uparrow+\infty} \frac{\varphi(\lambda|v|)-\varphi(0)}{\lambda}=R_{\varphi}|v|
$$

and it follows that $B_{P}=B_{h}=B\left(0 ; R_{\varphi}\right)$. Let

$$
U_{P}:=B\left(0 ; R_{\varphi}\right) \backslash \bigcup_{s \in A_{P}} \bar{B}\left(F_{P}(s) ; p(s) \rho_{\varphi}\right)
$$

In this case we get the following special version of Proposition 2.4.

Proposition 2.6. (i) For all $P, F_{P}$ is a monotone continuous in $\mathbb{R}^{d} \backslash A_{P}$. For any $s \in A_{P}$,
(ii) $\quad F_{P}(\lambda v) \rightarrow\left\{\begin{array}{ll}R_{\varphi} v, & \text { for } R_{\varphi}<+\infty \\ \infty, & \text { otherwise }\end{array} \quad\right.$ as $\lambda \rightarrow \infty$ for $v \in S^{d-1}$.

If either $\varphi$ is strictly convex, or it is strictly increasing and P is not concentrated in a line, then:
(iii) $F_{P}$ is a strictly monotone one-to-one map from $\mathbb{R}^{d}$ onto $U_{P} \cup F_{P}\left(A_{P}\right)$ and it is a homeomorphism of $\mathbb{R}^{d} \backslash A_{P}$ and $U_{P}$. If $P$ is nonatomic, then $F_{P}$ is $a$ homeomorphism of $\mathbb{R}^{d}$ and $B\left(0 ; R_{\varphi}\right)$.

We extend the map $F_{P}^{-1}$ (defined on the set $U_{P}$ ) to the whole ball $B\left(0 ; R_{\varphi}\right)$ assuming that for $t \in B\left(0 ; R_{\varphi}\right) \backslash U_{P} F_{P}^{-1}(t)$ coincides with the unique $s \in A_{P}$ such that $t \in \bar{B}\left(F_{P}(s) ; p(s) \rho_{\varphi}\right)$. Then:
(iv) For any $t \in B\left(0 ; R_{\varphi}\right)$, the unique ( $M, t$ )-parameter of $P$ coincides with $F_{P}^{-1}(t)$.

The next example was considered by [K-D] and [Ch].
Example 2.7. Let $\mathbf{X}=\mathbb{R}^{d}, d \geq 1$, and let $f(s ; x)=|s-x|-|x|$. Consider the functional

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}(|s-x|-|x|) P(d x) \quad \text { for } s \in \mathbb{R}^{d} .
$$

In this case we have $\varphi(\lambda):=\lambda, \lambda>0$, and $\Phi(\lambda):=\varphi^{\prime}(\lambda)=1, \lambda>0$, and the map $F_{P}$, defined by

$$
F_{P}(s):=\int_{\{x \neq s\}} \frac{s-x}{|s-x|} P(d x)
$$

is a subgradient of $f_{P}$. In the case $d=1$ we have $(s-x) /|s-x|=\operatorname{sign}(s-x)$, and

$$
F_{P}(s)=P((-\infty, s])+P((-\infty, s))-1 \quad \text { for } s \in \mathbb{R}^{1}
$$

so $F_{P}$ is just a simple modification of the distribution function of $P$, while $F_{P}^{-1}$ is a simple modification of the quantile function of $P$.

For $d>1$, Proposition 2.6 (with $\rho_{\varphi}=R_{\varphi}=1$ ) gives the main properties of the maps $F_{P}$ and $F_{P}^{-1}$, first studied by [K-D] and [Ch].

Note that properties (iii) and (iv) in Proposition 2.6 immediately imply the existence and the uniqueness of the Haldane's spatial median (if $P$ is not concentrated in a line). See also Milasevic and Ducharme (1987).

Example 2.8. Let $\mathbf{X}=\mathbb{R}^{d}, d \geq 1$, and let $f(s ; x):=|s-x|^{r}-|x|^{r}$ for some $r>1$. We consider the functional

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}\left(|s-x|^{r}-|x|^{r}\right) P(d x) \quad \text { for } s \in \mathbb{R}^{d},
$$

which is well defined for all $P$ with $\int_{\mathbb{R}^{d}}|x|^{r-1} P(d x)<+\infty$. In this case $\varphi(\lambda)=$ $\lambda^{r}, \lambda \geq 0$. We have $R_{\varphi}=+\infty, \rho_{\varphi}=0$, and $\Phi(\lambda)=r \lambda^{r-1}, \lambda \geq 0$. A map $F_{P}$, defined by

$$
F_{P}(s):=r \int_{\{x \neq s\}}(s-x)|s-x|^{r-2} P(d x),
$$

is a subgradient of $f_{P}$. Proposition 2.6 implies the following:
(i) $F_{P}$ is a strictly monotone homeomorphism from $\mathbb{R}^{d}$ onto $\mathbb{R}^{d}$;
(ii) $F_{P}(\lambda v) \rightarrow \infty$ as $\lambda \rightarrow \infty$ for $v \in S^{d-1}$;
(iii) For any $t \in \mathbb{R}^{d}$, the unique ( $M, t$ )-parameter of $P$ coincides with $F_{P}^{-1}(t)$.

In Examples 2.7, 2.8 we have $h(x)=|x|^{r}, r \geq 1$. The Fourier transform of this function is given by the following formula [see Gelfand and Shilov (1964)]:

$$
\tilde{h}(\lambda)=2^{r+d} \pi^{d / 2} \frac{\Gamma((r+d) / 2)}{\Gamma(-r / 2)}|\lambda|^{-r-d} .
$$

If $r \neq 2 k, k=1,2, \ldots$, we have $\tilde{h}(\lambda) \neq 0$ for all $\lambda \neq 0$.
Given $r \geq 1$, denote $\mathscr{P}_{r}$ the set of all Borel probabilities with

$$
\int_{\mathbb{R}^{d}}|x|^{r-1} P(d x)<+\infty .
$$

Theorem 2.5 implies the following statement.
Corollary 2.9. If $r \geq 1$ and $r \neq 2 k, k=1,2, \ldots$, then, for any two probability measures $P, Q \in \mathscr{P}_{r}$, the equality $F_{P}=F_{Q}$ implies $P=Q$.
3. Equivariance and symmetry. In this section we study equivariance of $M$-distribution and $M$-quantile functions with respect to certain groups of transformations.
3.1. Group of translations. Given $\theta \in \mathbb{R}^{d}$, denote by $P_{\theta}$ the $\theta$-translation of the Borel probability measure $P$ in $\mathbb{R}^{d}$.

Proposition 3.1. Let $h: \mathbb{R}^{d} \mapsto \mathbb{R}^{1}$ be a convex function. Suppose that for some $s_{0} \in \mathbb{R}^{d}$ and for all $s \in \mathbb{R}^{d}$ functions $h(s-\cdot)-h\left(s_{0}-\cdot\right)$ are P-integrable. Let

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}\left(h(s-x)-h\left(s_{0}-x\right)\right) P(d x) \quad \text { for } s \in \mathbb{R}^{d} .
$$

Then for all $\theta \in \mathbb{R}^{d}$, (i) $\partial f_{P_{\theta}}(s)=\partial f_{P}(s-\theta)$ for $s \in \mathbb{R}^{d}$; (ii) $\partial f_{P_{\theta}}^{-1}(t)=\theta+\partial f_{P}^{-1}(t)$ for $t \in \mathbb{R}^{d}$.

If $h$ is continuously differentiable in $\mathbb{R}^{d} \backslash\{0\}$, we can choose a point $t_{0} \in$ $\partial h(0)$, and define a map

$$
H(s):= \begin{cases}\nabla h(s), & \text { for } s \neq 0 \\ t_{0}, & \text { for } s=0\end{cases}
$$

Then the map $F_{P}$, defined by

$$
F_{P}(s):=\int_{\mathbb{R}^{d}} H(s-x) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

is a subgradient of $f_{P}$, and Proposition 2.4 gives its properties. In this case we have the following proposition.

Proposition 3.2. (i) For all $\theta \in \mathbb{R}^{d}, F_{P_{\theta}}(s)=F_{P}(s-\theta), s \in \mathbb{R}^{d}$. (ii) For all $\theta \in \mathbb{R}^{d}$, the $(M, t)$-parameter of $P_{\theta}$ is the $\theta$-translation of the $(M, t)$-parameter of $P$ :

$$
F_{P_{\theta}}^{-1}(t)=\theta+F_{P}^{-1}(t) \quad \text { for } t \in B_{h}
$$

3.2. Orthogonal group. Let us consider the function

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}\left(\varphi(|s-x|)-\varphi\left(\left|s_{0}-x\right|\right)\right) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

where $\varphi$ is a convex function on $[0,+\infty), s_{0} \in \mathbb{R}^{d}$ is a fixed point and the functions $\varphi(|s-\cdot|)-\varphi\left(\left|s_{0}-\cdot\right|\right)$ are $P$-integrable for all $s \in \mathbb{R}^{d}$.

Proposition 3.3. The following hold for all orthogonal transformations $O$ of $\mathbb{R}^{d}$ :
(i) $\partial f_{P_{1} O^{-1}}(s)=O \partial f_{P}\left(O^{-1} s\right)$ for $s \in \mathbb{R}^{d}$;
(ii) $\partial f_{P \circ O^{-1}}^{-1}(t)=O \partial f_{P}^{-1}\left(O^{-1} t\right)$ for $t \in \mathbb{R}^{d}$.

Denote

$$
\rho(x ; \lambda):=\sqrt{\left(\lambda-x^{(1)}\right)^{2}+\left(x^{(2)}\right)^{2}+\cdots+\left(x^{(d)}\right)^{2}} \quad \text { for } x \in \mathbb{R}^{d} \text { and } \lambda \geq 0
$$

Proposition 3.4. Let $P$ be a Borel probability measure in $\mathbb{R}^{d}$, spherically symmetric with respect to the center $a \in \mathbb{R}^{d}$. Then we have the following:

$$
\partial f_{P}(s)= \begin{cases}\psi_{P}(|s-a|) \frac{s-a}{|s-a|}, & \text { for } s \neq a  \tag{i}\\ \bar{B}\left(0 ; \rho_{\varphi}\right), & \text { for } s=a\end{cases}
$$

where $\psi_{P}$ is a monotone multivalued map from $(0,+\infty)$ into $\mathbb{R}^{1}$, defined by
(ii)

$$
\begin{gathered}
\psi_{P}(\lambda):=\int_{\mathbb{R}^{d}} \frac{\lambda+a^{(1)}-x^{(1)}}{\rho(x-a ; \lambda)}(\partial \varphi)(\rho(x-a ; \lambda)) P(d x) \\
\partial f_{P}^{-1}(t)= \begin{cases}a+\left(\psi_{P}\right)^{-1}(|t|) \frac{t}{|t|}, & \text { for }|t|>\rho_{\varphi} \\
a, & \text { otherwise }\end{cases}
\end{gathered}
$$

If $\varphi$ is differentiable in $(0,+\infty)$, and $\Phi(\lambda)=\varphi^{\prime}(\lambda), \lambda>0$, then the map $F_{P}$, defined by

$$
F_{P}(s):=\int_{\{x \neq s\}} \Phi(|s-x|) \frac{s-x}{|s-x|} P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

is a subgradient of $f_{P}$. Proposition 2.6 gives the properties of this map. In this case we have

Proposition 3.5. (i) For any orthogonal transformation $O$ of $\mathbb{R}^{d}$, $F_{P \circ O^{-1}}(s)=O F_{P}\left(O^{-1} s\right)$ for $s \in \mathbb{R}^{d}$.
(ii) If $\varphi$ is strictly convex, or it is strictly increasing and $P$ is not concentrated in a line, then $F_{P \circ O^{-1}}^{-1}(t)=O F_{P}^{-1}\left(O^{-1} t\right)$ for $t \in B\left(0 ; R_{\varphi}\right)$.

Proposition 3.6. Let $P$ be a Borel probability measure in $\mathbb{R}^{d}$, spherically symmetric with respect to the center $a \in \mathbb{R}^{d}$. Then we have the following:

$$
F_{P}(s)= \begin{cases}\psi_{P}(|s-a|) \frac{s-a}{|s-a|}, & \text { for } s \neq a  \tag{i}\\ 0, & \text { for } s=a\end{cases}
$$

where $\psi_{P}$ is a nondecreasing function from $(0,+\infty)$ into $\mathbb{R}^{1}$, defined by
(ii)

$$
\begin{aligned}
& \psi_{P}(\lambda):=\int_{\mathbb{R}^{d}} \frac{\lambda+a^{(1)}-x^{(1)}}{\rho(x-a ; \lambda)} \varphi^{\prime}(\rho(x-a ; \lambda)) P(d x) ; \\
& F_{P}^{-1}(t)= \begin{cases}a+\left(\psi_{P}\right)^{-1}(|t|) \frac{t}{|t|}, & \text { for }|t|>\rho_{\varphi}, \\
a, & \text { otherwise. } .\end{cases}
\end{aligned}
$$

Note that, given $a \in \mathbb{R}^{d}$, the transformation $s \mapsto 2 a-s$ of $\mathbb{R}^{d}$ (reflection of the space $\mathbb{R}^{d}$ with respect to the center $a$ ) is the superposition of a translation and the orthogonal transformation $s \mapsto-s$. Suppose that $P$ is symmetric with respect to the reflection about $a$ (which means that $X-a$ and $a-X$ are identically distributed, given a random vector $X$ with distribution $P$ ). In this case, Propositions 3.2 and 3.5 imply that

$$
F_{P}(s)=-F_{P}(2 a-s) \quad \text { for } s \in \mathbb{R}^{d}
$$

and

$$
F_{P}^{-1}(-t)=2 a-F_{P}^{-1}(t) \quad \text { for } t \in B\left(0 ; R_{\varphi}\right) .
$$

3.3. Affine group. Unfortunately, the $M$-distribution and $M$-quantile functions described above are not affine equivariant. In order to define affine equivariant versions of these functions, one should use in their definitions the kernel $f$, which depends on $P$. One of the possible choices of such a kernel is based on standardization of $P$ using the square root of its covariance matrix [Rao (1988) suggested doing this in order to define an affine equivariant version of Hadane's median]. Namely, one can use the kernel

$$
f(s, x):=\varphi\left(|s-x|_{P}\right)-\varphi\left(\left|s_{0}-x\right|_{P}\right),
$$

where $|s|_{P}:=\left|\operatorname{inv}\left(\Sigma_{P}^{1 / 2}\right) s\right|=\left\langle\operatorname{inv}\left(\Sigma_{P}\right) s, s\right\rangle^{1 / 2}, s \in \mathbb{R}^{d}, \Sigma_{P}$ being the covariance operator of measure $P$.
4. Statistical applications of $\boldsymbol{M}$-distributions and $\boldsymbol{M}$-quantiles. In this section we consider some statistical applications of $M$-distribution and $M$-quantile functions, introduced by [K-D] and [Ch]; see our Example 2.7. In this case $\Phi(\lambda) \equiv 1$, so we have

$$
\begin{equation*}
F_{P}(s):=\int_{\{x \neq s\}} \frac{s-x}{|s-x|} P(d x) . \tag{4.1}
\end{equation*}
$$

In addition to the properties of the map $F_{P}$, summarized in Sections 2 and 3, we need the following proposition. Denote $\mathbf{I}_{d}$ the identity operator in $\mathbb{R}^{d}$.

Proposition 4.1. Suppose that $P$ has a bounded density in $\mathbb{R}^{d}$ with $d \geq 2$. Then $F_{P}$ is continuously differentiable in $\mathbb{R}^{d}$ with derivative

$$
F_{P}^{\prime}(s)=\int_{\{x \neq s\}} \frac{1}{|s-x|}\left[\mathbf{I}_{d}-\frac{(s-x) \otimes(s-x)}{|s-x|^{2}}\right] P(d x) .
$$

Moreover, $F_{P}^{\prime}$ is uniformly continuous in $\mathbb{R}^{d}$, and $F_{P}^{\prime}(s)$ is positively definite for all $s \in \mathbb{R}^{d}$.

All asymptotic results below are based on limit theorems for empirical $M$ distribution and $M$-quantile functions, obtained in Section 5.
4.1. $L$-estimators. In the one-dimensional case, $L$-estimators (also called $L$-statistics), based on a sample ( $X_{1}, \ldots, X_{n}$ ), are defined as linear combinations of order statistics $X_{[1]} \leq X_{[2]} \leq \cdots \leq X_{[n]}$, or, more generally, $L_{n}:=\sum_{1}^{n} m_{n, j} h\left(X_{[j]}\right)$ where $m_{n, j}$ are certain weights, and $h$ is a transformation function. It is hard to extend this definition to the multivariate case, since there is no really natural ordering in $\mathbb{R}^{d}, d \geq 1$. But in many cases an $L$ estimator could be represented as a functional of the empirical measure, $L_{n}=$ $L\left(P_{n}\right)$, where the functional $L(P)$ is defined as $L(P):=\int_{0}^{1} h\left(F_{P}^{-1}(t)\right) m(t) d t$ with a certain weight function $m$. Here $F_{P}$ is the distribution function of $P$, and $F_{P}^{-1}$ is its quantile function. It is natural to call $L(P)$ an $L$-parameter of $P$. [Ch] suggested using a similar representation to define $L$-parameters and $L$-estimators in the multivariate case. His suggestion was to replace $F_{P}^{-1}$ for the inverse of the map $F_{P}$, defined by (4.1). In fact, any $M$-quantile function considered in previous sections could be used as well, so one could relate $L$-parameters and $L$-estimators to any $M$-parameter defined by convex minimization. To be specific, we assume that $F_{P}$ is defined by (4.1). In this case

$$
L(P):=L_{h, \mu}(P):=\int_{B_{d}} h\left(F_{P}^{-1}(t)\right) \mu(d t),
$$

where $\mu$ is a signed measure on $B_{d}$ with finite total variation, and $h: \mathbb{R}^{d} \mapsto \mathbb{R}^{m}$ is a vector-valued Borel function, such that the integral in the right-hand side is well defined.

For instance, one could be interested in the uniform distribution $\mu$ on $B(0, r)$ for a certain value $0<r<1$. Such an $L$-parameter could be viewed as a multivariate extension of the trimmed mean. Similarly, linearly trimmed means could be considered.

## Theorem 4.2. Suppose that we have the following conditions:

(i) P has a bounded density in $\mathbb{R}^{d}, d \geq 2$;
(ii) $\mu$ is supported in a ball $B(0, r)$ with some $0<r<1$;
(iii) $h$ is continuously differentiable in a neighborhood of $F_{P}^{-1}(B(0, r))$.

Then an L-estimator $L_{n}:=\int_{B_{d}} h\left(F_{P_{n}}^{-1}(t)\right) \mu(d t)$ is asymptotically normal. The limit distribution of the sequence $n^{1 / 2}\left(L_{n}-L(P)\right)$ is that of the random vector $-\int_{\mathbb{R}^{d}} h^{\prime}(s) \operatorname{inv}\left(F_{P}^{\prime}(s)\right) \xi_{P}(s)\left(\mu \circ F_{P}\right)(d s)$. Moreover, the following asymptotic representation holds:

$$
\begin{aligned}
& L_{n}=L(P)-\int_{\mathbb{R}^{d}} h^{\prime}(s) \operatorname{inv}\left(F_{P}^{\prime}(s)\right)\left(F_{P_{n}}(s)-F_{P}(s)\right)\left(\mu \circ F_{P}\right)(d s)+o_{p}\left(n^{-1 / 2}\right) \\
& \text { as } n \rightarrow \infty .
\end{aligned}
$$

The proof follows directly from Theorem 5.7. Assume that $h(s) \equiv s$, so that $L(P):=\int_{B_{d}} F_{P}^{-1}(t) \mu(d t)$. Assume also that $\mu\left(B_{d}\right)=1$. Let $P$ be a symmetric distribution with center $a \in \mathbb{R}^{d}$ in the sense that the random vectors $X-a$ and $a-X$ are identically distributed, given $X$ has the distribution $P$. If $\mu(A)=$ $\mu(-A)$ for any Borel subset $A$ of $B_{d}$, then it is easy to show that $L(P)=a$. In particular, if $P$ is spherically symmetric about $a \in \mathbb{R}^{d}$, then $L(P)=a$. Thus $L$-statistics described above could be used to estimate the center of symmetry of the distribution $P$.

In the one-dimensional case, $\mu$ is often absolutely continuous on $[0,1]$ : $\mu(d t)=m(t) d t$. Note that in this case, if $P$ is nonatomic, then $P \circ F_{P}^{-1}$ is the uniform distribution on $[0,1]$. It means that, in fact, an $L$-parameter can be expressed as

$$
L(P)=\int_{[0,1]} h\left(F_{P}^{-1}(t)\right) m(t)\left(P \circ F_{P}^{-1}\right)(d t) .
$$

Koltchinskii (1995a) extended this to the multivariate case, defining

$$
L(P):=\int_{B_{d}} h\left(F_{P}^{-1}(t)\right) m(t)\left(P \circ F_{P}^{-1}\right)(d t)=\int_{\mathbb{R}^{d}} h(s) m\left(F_{P}(s)\right) P(d s),
$$

and studied the asymptotic properties of the corresponding $L$-estimators.
4.2. Statistical tests. Suppose that $F_{P}$ is an $M$-distribution function, defined by (4.1). Let $\pi$ be a Borel probability on $\mathbb{R}^{d}$. In order to test the hypothesis $P=\pi$ against the alternative $P \neq \pi$ one could use the test statistics $D_{n}:=n^{1 / 2}\left\|F_{P_{n}}-F_{\pi}\right\|_{\mathbb{R}^{d}}$ or $W_{n}^{2}:=n \int_{\mathbb{R}^{d}}\left|F_{P_{n}}(s)-F_{\pi}(s)\right|^{2} \pi(d s)$, which could be viewed as multivariate extensions of Kolmogorov's and $\omega^{2}$-tests, respectively. If $\pi$ is nonatomic, then Theorem 5.6 implies that the limit distribution of $D_{n}$
coincides with the distribution of $\left\|\xi_{\pi}\right\|_{\mathbb{R}^{d}}$, while the limit distribution of $W_{n}^{2}$ is that of $\int\left|\xi_{\pi}(s)\right|^{2} \pi(d s)$. One could use bootstrap [which is justified in view of the theorem of Giné and Zinn (1990)] in order to determine the rejection region.

Suppose now we have a sample ( $X_{1}, \ldots, X_{n}$ ) from a probability distribution $P$ in $\mathbb{R}^{d}$, and our goal is to test the hypothesis that $P$ is spherically symmetric. Many authors suggested different tests for symmetry [see, for example, Arcones and Giné (1991) for some tests in the one-dimensional case and Beran (1979), Baringhaus (1991) for the multivariate case].

If $P$ is spherically symmetric and $a \in \mathbb{R}^{d}$ is the center of symmetry, then, by Proposition 3.6, we have the following representation of the $M$-distribution function $F_{P}$ :

$$
F_{P}(s)= \begin{cases}\psi_{P}(|s-a|) \frac{s-a}{|s-a|}, & \text { for } s \neq a \\ 0, & \text { for } s=a\end{cases}
$$

where $\psi_{P}(\lambda):=\psi_{P}(a ; \lambda)$ is defined by

$$
\psi_{P}(\lambda):=\int_{\mathbb{R}^{d}} \frac{\lambda+a^{(1)}-x^{(1)}}{\sqrt{\left(\lambda+a^{(1)}-x^{(1)}\right)^{2}+\left(a^{(2)}-x^{(2)}\right)^{2}+\cdots+\left(a^{(d)}-x^{(d)}\right)^{2}}} P(d x) .
$$

It's easy to see that $\psi_{P}(a ; \lambda):=\left\langle F_{P}\left(\lambda e_{1}+a\right), e_{1}\right\rangle$, where $e_{1}$ is the first vector of the canonical basis of $\mathbb{R}^{d}$.

For an arbitrary probability distribution $P$ in $\mathbb{R}^{d}$, define

$$
\tilde{\psi}_{P}(\lambda):=\int_{S^{d-1}}\left\langle F_{P}\left(\lambda v+F_{P}^{-1}(0)\right), v\right\rangle m(d v),
$$

where $m$ denotes the uniform distribution on $S^{d-1}$. Note that for a spherically symmetric $P \tilde{\psi}_{P} \equiv \psi_{P}$.

We introduce a measure of asymmetry of $P$ as follows:

$$
\gamma(P):=\sup _{s \in \mathbb{R}^{d} \mid}\left|F_{P}\left(s+F_{P}^{-1}(0)\right)-\tilde{\psi}_{P}(|s|) \frac{s}{|s|}\right|
$$

It's easy to see that the functional $\gamma(P)$ is invariant with respect to all orthogonal transformations of $\mathbb{R}^{d}$ [i.e., $\left.\gamma\left(P \circ O^{-1}\right)=\gamma(P)\right]$, and $\gamma(P)=0$ iff $P$ is spherically symmetric. Therefore, a test for spherical symmetry of $P$ can be based on the measure of asymmetry $\gamma\left(P_{n}\right)$ of the empirical distribution.

Theorem 4.3. If P has a bounded spherically symmetric density in $\mathbb{R}^{d}$ with $d \geq 2$, then the sequence of distributions of $n^{1 / 2} \gamma\left(P_{n}\right)$ converges weakly to the distribution of the random variable $\left\|\delta_{P}\right\|_{\mathbb{R}^{d}}$, where

$$
\begin{aligned}
& \delta_{P}(s):= \begin{cases}\zeta_{P}(s)-\int_{S^{d-1}}\left\langle\xi_{P}\left(|s| v+F_{P}^{-1}(0)\right), v\right| m(d v) \frac{s}{|s|}, & \text { for } s \neq 0, \\
0, & \text { for } s=0,\end{cases} \\
& \zeta_{P}(s):=\xi_{P}\left(s+F_{P}^{-1}(0)\right)-F_{P}^{\prime}\left(s+F_{P}^{-1}(0)\right) \operatorname{inv}\left(F_{P}^{\prime}\left(F_{P}^{-1}(0)\right)\right) \xi_{P}\left(F_{P}^{-1}(0)\right) .
\end{aligned}
$$

If $P$ is not spherically symmetric, then for all $t \geq 0$,

$$
\lim _{n \rightarrow \infty} \operatorname{Pr}\left\{n^{1 / 2} \gamma\left(P_{n}\right) \leq t\right\}=0 .
$$

We describe a bootstrap version of this test. Any spherically symmetric distribution $P$ is completely characterized by its center $a \in \mathbb{R}^{d}$ and the distribution $\pi$ of the random variable $|X-a|$, given $X$ with distribution $P$. We call $(a, \pi)$ parameters of $P$. Given a sample ( $X_{1}, \ldots, X_{n}$ ) from the unknown distribution $P$ in $\mathbb{R}^{d}$, denote by $\pi_{n}$ the empirical distribution based on the sample $\left(\left|X_{1}-F_{P_{n}}^{-1}(0)\right|, \ldots,\left|X_{n}-F_{P_{n}}^{-1}(0)\right|\right)$. Let $P_{n}^{s}$ be the spherically symmetric distribution in $\mathbb{R}^{d}$ with parameters $\left(F_{P_{n}}^{-1}(0), \pi_{n}\right)$. Define on a probability space $\left(\Omega_{s}, \Sigma_{s}, \operatorname{Pr}_{s}\right)$ an i.i.d. sample $\left(\hat{X}_{1}, \ldots, \hat{X}_{n}\right)$ from the distribution $P_{n}^{s}$. Let $\hat{P}_{n}$ be the empirical distribution of this sample. It could be shown that if $P$ has a bounded density in $\mathbb{R}^{d}$ with $d \geq 2$, then the sequence of distributions (with respect to $\operatorname{Pr}_{s}$ ) of the random variables $n^{1 / 2} \gamma\left(\hat{P}_{n}\right)$ converges weakly to the distribution of $\left\|\delta_{P^{s}}\right\|_{\mathbb{R}^{d}}$ in probability Pr. Here $P^{s}$ denotes the spherically symmetric distribution in $\mathbb{R}^{d}$ with parameters $\left(F_{P}^{-1}(0), \pi\right)$, where $\pi$ is the distribution of the random variable $\left|X-F_{P}^{-1}(0)\right|$, given $X$ with distribution $P$. In particular, if $P$ is spherically symmetric, then $P^{s}=P$. This allows us to construct a bootstrap test for spherical symmetry, based on the test statistic $n^{1 / 2} \gamma\left(P_{n}\right)$. The proofs go beyond this paper [see Koltchinskii and Lang Li (1996)].
4.3. Regression quantiles. Let $(X, Y)$ be a vector in the space $\mathbb{R}^{m} \times \mathbb{R}^{d}=$ $\mathbb{R}^{m+d}$, satisfying the following model:

$$
\begin{equation*}
Y=S_{0} X+\xi, \tag{4.2}
\end{equation*}
$$

where random vectors $X$ and $\xi$ are independent, and $S_{0}$ is a linear operator from $\mathbb{R}^{m}$ into $\mathbb{R}^{d}$. In the case $d=1$, of course, $S_{0} x=\left\langle s_{0}, x\right\rangle$ for some $s_{0} \in \mathbb{R}^{m}$.

Denote by $R$ the distribution of the random vector ( $X, Y$ ), and let $R_{X}, R_{Y}, P$ be the distributions of $X, Y, \xi$, respectively.

The space of all linear operators from $\mathbb{R}^{m}$ into $\mathbb{R}^{d}$ could be identified with $\mathbb{R}^{m d}$ (if $u_{i}, i=1, \ldots, d$ is the canonical basis in $\mathbb{R}^{d}$ and $v_{j}, j=1, \ldots, m$ is the canonical basis in $\mathbb{R}^{m}$, then we take $\left\{u_{i} \otimes v_{j}, i=1, \ldots, d, j=1, \ldots, m\right\}$ as a basis in $\mathbb{R}^{\text {md }}$ ).

We introduce a kernel $g: \mathbb{R}^{m d} \times \mathbb{R}^{m} \times \mathbb{R}^{d} \mapsto \mathbb{R}^{1}, g(S ; x, y):=|S x-y|-|y|$. Clearly, for all $x \in \mathbb{R}^{m}, y \in \mathbb{R}^{d} g$ is convex in $S \in \mathbb{R}^{m d}$. Assuming that $\mathbb{E}|X|=$ $\int_{\mathbb{R}^{m}}|x| R_{X}(d x)<+\infty$, define a convex function $g_{R}: \mathbb{R}^{m d} \mapsto \mathbb{R}^{1}$,

$$
g_{R}(S):=\int_{\mathbb{R}^{m} \times \mathbb{R}^{d}} g(S ; x, y) R(d x, d y)
$$

For $t \in B_{d}$, we define a regression $t$-quantile as a minimal point of the following functional in $S \in \mathbb{R}^{m d}$ :

$$
\mathbb{E}(|S X-Y|-|Y|-\langle S X, t\rangle)=g_{R}(S)-\langle S, t \otimes \mathbb{E} X\rangle
$$

This is an extension of the original definition of Koenker and Basset (1978) for $d=1$.

A subgradient $G_{R}: \mathbb{R}^{m d} \mapsto \mathbb{R}^{m d}$ of the function $g_{R}$ is given by:

$$
G_{R}(S):=\int_{\mathbb{R}^{m+d}} G(S ; x, y) R(d x, d y)
$$

where $G: \mathbb{R}^{m d} \times \mathbb{R}^{m} \times \mathbb{R}^{d}$ is defined by

$$
G(S ; x, y):= \begin{cases}\frac{S x-y}{|S x-y|} \otimes x, & \text { for } S x \neq y \\ 0, & \text { otherwise }\end{cases}
$$

Under condition (4.2),

$$
\begin{equation*}
G_{R}(S)=\mathbb{E} \frac{\left(S-S_{0}\right) X-\xi}{\left|\left(S-S_{0}\right) X-\xi\right|} \otimes X=\mathbb{E} F_{P}\left(\left(S-S_{0}\right) X\right) \otimes X \tag{4.3}
\end{equation*}
$$

Calculations show that

$$
\partial g(S ; x, y)= \begin{cases}\left\{\frac{S x-y}{|S x-y|} \otimes x\right\}, & \text { if } S x \neq y \\ \{u \otimes x:|u| \leq 1\}, & \text { otherwise }\end{cases}
$$

which, in view of Theorem 2.1, implies

$$
\partial g_{R}(S)=G_{R}(s)+\left\{u \otimes \int_{\{(x, y): S x=y\}} x R(d x, d y):|u| \leq 1\right\}
$$

We need some properties of $g_{R}$ and $G_{R}$ (most of them follow from the general properties of $M$-distribution and $M$-quantile functions; see Section 2). Relationship (4.4) below follows from (4.3) and Proposition 4.1.

Denote

$$
B_{R}:=\bigcap_{V \in \mathbb{R}^{m d}, V \neq 0}\{T:\langle V, T\rangle<\mathbb{E}|V X|\}
$$

Then $B_{R}$ is an open convex subset of the ball in $\mathbb{R}^{m d}$ with center 0 and radius $E|X|$.

We introduce the following conditions:

1. $X$ has an absolutely continuous distribution with $\mathbb{E}|X|^{2}<+\infty$;
2. $P$ has a bounded density $p$ in $\mathbb{R}^{d}$;
3. if $d=1$, assume, in addition, that $p$ is continuous and $R_{X}(\{x: p(\langle s, x\rangle)>$ $0\})>0$ for all $s \in \mathbb{R}^{m}$.
Given $S_{1}, S_{2} \in \mathbb{R}^{m d}$ with $S_{1} \neq S_{2}$, denote

$$
\mathscr{L}\left(S_{1}, S_{2} ; x\right):=\left\{S_{1} x+\lambda\left(S_{2}-S_{1}\right) x: \lambda \in \mathbb{R}^{1}\right\} \quad \text { for } x \in \mathbb{R}^{m}
$$

Proposition 4.4. For all $T \in B_{R}$ there exists an $(M, T)$-parameter of $R$. If $d>1$ and for all $S_{1}, S_{2} \in \mathbb{R}^{m d}$ with $S_{1} \neq S_{2}$ we have $R(\{(x, y): y \notin$ $\left.\left.\mathscr{L}\left(S_{1}, S_{2} ; x\right)\right\}\right)>0$, then $g_{R}$ is strictly convex and the $(M, T)$-parameter is unique for all $T \in B_{R}$. If conditions (1)-(3) hold, then $g_{R}$ is a strictly convex continuously differentiable function with gradient $G_{R}$, which is a strictly
monotone one-to-one map from $\mathbb{R}^{\text {md }}$ onto the open convex set $B_{R}$. Moreover, $G_{R}$ is a homeomorphism of $\mathbb{R}^{\text {md }}$ and $B_{R}$, and it is continuously differentiable in $\mathbb{R}^{\text {md }}$ with nonsingular and strictly positive derivative

$$
\begin{equation*}
G_{R}^{\prime}(S)=\mathbb{E} F_{P}^{\prime}\left(\left(S-S_{0}\right) X\right) \otimes X \otimes X \tag{4.4}
\end{equation*}
$$

For all $T \in B_{R}$, the unique ( $M, T$ )-parameter of $R$ coincides with $G_{R}^{-1}(T)$.
Since for $t \in B_{d}$,

$$
\langle V, t \otimes \mathbb{E} X\rangle=\langle t, \mathbb{E} V X\rangle<\mathbb{E}|V X|
$$

we have $t \otimes \mathbb{E} X \in B_{R}$, and it is clear that under conditions (1)-(3) the regression $t$-quantile exists, is unique and coincides with $G_{R}^{-1}(t \otimes \mathbb{E} X)$.

For $S=S_{0}$, we have $G_{R}\left(S_{0}\right)=F_{P}(0) \otimes \mathbb{E} X$, which implies $S_{0}=$ $G_{R}^{-1}\left(F_{P}(0) \otimes \mathbb{E} X\right)$. If either $F_{P}(0)=0$ (which means that the spatial median of $\xi$ is 0 ), or $\mathbb{E} X=0$, then $S_{0}=G_{R}^{-1}(0)$.

Consider a sample $\left(X_{i}, Y_{i}\right), i=1, \ldots, n$ from the distribution $R$, and let $R_{n}$ be the empirical distribution based on this sample. Denote $\bar{X}_{n}$ := $\left(X_{1}+\cdots+X_{n}\right) / n$. Since for $t \in B_{d}$,

$$
\left\langle V, t \otimes \bar{X}_{n}\right\rangle=\left\langle t, n^{-1}\left(V X_{1}+\cdots+V X_{n}\right)\right\rangle<n^{-1}\left(\left|V X_{1}\right|+\cdots+\left|V X_{n}\right|\right)
$$

we have

$$
t \otimes \bar{X}_{n} \in B_{R_{n}}:=\bigcap_{V \in \mathbb{R}^{m d}, V \neq 0}\left\{T:\langle V, T\rangle<n^{-1}\left(\left|V X_{1}\right|+\cdots+\left|V X_{n}\right|\right)\right\} .
$$

Thus, for all $t \in B_{d}$, the empirical regression $t$-quantiles exist. Moreover, for $d>1$, under condition (2) we have $R_{n}\left(\left\{(x, y): y \in \mathscr{L}\left(S_{1}, S_{2} ; x\right)\right\}\right)=0$ for all $S_{1}, S_{2} \in \mathbb{R}^{m d}$ with $S_{1} \neq S_{2}$ a.s. Thus, due to Proposition 4.4, for all $t \in B_{d}$ the empirical regression $t$-quantile exists, is unique and coincides with $G_{R_{n}}^{-1}\left(t \otimes \bar{X}_{n}\right)$.

Let $Q: B_{d} \mapsto \mathbb{R}^{m d}$ be the regression quantile function [which is uniquely defined under the conditions (1)-(3)], and let $Q_{n}: B_{d} \mapsto \mathbb{R}^{m d}$ be a version of the empirical regression quantile function [which always exists and, for $d>1$, under conditions (1), (2) is a.s. unique]. Define

$$
A(t):=\operatorname{inv}\left(G_{R}^{\prime} \circ G_{R}^{-1}\right)(t \otimes \mathbb{E} X) \quad \text { for } t \in B_{d}
$$

THEOREM 4.5. If conditions (1)-(3) hold, then the sequence of empirical regression quantile processes $\left\{n^{1 / 2}\left(Q_{n}-Q\right)(t)\right\}_{n \geq 1}$ converges weakly locally in $B_{d}$ to the Gaussian process

$$
A(t) \int_{R^{m+d}}(t \otimes x-G(Q(t) ; x, y)) W_{R}^{\circ}(d x, d y), \quad t \in B_{d}
$$

Moreover, locally uniformly in $t \in B_{d}$

$$
\begin{equation*}
Q_{n}(t)=Q(t)+A(t)\left(t \otimes\left(\bar{X}_{n}-\mathbb{E} X\right)\right)-A(t)\left(G_{R_{n}}-G_{R}\right) \circ Q(t)+o_{p}\left(n^{-1 / 2}\right) . \tag{4.5}
\end{equation*}
$$

It is worth mentioning that the definition of the regression quantiles and their asymptotic theory can be easily extended to the case of a discrepancy function $g(S ; x, y):=h(S x-y)-h(-y)$ with an arbitrary convex $h: \mathbb{R}^{d} \mapsto$ $\mathbb{R}^{1}$. This extension is also based on the theory of $M$-quantiles developed in this paper. See Bai, Rao and Wu (1992) for the asymptotic properties of $M$ estimators of linear regression parameters under such a discrepancy function.
5. Asymptotics of empirical $M$-distribution and $M$-quantile functions. First we consider asymptotics of general empirical $M$-quantile functions at a given point. Let $f: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{1}$ be a $P$-convex kernel. Assume that it is $P$-integrable for all $s \in \mathbb{R}^{d}$. Consider a convex function $f_{P}$, defined by

$$
f_{P}(s):=\int_{\mathbf{X}} f(s, x) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

Let $F: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{d}$ be such that $F(s, x) \in \partial f(s, x)$ for all $s \in \mathbb{R}^{d}$ and $P$-almost all $x \in \mathbf{X}$. Then

$$
F_{P}(s):=\int_{\mathbf{X}} F(s, x) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

is a subgradient of $f_{P}$ (see Theorem 2.1).
Let $\left\{C_{n}\right\}_{n \geq 1}$ be a sequence of random sets $\omega \mapsto C_{n}(\omega) \subset \mathbb{R}^{d}$, and let $\eta_{n}$ be a sequence of random vectors in $\mathbb{R}^{d}$. Given a sequence $\left\{\delta_{n}\right\}_{n \geq 1}$ of positive numbers, we write

$$
C_{n}=\eta_{n}+o_{p}\left(\delta_{n}\right) \quad \text { as } n \rightarrow \infty
$$

iff for all $\varepsilon>0$,

$$
\limsup _{n \rightarrow \infty} \operatorname{Pr}^{*}\left(\left\{C_{n} \not \subset B\left(\eta_{n} ; \varepsilon \delta_{n}\right)\right\}\right)=0
$$

TheOrem 5.1. Let $t_{0} \in \mathbb{R}^{d}$. Suppose that we have the following:
(i) $P$ has the unique $\left(M, t_{0}\right)$-parameter with respect to $f$ (which means that $f_{P}^{*}$ is differentiable at the point $\left.t_{0}\right)$. We denote it by $s_{0}\left[\right.$ then $\left.F_{P}\left(s_{0}\right)=t_{0}\right]$;
(ii) $\int_{\mathbf{X}}|F(s, x)|^{2} P(d x)<+\infty$ for all $s$ in a neighborhood of $s_{0}$;
(iii) ${F_{P}}_{P}$ is continuously differentiable at the point $s_{0}$ and the derivative $F_{P}^{\prime}\left(s_{0}\right)$ is a positively definite operator in $\mathbb{R}^{d}$.

Then the following representation holds for empirical ( $M, t_{0}$ )-quantiles

$$
\partial f_{P_{n}}^{-1}\left(t_{0}\right)=s_{0}-\operatorname{inv}\left(F_{P}^{\prime}\left(s_{0}\right)\right)\left(F_{P_{n}}\left(s_{0}\right)-F_{P}\left(s_{0}\right)\right)+o_{p}\left(n^{-1 / 2}\right) \quad \text { as } n \rightarrow \infty
$$

If $s_{n}$ is an $\left(M, t_{0}\right)$-quantile of $P_{n}$, then $n^{1 / 2}\left(s_{n}-s_{0}\right)$ is asymptotically normal. Its limit distribution is that of the Gaussian random vector $-\operatorname{inv}\left(F_{P}^{\prime}\left(s_{0}\right)\right) \int_{\mathbf{x}} F\left(s_{0}, x\right) W_{P}^{\circ}(d x)$.

The following corollary is a result of Haberman (1989) [see also Niemiro (1992)].

Corollary 5.2. Suppose that $P$ has the unique $M$-parameter $s_{0}$ with respect to $f$ (which means that $f_{P}^{*}$ is differentiable at point 0 ) and conditions (ii) and (iii) of Theorem 5.1 hold. Then the following representation holds for any $M$-estimator $s_{n}$ :

$$
s_{n}=s_{0}-\operatorname{inv}\left(F_{P}^{\prime}\left(s_{0}\right)\right) F_{P_{n}}\left(s_{0}\right)+o_{p}\left(n^{-1 / 2}\right) \quad \text { as } n \rightarrow \infty .
$$

In particular, $n^{1 / 2}\left(s_{n}-s_{0}\right)$ is asymptotically normal. Its limit distribution is that of the Gaussian random vector $-\operatorname{inv}\left(F_{P}^{\prime}\left(s_{0}\right)\right) \int_{\mathbf{X}} F\left(s_{0}, x\right) W_{P}^{\circ}(d x)$.

We consider now uniform asymptotic results for empirical $M$-distribution and $M$-quantile functions. To do this, we need some facts on empirical processes and on weak convergence (see the Appendix).

Suppose that $F: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{d}$ is a function, such that for all $s \in \mathbb{R}^{d}$ and for all $1 \leq j \leq d$ the $j$ th component $F^{(j)}(s ; \cdot)$ of $F(s ; \cdot)$ is $P$-square integrable: $F^{(j)}(s ; \cdot) \in L_{2}(\mathbf{X} ; d P)$. Consider

$$
F_{P}(s):=\int_{\mathbf{X}} F(s ; x) P(d x)
$$

and

$$
F_{P_{n}}(s):=\int_{\mathbf{X}} F(s ; x) P_{n}(d x)=n^{-1} \sum_{1}^{n} F\left(s ; X_{k}\right) \quad \text { for } s \in \mathbb{R}^{d} .
$$

We call $F$ P-Glivenko-Cantelli iff the class of functions $\mathscr{T}:=\bigcup_{j=1}^{n}\left\{F^{(j)}(s ; \cdot): s \in\right.$ $\left.\mathbb{R}^{d}\right\} \in G C(P)$. We call $F P$-Donsker iff, for any $1 \leq j \leq d, F^{(j)}$ is $\rho_{P}$-continuous on $\mathbb{R}^{d}$ and $\mathscr{F} \in \operatorname{CLT}(P)$.

Let $\phi$ be a nondecreasing continuous function from $[0,+\infty)$ into itself, and let $\psi(s):=\phi(|s|), s \in \mathbb{R}^{d}$. The following facts follow from the properties of empirical processes (see the Appendix and further references there).

Theorem 5.3. If $F / \psi$ is $P$-Glivenko-Cantelli, then a.s. $\left\|F_{P_{n}}-F_{P}\right\|_{\mathbb{R}^{d}, \psi} \rightarrow 0$ as $n \rightarrow \infty$. In particular, a.s. for all compacts $K \subset \mathbb{R}^{d},\left\|F_{P_{n}}-F_{P}\right\|_{K} \rightarrow 0$ as $n \rightarrow \infty$. If $\psi$ is bounded and $F$ is P-Glivenko-Cantelli, then a.s. $\| F_{P_{n}}-$ $F_{P} \|_{\mathbb{R}^{d}} \rightarrow 0$ as $n \rightarrow \infty$. If $F / \psi$ is $P$-Donsker, then

$$
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w} \xi_{P} \text { in } l_{\psi}^{\infty}\left(\mathbb{R}^{d}\right) \text { as } n \rightarrow \infty,
$$

where $\xi_{P}(s):=\int_{\mathbf{X}} F(s, x) W_{P}^{\circ}(d x)$ for $s \in \mathbb{R}^{d}$. In particular,

$$
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w, \text { loc }} \xi_{P} \quad \text { as } n \rightarrow \infty .
$$

If $\psi$ is bounded and $F$ is $P$-Donsker, then

$$
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w} \xi_{P} \quad \text { as } n \rightarrow \infty .
$$

Let now $f: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{1}$ be a $P$-strictly convex kernel. Assume that it is $P$-integrable for all $s \in \mathbb{R}^{d}$. Consider a convex function $f_{P}$, defined by

$$
f_{P}(s):=\int_{\mathbf{X}} f(s, x) P(d x) \quad \text { for } s \in \mathbb{R}^{d} .
$$

Let $F: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{d}$ be such that $F(s, x) \in \partial f(s, x)$ for all $s \in \mathbb{R}^{d}$ and $P$-almost all $x \in \mathbf{X}$. Then

$$
F_{P}(s):=\int_{\mathbf{X}} F(s, x) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

is a subgradient of $f_{P}$ (see Theorem 2.1). Let $\delta_{f}(s ; x)$ be the diameter of the bounded convex set $\partial f(s, x): \delta_{f}(s ; x):=\operatorname{Diam}(\partial f(s ; x))$.

Denote $Q:=F_{P}^{-1}=\left(\nabla f_{P}\right)^{-1}: U_{P} \mapsto D_{P}$ [see properties (5), (6) in Section 2.1].

THEOREM 5.4. Let $S$ be an open subset of $\mathbb{R}^{d}$ and let $T:=F_{P}(S)$. Suppose the following:
(i) $F_{P}$ is continuously differentiable in $S$ with nonsingular derivative $F^{\prime}{ }_{P}(s), s \in S$;
(ii) for some $c>0$, either $\psi \equiv c$, or for all $s \in \mathbb{R}^{d}$ with $|s|$ sufficiently large, $F_{P}(s) \geq c \psi(s) ;$
(iii) $F / \psi$ is $P$-Donsker;
(iv) $\sup _{s \in \mathbb{R}^{d}} \int_{\mathbf{X}} \delta_{f}(s ; x) P_{n}(d x)=o_{P}\left(n^{-1 / 2}\right)$ as $n \rightarrow \infty$;
(v) for any compact set $K \subset B_{P} \lim \sup _{n \rightarrow \infty} \operatorname{Pr}^{*}\left(\left\{K \not \subset B_{P_{n}}\right\}\right)=0$.

Then for any $Q_{n}: \Omega \times T \mapsto \mathbb{R}^{d}$, such that $Q_{n}(t) \in \partial f_{P_{n}}^{-1}(t)$ for $t \in T$ [so that $Q_{n}(t)$ is an ( $M, t$ )-estimator], we have

$$
n^{1 / 2}\left(Q_{n}-Q\right) \rightarrow_{w, \text { loc }} \eta_{P} \quad \text { in } S \text { as } n \rightarrow \infty
$$

where $\eta_{P}:=-\operatorname{inv}\left(F_{P}^{\prime} \circ Q\right)\left(\xi_{P} \circ Q\right)$. Moreover,

$$
n^{1 / 2}\left(Q_{n}-Q+\operatorname{inv}\left(F_{P}^{\prime} \circ Q\right)\left(F_{P_{n}}-F_{P}\right) \circ Q\right) \rightarrow_{p, \text { loc }} 0 \quad \text { in } S \text { as } n \rightarrow \infty
$$

Now we consider in some detail asymptotics of empirical spatial distributions and quantiles. Let $\varphi$ be a convex increasing function on $[0,+\infty)$, continuously differentiable in $(0,+\infty)$. Let $\Phi:=\varphi^{\prime}$. Suppose that for some $C>0$,

$$
\begin{equation*}
\Phi\left(\lambda_{1}+\lambda_{2}\right) \leq C\left[\Phi\left(\lambda_{1}\right)+\Phi\left(\lambda_{2}\right)\right] \text { for } \lambda_{1}, \lambda_{2}>0 \tag{5.1}
\end{equation*}
$$

and for all $c>0$,

$$
\begin{equation*}
\frac{\Phi(\lambda+c)}{\Phi(\lambda)} \rightarrow 1 \quad \text { as } \lambda \rightarrow \infty \tag{5.2}
\end{equation*}
$$

Denote $\psi(s):=\Phi(|s|) \bigvee 1$ for $s \in \mathbb{R}^{d}$. We consider an $M$-distribution function

$$
F_{P}(s):=\int_{\{x \neq s\}} \Phi(|s-x|) \frac{s-x}{|s-x|} P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

which is a subgradient of the convex function

$$
f_{P}(s):=\int_{\mathbb{R}^{d}}(\varphi(|s-x|)-\varphi(|x|)) P(d x) \quad \text { for } s \in \mathbb{R}^{d}
$$

Theorem 5.5. If

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} \Phi(|x|) P(d x)<+\infty, \tag{5.3}
\end{equation*}
$$

then a.s.

$$
\begin{equation*}
\left\|F_{P_{n}}-F_{P}\right\|_{\mathbb{R}^{d}, \psi} \rightarrow 0 \text { as } n \rightarrow \infty . \tag{5.4}
\end{equation*}
$$

In particular, a.s. for all compacts $K \subset \mathbb{R}^{d}$,

$$
\begin{equation*}
\left\|F_{P_{n}}-F_{P}\right\|_{K} \rightarrow 0 \text { as } n \rightarrow \infty . \tag{5.5}
\end{equation*}
$$

If $\Phi(+\infty)<+\infty$, then a.s.

$$
\begin{equation*}
\left\|F_{P_{n}}-F_{P}\right\|_{\mathbb{R}^{d}} \rightarrow 0 \text { as } n \rightarrow \infty \tag{5.6}
\end{equation*}
$$

Theorem 5.6. If $P$ is nonatomic and

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} \Phi^{2}(|x|) P(d x)<+\infty, \tag{5.7}
\end{equation*}
$$

then

$$
\begin{equation*}
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w} \xi_{P} \quad \text { in } l_{\psi}^{\infty}\left(\mathbb{R}^{d}\right) \text { as } n \rightarrow \infty, \tag{5.8}
\end{equation*}
$$

where

$$
\xi_{P}(s):=\int_{\{x \neq s\}} \Phi(|s-x|) \frac{s-x}{|s-x|} W_{P}^{\circ}(d x) \quad \text { for } s \in \mathbb{R}^{d} .
$$

In particular,

$$
\begin{equation*}
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w, \text { loc }} \xi_{P} \quad \text { as } n \rightarrow \infty . \tag{5.9}
\end{equation*}
$$

If $\Phi(+\infty)<+\infty$, then

$$
\begin{equation*}
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w} \xi_{P} \quad \text { as } n \rightarrow \infty . \tag{5.10}
\end{equation*}
$$

Suppose that $P$ is nonatomic and it is not concentrated in a line. Then, by Proposition 2.6, $F_{P}^{-1}$ is a well defined map from $B\left(0 ; R_{\varphi}\right)$ into $\mathbb{R}^{d}$. Denote $Q:=F_{P}^{-1}$.

Theorem 5.7. Let $S$ be an open subset of $\mathbb{R}^{d}$ and let $T:=F_{P}(S)$. If $F_{P}$ is continuously differentiable in $S$ with nonsingular derivative $F_{P}^{\prime}(s), s \in S$, and if condition (5.7) holds, then, for any $Q_{n}$ from $\Omega \times T$ into $\mathbb{R}^{d}$, such that $Q_{n}(t) \in \partial f_{P_{n}}^{-1}(t)$ for all $t \in T$ [so that $Q_{n}(t)$ is an ( $\left.M, t\right)$-parameter of $P$ ],

$$
n^{1 / 2}\left(Q_{n}-Q\right) \rightarrow_{w, \text { loc }} \eta_{P} \quad \text { in } S \text { as } n \rightarrow \infty,
$$

where $\eta_{P}:=-\operatorname{inv}\left(F_{P}^{\prime} \circ F_{P}^{-1}\right)\left(\xi_{P} \circ F_{P}^{-1}\right)$. Moreover,

$$
n^{1 / 2}\left(Q_{n}-Q+\operatorname{inv}\left(F_{P}^{\prime} \circ F_{P}^{-1}\right)\left(F_{P_{n}}-F_{P}\right) \circ Q\right) \rightarrow_{p, \text { loc }} 0 \quad \text { in } S \text { as } n \rightarrow \infty .
$$

Remarks. Since $P$ is nonatomic and it is not concentrated in a line, $P_{n}$ is not concentrated in a line either with probability tending to 1 as $n \rightarrow \infty$. By Proposition 2.6, $F_{P_{n}}^{-1}$ is well defined on $B\left(0 ; R_{\varphi}\right)$ and $Q_{n} \equiv F_{P_{n}}^{-1}$ with probability tending to 1 as $n \rightarrow \infty$. If, moreover, $P(L)=0$ for all straight lines $L$ in $\mathbb{R}^{d}$, then $P_{n}$ is not concentrated in a line and $Q_{n} \equiv F_{P_{n}}^{-1}$ for almost all samples $\left(X_{1}, \ldots, X_{n}\right)$ from $P$ with size $n>2$.

Note that the proof of Theorem 5.7 (see Theorems 5.4 and A.4) includes, in fact, a differentiation of the map $R \mapsto F_{R}^{-1}$. Suppose that $P$ is nonatomic. Given $R \in \mathscr{P}\left(\mathbb{R}^{d}\right)$, denote $P_{\varepsilon}:=(1-\varepsilon) P+\varepsilon R$, where $\varepsilon \in[0,1]$. Let $t \in$ $B\left(0 ; R_{\varphi}\right)$ be a point such that $F_{P}$ is continuously differentiable in a neighborhood of the point $F_{P}^{-1}(t)$ with nonsingular derivative. Then

$$
\frac{F_{P_{\varepsilon}}^{-1}(t)-F_{P}^{-1}(t)}{\varepsilon} \rightarrow-\operatorname{inv}\left(F_{P}^{\prime} \circ F_{P}^{-1}(t)\right)\left(F_{P}-F_{R}\right) \circ F_{P}^{-1}(t) \quad \text { as } \varepsilon \rightarrow 0 .
$$

It follows, in particular, that the influence curve of the functional $P \mapsto F_{P}^{-1}(t)$ is given by the following expression:

$$
\begin{aligned}
& I C\left(x ; P ; F_{P}^{-1}(t)\right)=-\operatorname{inv}\left(F_{P}^{\prime} \circ F_{P}^{-1}(t)\right)\left(F_{P}-F_{\delta_{x}}\right) \circ F_{P}^{-1}(t) \\
& \text { for all } t \neq F_{P}(x), x \in \mathbb{R}^{d} .
\end{aligned}
$$

Koltchinskii (1994a, b) studied Bahadur-Kiefer representations of empirical spatial quantiles.

## APPENDIX

Notation. $\langle\cdot, \cdot\rangle$ denotes the canonical inner product in the finitedimensional space $\mathbb{R}^{d} ;|\cdot|=\langle\cdot, \cdot\rangle^{1 / 2}$ is the corresponding norm; $\|\cdot\|$ (sometimes with indices) is used for the norms of functions or linear transformations (matrices). If $x \in \mathbb{R}^{d}$, then $x^{(j)}$ denotes the $j$ th coordinate of $x$. Similarly, if $F$ is an $\mathbb{R}^{d}$-valued function, $F^{(j)}$ is its $j$ th coordinate.

We denote the ball in $\mathbb{R}^{d}$ with center $s \in \mathbb{R}^{d}$ and radius $R>0$ by $B(s ; R)$. In particular, $B_{d}:=B(0 ; 1) . S^{d-1}$ denotes the unit sphere in $\mathbb{R}^{d}: S^{d-1}:=\{s \in$ $\left.\mathbb{R}^{d}:|s|=1\right\}$. Given a set $A \subset \mathbb{R}^{d}, \bar{A}$ denotes its closure.

We use the notation $\mathscr{L}\left(\mathbb{R}^{d}\right)$ for the space of all linear transformations of $\mathbb{R}^{d}$. Given two vectors $u \in \mathbb{R}^{d}$ and $v \in \mathbb{R}^{m}, u \otimes v$ denotes the linear transformation $x \mapsto\langle v, x\rangle u$ from $\mathbb{R}^{m}$ into $\mathbb{R}^{d}$.

Regularly, we denote by $\nabla f$ the gradient of a function $f$ on an open subset of $\mathbb{R}^{d}$, and by $F^{\prime}$ the derivative of a vector-valued map $F$ from an open subset of $\mathbb{R}^{d}$ into $\mathbb{R}^{d}$. So $\nabla f(s) \in \mathbb{R}^{d}$ and $F^{\prime}(s) \in \mathscr{L}\left(\mathbb{R}^{d}\right)$.

Given an invertible $A \in \mathscr{L}\left(\mathbb{R}^{d}\right)$, the notation $\operatorname{inv}(A)$ will be used for its inverse. Given an $\mathscr{L}\left(\mathbb{R}^{d}\right)$-valued function $\Psi$, defined on a subset of $\mathbb{R}^{d}, \operatorname{inv}(\Psi)$ denotes the function $s \mapsto \operatorname{inv}(\Psi(s))$ (assuming, of course, that $\Psi(s)$ is invertible). If $\zeta$ is a function from a subset of $\mathbb{R}^{d}$ into $\mathbb{R}^{d}$, and if $\Psi$ is an $\mathscr{L}\left(\mathbb{R}^{d}\right)$-valued function defined on the domain of $\zeta$, then the "product" $\Psi \zeta$ is the function $s \mapsto \Psi(s) \zeta(s)$.

We denote by $(\Omega, \Sigma, \operatorname{Pr})$ the main probability space which carries all random elements under consideration; $\mathbb{E}$ denotes the expectation with respect to Pr. Define

$$
\begin{gathered}
\operatorname{Pr}^{*}(A):=\inf \{\operatorname{Pr}(B): B \in \Sigma, B \supset A\} \\
\mathbb{E}^{*} \xi:=\inf \{\mathbb{E} \eta: \eta-\Sigma \text {-measurable, } \eta \geq \xi\} .
\end{gathered}
$$

Given a measurable space $(\mathbf{X}, \mathscr{A}), \mathscr{P}(\mathbf{X})$ denotes the set of all probability measures on $\mathbf{X}$.

Facts on convex analysis. The function $f: \mathbb{R}^{d} \mapsto \mathbb{R}^{1}$ is called convex iff $f\left(\alpha s_{1}+(1-\alpha) s_{2}\right) \leq \alpha f\left(s_{1}\right)+(1-\alpha) f\left(s_{2}\right)$ for all $s_{1}, s_{2} \in \mathbb{R}^{d}$ and $\alpha \in[0,1]$, and it is called strictly convex iff the above inequality is strict for all $s_{1} \neq s_{2}, \alpha \in$ $(0,1)$. Any convex function is continuous on $\mathbb{R}^{d}([\mathrm{Ro}]$, Theorem 10.1).

For $s \in \mathbb{R}^{d}$ and $v \in \mathbb{R}^{d}, v \neq 0$, denote $\Delta f(s ; v ; \lambda):=\lambda^{-1}(f(s+\lambda v)-f(s))$. The directional derivative of $f$ at the point $s$ in the direction $v$ is defined by

$$
d f(s ; v):=\inf _{\lambda \in(0,+\infty)} \Delta f(s ; v ; \lambda)=\lim _{\lambda \downarrow 0} \Delta f(s ; v ; \lambda) .
$$

A set $\partial f(s)$ of all points $t \in \mathbb{R}^{d}$ such that $f\left(s^{\prime}\right) \geq f(s)+\left\langle s^{\prime}-s, t\right\rangle$ for all $s^{\prime} \in \mathbb{R}^{d}$ is called a subdifferential of $f$ at a point $s \in \mathbb{R}^{d}$. A function $f$ is called subdifferentiable at a point $s \in \mathbb{R}^{d}$ iff $\partial f(s) \neq \varnothing$. Any point $t \in \partial f(s)$ is called a subgradient of $f$ at point $s$. If $f$ is differentiable at a point $s$, we denote $\nabla f(s)$ the gradient of $f$ at $s$.

Given a convex function $f$, define a set $D_{f}:=\left\{s \in \mathbb{R}^{d}: d f(s ; v)\right.$ is linear in $\left.v\right\}$. Then $f$ is differentiable at any point $s \in D_{f}$ and $\partial f(s)=\{\nabla f(s)\}, s \in D_{f} . D_{f}$ is a dense subset of $\mathbb{R}^{d}$ and its complement $\mathbb{R}^{d} \backslash D_{f}$ has Lebesgue measure 0. The gradient $\nabla f: s \mapsto \nabla f(s)$ is continuous on $D_{f}$ ([Ro], Theorems 25.2, 25.5).

A map $F: \mathbb{R}^{d} \mapsto \mathbb{R}^{d}$ is called a subgradient of $f$ iff $F(s) \in \partial f(s)$ for all $s \in \mathbb{R}^{d}$.
A map $G$ from $\mathbb{R}^{d}$ into $2^{\mathbb{R}^{d}}$ is called a multivalued map. For instance, $s \mapsto \partial f(s)$ is a multivalued map. The inverse of a multivalued map $G$ is the multivalued map $G^{-1}$ such that $G^{-1}(t)=\left\{s \in \mathbb{R}^{d}: t \in G(s)\right\}$ for all $t \in \mathbb{R}^{d}$.

A multivalued map $G$ is called single-valued iff for any $s \in \mathbb{R}^{d}$ the set $G(s)$ contains only one element. Of course, any such map could be identified with a regular map from $\mathbb{R}^{d}$ into $\mathbb{R}^{d}$.

A multivalued map $G$ is called monotone iff for any $s_{1}, s_{2} \in \mathbb{R}^{d}$ and for any $t_{1} \in G\left(s_{1}\right), t_{2} \in G\left(s_{2}\right)\left\langle s_{2}-s_{1}, t_{2}-t_{1}\right\rangle \geq 0$. It is called strictly monotone iff for any $s_{1}, s_{2} \in \mathbb{R}^{d}$ such that $s_{1} \neq s_{2}$ and for any $t_{1} \in G\left(s_{1}\right), t_{2} \in G\left(s_{2}\right)$, $\left\langle s_{2}-s_{1}, t_{2}-t_{1}\right\rangle>0$.

In particular, a map $F: \mathbb{R}^{d} \mapsto \mathbb{R}^{d}$ is called monotone (strictly monotone) iff the multivalued map $\mathbb{R}^{d} \ni s \mapsto\{F(s)\}$ is monotone (strictly monotone).

A multivalued map $G$ is said to be upper semicontinuous at a point $s \in \mathbb{R}^{d}$ iff for any neighborhood $V$ of the set $G(s)$ there exists a neighborhood $U$ of the point $s$ such that for any $s^{\prime} \in U$ we have $G\left(s^{\prime}\right) \subset V$.

A function $f^{*}: \mathbb{R}^{d} \mapsto \mathbb{R}^{1} \cup\{+\infty\}$, defined by $f^{*}(t):=\sup _{s \in \mathbb{R}^{d}}[\langle s, t\rangle-f(s)]$, is called the Young-Fenchel conjugate of $f$. It is a convex function and $\partial f^{*}=$ $(\partial f)^{-1}([\mathrm{Ro}]$, Theorem 23.5).

The recession function $r_{f}$ of a convex function $f$ on $\mathbb{R}^{d}$ is defined as

$$
r_{f}(v):=\sup _{s \in \mathbb{R}^{d}}[f(s+v)-f(s)] \quad \text { for } v \in \mathbb{R}^{d} .
$$

Empirical processes and weak convergence. Let $S$ be a set, and let ( $L,|\cdot|$ ) be a separable Banach space. Denote by $l^{\infty}(S)=l^{\infty}(S ; L)$ the normed space of all uniformly bounded functions $Y: S \mapsto L$ with norm

$$
\|Y\|_{S}:=\sup _{s \in S}|Y(s)| \quad \text { for } Y \in l^{\infty}(S) .
$$

More generally, if $\psi$ is a positive function on $S$, we define a norm

$$
\|Y\|_{S, \psi}:=\sup _{s \in S} \frac{|Y(s)|}{\psi(s)}
$$

and denote by $l_{\psi}^{\infty}(S)$ the normed space of all functions $Y$ with $\|Y\|_{S, \psi}<+\infty$.
Let $\left\{\xi_{n}\right\}_{n \geq 1}$ be a sequence of random functions $\xi_{n}: S \mapsto L$. We say that $\left\{\xi_{n}\right\}_{n \geq 1}$ converges weakly in the space $l^{\infty}(S)$ to a random function $\xi: S \mapsto L$, or

$$
\xi_{n} \rightarrow_{w} \xi \quad \text { as } n \rightarrow \infty,
$$

iff there exists a separable subspace $\mathscr{N} \subset l^{\infty}(S)$ such that $\xi \in \mathscr{N}$ a.s. and for any bounded continuous functional $\Phi: l^{\infty}(S) \mapsto \mathbb{R}$ we have $\mathbb{E}^{*} \Phi\left(\xi_{n}\right) \rightarrow \mathbb{E} \Phi(\xi)$ as $n \rightarrow \infty$.

The sequence $\left\{\xi_{n}\right\}_{n \geq 1}$ converges weakly in $l_{\psi}^{\infty}(S)$ to a random function $\xi$ iff

$$
\frac{\xi_{n}}{\psi} \rightarrow_{w} \frac{\xi}{\psi} \quad \text { as } n \rightarrow \infty
$$

It's easy to check that if $\xi_{n} \rightarrow_{w} \xi$ in $l_{\psi}^{\infty}(S)$, then the same holds in $l_{\psi_{1}}^{\infty}(S)$ for any $\psi_{1}$ such that $c^{-1} \psi \leq \psi_{1} \leq c \psi$ for some $c>0$.

Let $S$ be an open subset of a metric space. The sequence $\left\{\xi_{n}\right\}_{n \geq 1}$ converges weakly locally in $S$ to a random function $\xi$, or $\xi_{n} \rightarrow_{w, \text { loc }} \xi$ as $n \rightarrow \infty$, iff $\xi$ is a.s. continuous in $S$ and for any compact subset $K$ of $S,\left\{\xi_{n}\right\}_{n \geq 1}$ converges weakly to $\xi$ in the space $l^{\infty}(K)$.

Let $(\mathbf{X}, \mathscr{A}, P)$ be a probability space. Denote $(\Omega, \Sigma, \operatorname{Pr})$ a countable product of copies of $(\mathbf{X}, \mathscr{A}, P)$ times the unit interval $[0,1]$ with Borel $\sigma$-algebra and Lebesgue measure. Let $X_{i}, i \geq 1$ be coordinates on $(\Omega, \Sigma, \operatorname{Pr})$. We define the empirical measure $P_{n}$, based on a sample ( $X_{1}, \ldots, X_{n}$ ), as

$$
P_{n}:=n^{-1} \sum_{1}^{n} \delta_{X_{i}},
$$

where $\delta_{x}$ is a unit point mass at $x$. A signed random measure $Z_{n}:=n^{1 / 2}\left(P_{n}-\right.$ $P$ ) is called an empirical process. Usually $Z_{n}$ is considered as a stochastic process $Z_{n}(f)$, indexed by $\mathscr{A}$-measurable functions on $\mathbf{X}$. Here $\nu(f):=\int_{\mathbf{X}} f d \nu$ for a signed measure $\nu$ on $\mathscr{A}$.

The limit process (in the sense of finite-dimensional distributions) of the sequence $Z_{n}(f), f \in L_{2}(\mathbf{X}, d P)$ of empirical processes is the so called $P$ Brownian bridge, that is, a Gaussian process $W_{P}^{\circ}(f), f \in L_{2}(\mathbf{X}, d P)$ with mean 0 and covariance

$$
\mathbb{E} W_{P}^{\circ}(f) W_{P}^{\circ}(g)=P(f g)-P(f) P(g)
$$

The $P$-Brownian bridge can be also viewed as a stochastic integral $W_{P}^{\circ}(f)=$ $\int_{\mathbf{x}} f d W_{P}^{\circ}$.

A class $\mathscr{F} \subset L_{1}(\mathbf{X}, d P)$ is called $P$-Glivenko-Cantelli, or $\mathscr{F} \in G C(P)$, iff the sequence of random variables $\left\{\left\|P_{n}-P\right\|_{\mathscr{F}}\right\}$ tends to 0 as $n \rightarrow \infty$ a.s.

We introduce the following metric in $L_{2}(\mathbf{X}, d P)$ :

$$
\rho_{P}(f, g):=\left(P(f-g)^{2}-(P(f-g))^{2}\right)^{1 / 2}
$$

A class $\mathscr{T} \subset L_{2}(\mathbf{X}, d P)$ is called $P$-Donsker, or $\mathscr{F} \in \operatorname{CLT}(P)$, iff

$$
Z_{n} \rightarrow_{w} W_{P}^{\circ} \text { in } l^{\infty}(\mathscr{F}) \text { as } n \rightarrow \infty
$$

A measurable real-valued function $E$ on $\mathbf{X}$ is called an envelope of the class $\mathscr{F}$ iff $|f(x)| \leq E(x)$ for all $x \in \mathbf{X}, f \in \mathscr{F}$.

Given a (pseudo)metric space ( $S, d$ ), $N_{d}(S ; \varepsilon)$ denotes the minimal number of balls of radiuses $\varepsilon$, covering $S$. Given $Q \in \mathscr{P}(\mathbf{X})$ and $1 \leq r<+\infty$, denote $d_{Q, r}$ the metric of the space $L_{r}(\mathbf{X} ; d Q)$. We define

$$
N_{r}(\mathscr{F} ; \varepsilon):=N_{r, E}(\mathscr{F} ; \varepsilon):=\sup _{Q \in \mathscr{P}(\mathbf{X})} N_{d_{Q, r}}\left(\mathscr{F} ; \varepsilon\left(\int_{\mathbf{X}} E^{r} d Q\right)^{1 / r}\right) \text { for } \varepsilon>0
$$

We set $H_{r}(\mathscr{F} ; \varepsilon):=\log N_{r}(\mathscr{F} ; \varepsilon)$.
Under some measurability, the conditions $H_{1}(\mathscr{F} ; \varepsilon)<+\infty$ for all $\varepsilon>0$ and $\int_{\mathbf{x}} E d P<+\infty$ imply $\mathscr{T} \in G C(P)$. The conditions

$$
\int_{0}^{\infty} H_{2}^{1 / 2}(\mathscr{F} ; \varepsilon) d \varepsilon<+\infty
$$

and $\int_{\mathbf{X}} E^{2} d P<+\infty$ imply $\mathscr{F} \in \operatorname{CLT}(P)$. For the so-called VC-subgraph classes $\mathscr{F}$, there exist constants $V=V(\mathscr{F})$ and $C=C(V ; r)$ such that

$$
N_{r}(\mathscr{F} ; \varepsilon) \leq C \varepsilon^{-V} \quad \text { for all } \varepsilon>0
$$

In particular, if there exists a finite-dimensional space $\mathscr{H}$ of real-valued functions on $\mathbf{X} \times \mathbb{R}^{1}$, such that any set of the form $\{(x, t): f(x) \geq t \geq 0\}$, or $\{(x, t): f(x) \leq t \leq 0\}$, with some $f \in \mathscr{F}$, can be represented as $\{h \geq 0\}$ for some $h \in \mathscr{H}$, then $\mathscr{T}$ is a VC-subgraph class with $V(\mathscr{F})$ depending only on the dimension of $\mathscr{H}$. See Dudley (1984), Pollard (1990) and van der Vaart and Wellner (1996) for these and further facts on empirical processes and VC-subgraph classes.

## Proofs of the main results.

Proof of Theorem 2.1. The first statement follows immediately from Theorem 8.3.4 of Ioffe and Tihomirov (1974). If $F_{P}(s)$ is a subgradient of $f_{P}$ at a point $s \in \mathbb{R}^{d}$, then $F_{P}(s) \in \partial f_{P}(s)$. The first statement along with the definition of the integral of multivalued maps implies that there exists a map $F: \mathbb{R}^{d} \times \mathbf{X} \mapsto \mathbb{R}^{d}$, such that $F(s, x) \in \partial f(s, x) P$ a.e., $F(s, \cdot)$ is $P$-integrable, and (2.1) holds. Then it is easy to check that the map $F_{P}: \mathbb{R}^{d} \mapsto \mathbb{R}^{d}$, defined by (2.1), is a subgradient of the convex function $f_{P}$. Indeed, we have for all $s \in \mathbb{R}^{d}, v \in S^{d-1}\langle v, F(s, x)\rangle \leq d f(s ; v ; x) P$ a.e. Therefore

$$
\left\langle v, F_{P}(s)\right\rangle=\int_{\mathbf{X}}\langle v, F(s, x)\rangle P(d x) \leq \int_{\mathbf{x}} d f(s ; v ; x) P(d x)=d f_{P}(s ; v)
$$

and the second statement follows.

Proof of Theorem 2.5. First note, that $H^{(j)}=\left(\partial / \partial x_{j}\right) h$ (in the sense of generalized functions). Using the properties of the Fourier transform, we get

$$
\begin{equation*}
\tilde{H}^{(j)}(\lambda)=-i \lambda_{j} \tilde{h}(\lambda) \text { for } \lambda \in \mathbb{R}^{d} . \tag{A.1}
\end{equation*}
$$

We show that for $P \in \mathscr{P}_{h}$,

$$
\begin{equation*}
\tilde{F}_{P}^{(j)}=\tilde{H}^{(j)} \tilde{P}, \quad 1 \leq j \leq d \tag{A.2}
\end{equation*}
$$

Indeed, assume first that $P$ has a density $p$ belonging to the space $\mathscr{\rho}\left(\mathbb{R}^{d}\right)$. Then $F_{P}^{(j)}=H^{(j)} * p$, which implies $\tilde{F}_{P}^{(j)}=\tilde{H}^{(j)} \tilde{p}$, and (A.2) follows.

Next, assume that $P$ is any Borel probability measure with bounded support. Given a random vector $X$ with distribution $P$, let $Y$ be a random vector independent of $X$, with density in $\mathscr{\mathscr { L }}\left(\mathbb{R}^{d}\right)$. Define $Z_{n}:=X+n^{-1} Y$. Let $\Pi_{n}$ be the distribution of $Z_{n}$. Then $\Pi_{n}$ has the density $p_{n}$ belonging to $\mathscr{\rho}\left(\mathbb{R}^{d}\right)$, and we have

$$
\begin{equation*}
\tilde{F}_{\Pi_{n}}^{(j)}=\tilde{H}^{(j)} \tilde{\Pi}_{n}, \quad 1 \leq j \leq d . \tag{A.3}
\end{equation*}
$$

Using continuity of $F_{P}$ in $\mathbb{R}^{d} \backslash A_{P}$ (see Proposition 2.4), it is easy to show that for all $s \notin A_{P}$,

$$
F_{\Pi_{n}}(s)=\int_{\mathbb{R}^{d}} H(s-x) \Pi_{n}(d x)=\mathbb{E} H\left(s-X-n^{-1} Z\right)=\mathbb{E} F_{P}\left(s-n^{-1} Z\right) \rightarrow F_{P}(s),
$$

which implies that $F_{\Pi_{n}}^{(j)}$ converges to $F_{P}^{(j)}$ in $\mathscr{\rho}^{\prime}\left(\mathbb{R}^{d}\right)$ for all $1 \leq j \leq d$. Due to continuity of the Fourier transform, we also have that $\tilde{F}_{\Pi_{n}}^{(j)}$ converges to $\tilde{F}_{P}^{(j)}$ for all $1 \leq j \leq d$ in $\mathscr{\Omega}^{\prime}\left(\mathbb{R}^{d}\right)$. On the other hand, $\tilde{\Pi}_{n}$ converges to $\tilde{P}$. Passing to the limit in (A.3) as $n \rightarrow \infty$, we get (A.2) for an arbitrary $P$ with bounded support.

In the general case, given an $X$ with distribution $P \in \mathscr{P}_{h}$, denote $X_{n}$ the random vector obtained by conditioning of $X$ on $|X| \leq n$, and let $\Pi_{n}$ be the
distribution of $X_{n}$. Then (A.3) holds for such a $\Pi_{n}$. Under the condition $P \in \mathscr{P}_{h}$, we have for all $s \in \mathbb{R}^{d}$,

$$
F_{\Pi_{n}}(s)=\frac{1}{P\{x:|x| \leq n\}} \int_{\{x:|x| \leq n\}} H(s-x) P(d x) \rightarrow F_{P}(s) \quad \text { as } n \rightarrow \infty .
$$

It follows that for all $1 \leq j \leq d F_{\Pi_{n}}^{(j)}$ converges to $F_{P}^{(j)}$ in $\mathscr{\rho}^{\prime}\left(\mathbb{R}^{d}\right)$ as $n \rightarrow \infty$, which implies that $\tilde{F}_{\Pi_{n}}^{(j)}$ converges to $\tilde{F}_{P}^{(j)}$ in $\mathscr{\Omega}^{\prime}\left(\mathbb{R}^{d}\right)$. Since we also have $\tilde{\Pi}_{n}$ converges to $\tilde{P}$, by passing to the limit in (A.3), we get (A.2) in the general case.

Now, if we have two Borel probabilities $P$ and $Q$ such that $F_{P} \equiv F_{Q}$, we get, for all $1 \leq j \leq d, \tilde{F}_{P}^{(j)}=\tilde{F}_{Q}^{(j)}$. Therefore, by (A.1) and (A.2), $\lambda_{j} \tilde{h}(\lambda)(\tilde{P}(\lambda)-$ $\tilde{Q}(\lambda)) \equiv 0$. Under the condition $\tilde{h} \neq 0$ a.e. in $\mathbb{R}^{d}$, we have $\tilde{P}=\tilde{Q}$ a.e., and since $\tilde{P}$ and $\tilde{Q}$ are both continuous, we can conclude that $\tilde{P}=\tilde{Q}$ and $P=Q$.

Proof of Proposition 2.6. Almost everything follows from Proposition 2.4. One need prove just the second statement of (i), and (ii ). Let $s \in A_{P}$. Clearly, there exists a probability measure $Q$ on $\mathbb{R}^{d}$ such that $Q(\{s\})=0$ and $P=p(s) \delta_{s}+(1-p(s)) Q$, where $\delta_{s}$ is the unit mass at the point $s$. We have $F_{P}=p(s) F_{\delta_{s}}+(1-p(s)) F_{Q}$. It follows that

$$
F_{P}(s+\varepsilon v)=p(s) \Phi(\varepsilon) v+(1-p(s)) F_{Q}(s+\varepsilon v) .
$$

Since $F_{Q}$ is continuous at the point $s$ and $\Phi(0+)=\rho_{\varphi}$, we get

$$
F_{P}(s+\varepsilon v) \rightarrow p(s) \rho_{\varphi} v+(1-p(s)) F_{Q}(s) .
$$

It remains to note that since $F_{\delta_{s}}(s)=0$, we have $F_{P}(s)=(1-p(s)) F_{Q}(s)$, which proves (i).

Part (ii) follows from the representation

$$
F_{P}(\lambda v)=\int_{\{x \neq \lambda v\}} \Phi\left(\lambda\left|v-\frac{x}{\lambda}\right|\right) \frac{v-x / \lambda}{|v-x / \lambda|} P(d x),
$$

the fact that $\Phi(+\infty)=R_{\varphi}$ and the dominated convergence theorem.
Proof of Proposition 3.1. Indeed, for all $s \in \mathbb{R}^{d}$,

$$
\begin{aligned}
f_{P_{\theta}}(s)= & \int_{\mathbb{R}^{d}}\left(h(s-x)-h\left(s_{0}-x\right)\right) P_{\theta}(d x) \\
= & \int_{\mathbb{R}^{d}}\left(h(s-\theta-x)-h\left(s_{0}-\theta-x\right)\right) P(d x) \\
= & \int_{\mathbb{R}^{d}}\left(h(s-\theta-x)-h\left(s_{0}-x\right)\right) P(d x) \\
& +\int_{\mathbb{R}^{d}}\left(h\left(s_{0}-x\right)-h\left(s_{0}-\theta-x\right)\right) P(d x) \\
= & f_{P}(s-\theta)+f_{P_{\theta}}\left(s_{0}+\theta\right) .
\end{aligned}
$$

To prove (i), note that $t \in \partial f_{P_{\theta}}(s)$ iff for all $v \in \mathbb{R}^{d}, v \neq 0$ we have

$$
\langle v, t\rangle \leq d f_{P_{\theta}}(s ; v)=d f_{P}(s-\theta ; v) .
$$

It means that $t \in \partial f_{P}(s-\theta)$. Then we have for all $t \in \mathbb{R}^{d}$,

$$
\begin{aligned}
f_{P_{\theta}}^{*}(t) & =\sup _{s \in \mathbb{R}^{d}}\left[\langle s, t\rangle-f_{P_{\theta}}(s)\right]=\sup _{s \in \mathbb{R}^{d}}\left[\langle s, t\rangle-f_{P}(s-\theta)-f_{P_{\theta}}\left(s_{0}+\theta\right)\right] \\
& =\sup _{s \in \mathbb{R}^{d}}\left[\langle s+\theta, t\rangle-f_{P}(s)\right]-f_{P_{\theta}}\left(s_{0}+\theta\right) \\
& =\langle\theta, t\rangle+\sup _{s \in \mathbb{R}^{d}}\left[\langle s, t\rangle-f_{P}(s)\right]-f_{P_{\theta}}\left(s_{0}+\theta\right) \\
& =\langle\theta, t\rangle+f_{P}^{*}(t)-f_{P_{\theta}}\left(s_{0}+\theta\right) .
\end{aligned}
$$

Finally, $s \in \partial f_{P_{\theta}}^{*}(t)=\partial f_{P_{\theta}}^{-1}(t)$ iff for all $v \in \mathbb{R}^{d}, v \neq 0$

$$
\langle s, v\rangle \leq d f_{P_{\theta}}^{*}(t ; v)=\langle\theta, v\rangle+d f_{P}^{*}(t ; v),
$$

which is equivalent to

$$
\langle s-\theta, v\rangle \leq d f_{P}^{*}(t ; v) .
$$

It means that $s-\theta \in \partial f_{P}^{*}(t)=\partial f_{P}^{-1}(t)$, so (ii) holds.
Proof of Proposition 3.3. Note that $\partial f_{P}(s)$ does not depend on the choice of $s_{0}$ in the definition of $f_{P}$. So, we could assume $s_{0}=0$ and $f_{P}(s):=$ $\int_{\mathbb{R}^{d}}(\varphi(|s-x|)-\varphi(|x|)) P(d x)$. For an orthogonal transformation $O$ and $s \in \mathbb{R}^{d}$, we get

$$
\begin{aligned}
f_{P \circ O^{-1}}(s) & =\int_{\mathbb{R}^{d}}(\varphi(|s-x|)-\varphi(|x|)) P \circ O^{-1}(d x) \\
& =\int_{\mathbb{R}^{d}}(\varphi(|s-O x|)-\varphi(|O x|)) P(d x) \\
& =\int_{\mathbb{R}^{d}}\left(\varphi\left(\left|O^{-1} s-x\right|\right)-\varphi(|x|)\right) P(d x)=f_{P}\left(O^{-1} s\right) .
\end{aligned}
$$

Next, for $t \in \mathbb{R}^{d}$, we have

$$
\begin{aligned}
f_{P \circ O^{-1}}^{*}(t) & =\sup _{s \in \mathbb{R}^{d}}\left[\langle s, t\rangle-f P \circ O^{-1}(s)\right]=\sup _{s \in \mathbb{R}^{d}}\left[\langle s, t\rangle-f_{P}\left(O^{-1} s\right)\right] \\
& =\sup _{s \in \mathbb{R}^{d}}\left[\left\langle O^{-1} s, O^{-1} t\right\rangle-f_{P}\left(O^{-1} s\right)\right] \\
& =\sup _{s \in \mathbb{R}^{d}}\left[\left\langle s, O^{-1} t\right\rangle-f_{P}(s)\right]=f_{P}^{*}\left(O^{-1} t\right) .
\end{aligned}
$$

Relationships (i) and (ii) now follow by subdifferentiation.
Proof of Proposition 3.4. Without loss of generality we can and do assume that $s_{0}=a=0$ (since subdifferentials $\partial f_{P}$ and $\partial f_{P}^{*}$ do not depend on the choice of $s_{0}$ ). Since $P$ is invariant with respect to any orthogonal transformation $O$, it follows from the proof of Proposition 3.3 that $f_{P \circ O^{-1}}(s)=f_{P}\left(O^{-1} s\right)=$ $f_{P}(s), s \in \mathbb{R}^{d}$, so the convex function $f_{P}$ is invariant with respect to all orthogonal transformations. It means (see [Ro], page 110) that there exists a lower semicontinuous convex nondecreasing function $\varphi_{P}$ on $[0,+\infty)$ such that
$f_{P}(s)=\varphi_{P}(|s|)$. Let $\left(e_{1}, \ldots, e_{d}\right)$ be an orthonormal basis in $\mathbb{R}^{d}$. Given a vector $x \in \mathbb{R}^{d}$, let $\left(x_{1}, \ldots, x_{d}\right)$ be its coordinates in this basis. Then for $\lambda \geq 0$,

$$
\varphi_{P}(\lambda)=f_{P}\left(\lambda e_{1}\right)=\int_{\mathbb{R}^{d}} \varphi\left(\left|\lambda e_{1}-x\right|\right) P(d x)=\int_{\mathbb{R}^{d}} \varphi(\rho(x ; \lambda)) P(d x) .
$$

Subdifferentiating this relationship, we get (i).
Now we use the formula for the conjugate of the convex function $f_{P}(s)=$ $\varphi_{P}(|s|)$ (see [Ro], page 110) to get $f_{P}^{*}(t)=\varphi_{P}^{+}(|t|), t \in \mathbb{R}^{d}$, where the function $\varphi_{P}^{+}$is a monotone conjugate of $\varphi_{P}$ :

$$
\varphi_{P}^{+}(\lambda):=\sup _{\alpha \geq 0}\left[\lambda \alpha-\varphi_{P}(\alpha)\right] \text { for } \lambda \geq 0 .
$$

It is a lower semicontinuous convex nondecreasing function such that for all $\lambda \geq 0, \partial \varphi_{P}^{+}(\lambda)=\left(\partial \varphi_{P}\right)^{-1}(\lambda)$. Subdifferentiating the relationship for $f_{P}^{*}$ gives (ii).

In order to prove Theorem 5.1 we need a couple of facts about convex functions and their minimal points. The next lemma is just a stochastic version of a well-known fact of convex analysis (see [Ro], Theorem 10.8).

Lemma A.1. Let $f_{n}: \mathbb{R}^{d} \times \Omega \mapsto \mathbb{R}^{1}$ be a sequence of random functions, convex on $\mathbb{R}^{d}$ for $\operatorname{Pr}$-almost all $\omega \in \Omega$. Suppose that for all $s \in \mathbb{R}^{d}, f_{n}(s) \rightarrow f(s)$ as $n \rightarrow \infty$ in probability. Then for all compacts $K \subset \mathbb{R}^{d}\left\|f_{n}-f\right\|_{K} \rightarrow 0$ as $n \rightarrow \infty$ in probability.

Recall that $\partial f^{-1}(0)$ coincides with the set of all minimal points of $f$.
Lemma A.2. Let $f: \mathbb{R}^{d} \mapsto \mathbb{R}^{1}$ be a convex function and let $g: \mathbb{R}^{d} \mapsto \mathbb{R}^{1}$ be a quadratic function, defined by $g(s):=\langle s, a\rangle+\frac{1}{2}\langle A s, s\rangle, s \in \mathbb{R}^{d}$, where $a \in \mathbb{R}^{d}$ is a vector, and $A$ is a symmetric positively definite operator in $\mathbb{R}^{d}$ such that $\langle A s, s\rangle \geq c|s|^{2}$ for all $s \in \mathbb{R}^{d}$ with a constant $c>0$. Let $b:=-\operatorname{inv}(A) a$. Take $a$ number $r>|b|$. Denote $\delta:=\left(6 c^{-1}\|f-g\|_{\bar{B}(0 ; r)}\right)^{1 / 2}$. If $\delta \leq r-|b|$, then $\partial f^{-1}(0)$ belongs to the ball $\bar{B}(b ; \delta)$.

Proof. The unique minimum of $g$ is at the point $b$. By the Taylor expansion,

$$
g(s)=g(b)+\frac{1}{2}\langle A(s-b), s-b\rangle \quad \text { for } s \in \mathbb{R}^{d} .
$$

If $|s-b|=\delta$, then

$$
g(s) \geq g(b)+\frac{c}{2} \delta^{2}=g(b)+3\|f-g\|_{\bar{B}(0, r)} .
$$

Since for all $s$ with $|s-b|=\delta \bar{B}(b ; \delta) \subset \bar{B}(0 ; r), g(s) \leq f(s)+\|f-g\|_{\bar{B}(0, r)}$ and $g(b) \geq f(b)-\|f-g\|_{\bar{B}(0, r)}$. Therefore, for all $s$ on the sphere $|s-b|=\delta$, we have $f(s) \geq f(b)+\|f-g\|_{\bar{B}(0, r)}$. Since $f$ is convex, any minimal point of $f$ must be in the ball $\bar{B}(b, \delta)$.

The next statement is very useful in the study of asymptotics of statistical estimators defined by convex minimization. In fact, its versions have already been used by Pollard (1988) and Niemiro (1992). Let $f_{n}: \mathbb{R}^{d} \times \Omega \mapsto \mathbb{R}^{1}$ be a sequence of random functions. Assume that $f_{n}(s, \omega)$ is convex in $s$ on $\mathbb{R}^{d}$ for Pr-almost all $\omega \in \Omega$. Suppose that $\mathbb{E}\left|f_{n}(s)\right|<+\infty$, and there exists a function $f: \mathbb{R}^{d} \mapsto \mathbb{R}^{1}$ such that $\mathbb{E} f_{n}(s)=f(s)$, for all $s \in \mathbb{R}^{d}$ and all $n \geq 1$. Let $F_{n}: \mathbb{R}^{d} \times \Omega \mapsto \mathbb{R}^{d}$ be a subgradient of $f_{n}$ for Pr-almost all $\omega \in \Omega$. Then (see Theorem 2.1) $F(s):=\mathbb{E} F_{n}(s)$ is a subgradient of $f$ for all $s \in \mathbb{R}^{d}$.

Given $s_{0} \in \mathbb{R}^{d}$, denote

$$
\begin{aligned}
r_{n}\left(s_{0} ; u\right) & :=f_{n}\left(s_{0}+u\right)-f_{n}\left(s_{0}\right)-\left\langle F_{n}\left(s_{0}\right), u\right\rangle \text { for } u \in \mathbb{R}^{d}, \\
r\left(s_{0} ; u\right) & :=f\left(s_{0}+u\right)-f\left(s_{0}\right)-\left\langle F\left(s_{0}\right), u\right\rangle \text { for } u \in \mathbb{R}^{d} .
\end{aligned}
$$

## Theorem A.3. Suppose the following:

(i) the function $f$ has its unique minimum at the point $s_{0} \in \mathbb{R}^{d}$;
(ii) if is twice continuously differentiable at this point, so that the first derivative $f^{\prime}\left(s_{0}\right)=F\left(s_{0}\right)=0$ and the second derivative $f^{\prime \prime}\left(s_{0}\right)=F^{\prime}\left(s_{0}\right)$ is a positively definite operator in $\mathbb{R}^{d}$;
(iii) $r_{n}\left(s_{0} ; u / n^{1 / 2}\right)-r\left(s_{0} ; u / n^{1 / 2}\right)=o_{p}\left(n^{-1}\right)$ as $n \rightarrow \infty$ for all $u \in \mathbb{R}^{d}$;
(iv) the sequence $n^{1 / 2} F_{n}\left(s_{0}\right)$ is stochastically bounded.

Then the following asymptotic representation holds for the minimal set of $f_{n}$ :

$$
\partial f_{n}^{-1}(0)=\partial f_{n}^{*}(0)=s_{0}-\operatorname{inv}\left(F^{\prime}\left(s_{0}\right)\right) F_{n}\left(s_{0}\right)+o_{p}\left(n^{-1 / 2}\right) \quad \text { as } n \rightarrow \infty .
$$

Proof. Since $f$ is twice continuously differentiable at the point $s_{0}$, we have for all $u \in \mathbb{R}^{d}$,

$$
\begin{equation*}
n r\left(s_{0} ; \frac{u}{n^{1 / 2}}\right)=\frac{1}{2}\left\langle F^{\prime}\left(s_{0}\right) u, u\right\rangle+o(1) \quad \text { as } n \rightarrow \infty . \tag{A.4}
\end{equation*}
$$

It follows from conditions (ii), (iii) and (A.4) that for all $u \in \mathbb{R}^{d}$,

$$
n\left(f_{n}\left(s_{0}+\frac{u}{n^{1 / 2}}\right)-f_{n}\left(s_{0}\right)-\left\langle F_{n}\left(s_{0}\right), \frac{u}{n^{1 / 2}}\right\rangle\right)=\frac{1}{2}\left\langle F^{\prime}\left(s_{0}\right) u, u\right\rangle+o_{p}(1) \text { as } n \rightarrow \infty .
$$

Moreover, Lemma A. 1 implies that convergence in the last relationship is uniform over compacts in $\mathbb{R}^{d}$. Denote

$$
\tilde{f}_{n}(u):=n\left(f_{n}\left(s_{0}+\frac{u}{n^{1 / 2}}\right)-f_{n}\left(s_{0}\right)\right) \quad \text { for } u \in \mathbb{R}^{d}
$$

and

$$
\tilde{g}_{n}(u):=n^{1 / 2}\left\langle F_{n}\left(s_{0}\right), u\right\rangle+\frac{1}{2}\left\langle F^{\prime}\left(s_{0}\right) u, u\right\rangle \quad \text { for } u \in \mathbb{R}^{d} .
$$

Take $\varepsilon>0$. By condition (iv) there exists $r>0$ such that

$$
\left|n^{1 / 2} \operatorname{inv}\left(F^{\prime}\left(s_{0}\right)\right) F_{n}\left(s_{0}\right)\right|<r / 2
$$

with probability at least $1-\varepsilon$. On the other hand, we have

$$
\delta_{n}^{2}:=\left\|\tilde{f}_{n}-\tilde{g}_{n}\right\|_{\bar{B}(0 ; r)}=o_{p}(1) \quad \text { as } n \rightarrow \infty .
$$

By Lemma A. 2 the set $\partial \tilde{f}_{n}^{-1}(0)$ of all minimal points of $\tilde{f}_{n}$ belongs (with probability tending to 1 ) to a ball with center $-n^{1 / 2} \operatorname{inv}\left(F^{\prime}\left(s_{0}\right)\right) F_{n}\left(s_{0}\right)$ and with radius of the order $\delta_{n}$. It remains to observe that $\partial f_{n}^{-1}\left(t_{0}\right)=s_{0}+n^{-1 / 2} \partial \tilde{f}_{n}^{-1}(0)$ and the result follows.

Proof of Theorem 5.1. We apply Theorem A. 3 to the following functions:

$$
\begin{aligned}
f_{n}(s) & :=f_{P_{n}}(s)-\left\langle s, t_{0}\right\rangle, \quad F_{n}(s):=F_{P_{n}}(s)-t_{0}, \\
f(s) & :=f_{P}(s)-\left\langle s, t_{0}\right\rangle, \quad F(s):=F_{P}(s)-t_{0} \quad \text { for } s \in \mathbb{R}^{d} .
\end{aligned}
$$

Clearly, $s_{0}$ is the unique minimal point of $f$, and conditions (i) and (ii) of Theorem A. 3 hold. Due to the central limit theorem, condition (iv) of Theorem A. 3 follows from condition (ii) of Theorem 5.1, since $n^{1 / 2} F_{n}\left(s_{0}\right)=n^{1 / 2}\left(F_{P_{n}}\left(s_{0}\right)-\right.$ $\left.F_{P}\left(s_{0}\right)\right)$. To check condition (iii) of Theorem A.3, note that

$$
\mathbb{E} r_{n}\left(s_{0} ; \frac{u}{n^{1 / 2}}\right)=r\left(s_{0} ; \frac{u}{n^{1 / 2}}\right),
$$

and

$$
\begin{align*}
& \operatorname{Var}\left(n r_{n}\left(s_{0} ; \frac{u}{n^{1 / 2}}\right)\right) \\
&=\operatorname{Var}\left(\sum_{1}^{n}\left[f\left(s_{0}+\frac{u}{n^{1 / 2}}, X_{i}\right)-f\left(s_{0}, X_{i}\right)-\left\langle F\left(s_{0}, X_{i}\right), \frac{u}{n^{1 / 2}}\right\rangle\right]\right) \\
& \quad=\sum_{1}^{n} \operatorname{Var}\left(f\left(s_{0}+\frac{u}{n^{1 / 2}}, X_{i}\right)-f\left(s_{0}, X_{i}\right)-\left\langle F\left(s_{0}, X_{i}\right), \frac{u}{n^{1 / 2}}\right\rangle\right)  \tag{A.5}\\
& \quad \leq n \int_{\mathbf{X}}\left(f\left(s_{0}+\frac{u}{n^{1 / 2}}, x\right)-f\left(s_{0}, x\right)-\left\langle F\left(s_{0}, x\right), \frac{u}{n^{1 / 2}}\right\rangle\right)^{2} P(d x) .
\end{align*}
$$

By the definition of the subgradient of convex functions, we have the inequalities

$$
f\left(s_{0}+u, x\right)-f\left(s_{0}, x\right) \geq\left\langle F\left(s_{0}, x\right), u\right\rangle
$$

and

$$
f\left(s_{0}+u, x\right)-f\left(s_{0}, x\right) \leq\left\langle F\left(s_{0}+u, x\right), u\right\rangle
$$

which imply

$$
0 \leq f\left(s_{0}+u, x\right)-f\left(s_{0}, x\right)-\left\langle F\left(s_{0}, x\right), u\right\rangle \leq\left\langle F\left(s_{0}+u, x\right)-F\left(s_{0}, x\right), u\right\rangle .
$$

Thus

$$
\begin{aligned}
0 & \leq f\left(s_{0}+\frac{u}{n^{1 / 2}}, x\right)-f\left(s_{0}, x\right)-\left\langle F\left(s_{0}, x\right), \frac{u}{n^{1 / 2}}\right\rangle \\
& \leq n^{-1 / 2}\left\langle F\left(s_{0}+\frac{u}{n^{1 / 2}}, x\right)-F\left(s_{0}, x\right), u\right\rangle
\end{aligned}
$$

and it follows from (A.5) that

$$
\begin{equation*}
\operatorname{Var}\left(n r_{n}\left(s_{0} ; \frac{u}{n^{1 / 2}}\right)\right) \leq \mathbb{E}\left\langle F\left(s_{0}+\frac{u}{n^{1 / 2}}, x\right)-F\left(s_{0}, x\right), u\right\rangle^{2} . \tag{A.6}
\end{equation*}
$$

By the monotonicity of the subgradient of a convex function, the sequence of nonnegative functions

$$
\gamma_{n}(\cdot)=\left\langle F\left(s_{0}+\frac{u}{n^{1 / 2}}, \cdot\right)-F\left(s_{0}, \cdot\right), u\right\rangle, \quad n \geq 1
$$

is nonincreasing. Denote $\gamma: \mathbf{X} \mapsto[0,+\infty)$ its limit. Then

$$
\begin{align*}
\int_{\mathbf{X}} \gamma(x) P(d x) & =\lim _{n \rightarrow \infty} \int_{\mathbf{x}}\left\langle F\left(s_{0}+\frac{u}{n^{1 / 2}}, x\right)-F\left(s_{0}, x\right), u\right\rangle P(d x)  \tag{A.7}\\
& =\lim _{n \rightarrow \infty}\left\langle F_{P}\left(s_{0}+\frac{u}{n^{1 / 2}}\right)-F_{P}\left(s_{0}\right), u\right\rangle=0,
\end{align*}
$$

since $F_{P}$ is continuous at the point 0 . But $\gamma$ is nonnegative, so $\gamma=0 P$-a.e. Condition (ii) implies, by dominated convergence, that $\int_{\mathbf{X}} \gamma_{n}^{2} d P \rightarrow 0, n \rightarrow \infty$. Thus (A.6) and (A.7) imply condition (iii) of Theorem A.3. Since $\partial f_{n}^{-1}(0)=$ $\partial f_{P_{n}}^{-1}\left(t_{0}\right)$ and $F_{n}\left(s_{0}\right)=F_{P_{n}}\left(s_{0}\right)-F_{P}\left(s_{0}\right)$, this completes the proof.

We need one fact on asymptotics of inverses of random functions. Let $S \subset \mathbb{R}^{d}$ and $H: S \mapsto \mathbb{R}^{d}$. Given $\varepsilon>0$, a map $J$ from a subset $T \subset \mathbb{R}^{d}$ into $S$ will be called an $\varepsilon$-inverse of $H$ on $T$ iff $|t-H(J(t))|<\varepsilon$ for all $t \in T$. Let $\mathscr{J}_{H}(T ; S ; \varepsilon)$ be the set of all $\varepsilon$-inverses on $T$ of the map $H$ from $S$ into $\mathbb{R}^{d}$. The following statement was proved by Koltchinskii (1995b) (see Theorems 2.5 and 2.7 there).

ThEOREM A.4. Let $\left\{a_{n}\right\}_{n \geq 1}$ be a sequence of positive real numbers with $a_{n} \rightarrow \infty, n \rightarrow \infty$. Let $S$ and $T$ be open subsets of $\mathbb{R}^{d}$, and $G$ be a one-toone map from $S$ onto $T$, continuously differentiable in $S$ with nonsingular derivative $G^{\prime}(s), s \in \mathbb{R}^{d}$. Let $\left\{\varepsilon_{n}\right\}_{n \geq 1}$ be a sequence of positive real numbers with $\varepsilon_{n}=o\left(a_{n}^{-1}\right)$ as $n \rightarrow \infty$. Let $G_{n}$ be functions from $\Omega \times S$ into $\mathbb{R}^{d}$. Suppose there exists a stochastic process $\xi: S \mapsto \mathbb{R}^{d}$ such that

$$
a_{n}\left(G_{n}-G\right) \rightarrow_{w, \text { loc }} \xi \quad \text { in } S \text { as } n \rightarrow \infty .
$$

Let $J_{n}$ be a function from $\Omega \times T$ into $\mathbb{R}^{d}$, satisfying the following condition: for any compact $K \subset T$ there exists a compact $C \subset S$ such that

$$
\limsup _{n \rightarrow \infty} \operatorname{Pr}^{*}\left(\left\{J_{n} \notin \mathscr{J}_{G_{n}}\left(K ; C ; \varepsilon_{n}\right)\right\}\right)=0 .
$$

Then

$$
a_{n}\left(J_{n}-G^{-1}+\operatorname{inv}\left(G^{\prime} \circ G^{-1}\right)\left(G_{n}-G\right) \circ G^{-1}\right) \rightarrow_{p, \operatorname{loc}} 0 \text { in } T \text { as } n \rightarrow \infty .
$$

Proof of Theorem 5.4. We apply Theorem A. 4 to the maps $G:=F_{P}$, $G_{n}:=F_{P_{n}}$ and $J_{n}:=Q_{n}$. By the definitions of subdifferential and subgradient, we have $F_{P_{n}}\left(Q_{n}(t)\right) \in \partial f_{P_{n}}\left(Q_{n}(t)\right)$ and, since $Q_{n}(t) \in \partial f_{P_{n}}^{-1}(t)$, we have $t \in \partial f_{P_{n}}\left(Q_{n}(t)\right)$. Thus, for all $t \in B_{P_{n}} \cap T$,

$$
\begin{align*}
\left|F_{P_{n}}\left(Q_{n}(t)\right)-t\right| & \leq \operatorname{Diam}\left(\partial f_{P_{n}}\left(Q_{n}(t)\right)\right) \leq \int_{\mathbf{x}} \operatorname{Diam}\left(\partial f\left(Q_{n}(t), x\right)\right) P_{n}(d x) \\
& =\int_{\mathbf{X}} \delta_{f}\left(Q_{n}(t) ; x\right) P_{n}(d x) \leq \sup _{s \in \mathbb{R}^{d}} \int_{\mathbf{X}} \delta_{f}(s ; x) P_{n}(d x) . \tag{A.8}
\end{align*}
$$

Let $K$ be a compact subset of $T$. Then, clearly, $K \subset B_{P}$. Take $\varepsilon>0$. Denote

$$
\begin{align*}
A:=\left\{\omega: K \subset B_{f_{P_{n}}}\right\} & \cap\left\{\omega: \sup _{s \in \mathbb{R}^{d}} \frac{\left|F_{P_{n}}(s)-F_{P}(s)\right|}{\psi(s)}<\varepsilon\right\}  \tag{A.9}\\
& \cap\left\{\omega: \sup _{s \in \mathbb{R}^{d}} \int_{\mathbf{x}} \delta_{f}(s ; x) P_{n}(d x)<\varepsilon n^{-1 / 2}\right\} .
\end{align*}
$$

It follows from Theorem 5.3 and conditions (iv), (v) that for large $n, \operatorname{Pr}^{*}(\Omega \backslash A)$ $\leq \varepsilon$.

If $\psi \equiv c$, we have for all $\omega \in A$ and $t \in K$,

$$
\begin{equation*}
\left|F_{P_{n}}\left(Q_{n}(t)\right)-F_{P}\left(Q_{n}(t)\right)\right| \leq c \varepsilon . \tag{A.10}
\end{equation*}
$$

Otherwise, by condition (ii), we get for large $a>0$,

$$
\begin{equation*}
\inf _{|s| \geq a} \frac{\left|F_{P}(s)\right|}{\psi(s)} \geq c . \tag{A.11}
\end{equation*}
$$

If $\omega \in A$, we have, by (A.9) and (A.11), that for all $s$ with $|s| \geq a$,

$$
\begin{equation*}
\left|F_{P_{n}}(s)\right| \geq(c-\varepsilon) \phi(|s|) . \tag{A.12}
\end{equation*}
$$

We have either $\left|Q_{n}(t)\right| \leq a$ or $\left|Q_{n}(t)\right|>a$ and then, by inequality (A.12),

$$
\left|F_{P_{n}}\left(Q_{n}(t)\right)\right| \geq(c-\varepsilon) \phi\left(\left|Q_{n}(t)\right|\right) .
$$

It follows from (A.8) and (A.9) that for all $\omega \in A$ and $t \in K\left|F_{P_{n}}\left(Q_{n}(t)\right)\right| \leq$ $|t|+\varepsilon n^{-1 / 2}$. Therefore, for all $\omega \in A$ and $t \in K$,

$$
\left|Q_{n}(t)\right| \leq \phi^{-1}\left(\frac{|t|+\varepsilon n^{-1 / 2}}{c-\varepsilon}\right) \bigvee a
$$

Since $K$ is bounded, there exists $r>0$ such that $|t| \leq r$ for all $t \in K$. It follows that for all $\omega \in A$,

$$
\begin{equation*}
\sup _{t \in K}\left|Q_{n}(t)\right| \leq \phi^{-1}\left(\frac{r+\varepsilon}{c-\varepsilon}\right) \bigvee a \tag{A.13}
\end{equation*}
$$

Since, by (i), $S \subset D_{P}$, it follows that $F_{P}$ is a homeomorphism of $S$ and $T$ (see the properties of $M$-distribution function). Therefore, there exists a compact
$C \subset S$ and a positive number $\delta$ such that $F_{P}^{-1}\left(K_{\delta}\right) \subset C$. Using (A.13), we get for all $\omega \in A$ and $t \in K$,

$$
\begin{equation*}
\left|F_{P_{n}}\left(Q_{n}(t)\right)-F_{P}\left(Q_{n}(t)\right)\right| \leq \varepsilon \phi\left(\left|Q_{n}(t)\right|\right) \leq \varepsilon\left[\frac{r+\varepsilon}{c-\varepsilon} \bigvee \phi(a)\right] \tag{A.14}
\end{equation*}
$$

It follows from (A.8), (A.10) and (A.14) that for $\omega \in A$ and $t \in K$,

$$
\begin{equation*}
\left|t-F_{P}\left(Q_{n}(t)\right)\right| \leq \varepsilon\left[\frac{r+\varepsilon}{c-\varepsilon} \bigvee \phi(a) \bigvee c\right]+\varepsilon \tag{A.15}
\end{equation*}
$$

Clearly, $\varepsilon>0$ can be choosen in such a way that (A.15) yields $\left|t-F_{P}\left(Q_{n}(t)\right)\right|<$ $\delta$. It means that $F_{P}\left(Q_{n}(t)\right) \in K_{\delta}$, which implies $Q_{n}(t) \in C$. Thus we proved that for all $\omega \in A, Q_{n}(K) \subset C$. By condition (iv) and (A.8), there exists a sequence $\left\{\varepsilon_{n}\right\}$ of positive numbers such that $\varepsilon_{n}=o\left(n^{-1 / 2}\right)$ and for all large $n$, $\sup _{t \in K}\left|F_{P_{n}}\left(Q_{n}(t)\right)-t\right| \leq \varepsilon_{n}$. It shows that

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \operatorname{Pr}^{*}\left(\left\{Q_{n} \notin \mathscr{J}\left(K ; C ; \varepsilon_{n}\right)\right\}\right)=0 . \tag{A.16}
\end{equation*}
$$

Under condition (iii), Theorem 5.3 implies that

$$
n^{1 / 2}\left(F_{P_{n}}-F_{P}\right) \rightarrow_{w, \text { loc }} \xi_{P} \quad \text { as } n \rightarrow \infty \text { in } \mathbb{R}^{d},
$$

and the result follows from Theorem A.4.
Proof of Theorems 5.5 and 5.6. Let $\Delta_{k}:=\left(\Phi^{-1}\left(2^{k-1}\right), \Phi^{-1}\left(2^{k}\right)\right]$ for $k \geq$ 1 , and $\Delta_{0}:=\left[0, \Phi^{-1}(1)\right]$. Define

$$
\tilde{\psi}(s):=\sum_{k=0}^{\infty} 2^{k} I_{\left\{|s| \in \Delta_{k}\right\}} \quad \text { for } s \in \mathbb{R}^{d},
$$

and let

$$
F(s, x):= \begin{cases}\frac{\Phi(|s-x|)}{\tilde{\psi}(s)} \frac{s-x}{|s-x|}, & \text { for } x \neq s \\ 0, & \text { otherwise }\end{cases}
$$

Let $\mathscr{F}:=\bigcup_{1}^{d} \mathscr{F}^{(j)}$, where $\mathscr{F}^{(j)}:=\left\{F^{(j)}(s, \cdot): s \in \mathbb{R}^{d}\right\}$. Since $\tilde{\psi}(s) \geq 1$, condition (5.1) implies that the function $E_{\mathscr{F}}(x):=C_{1}+C_{2} \Phi(|x|)$ is an envelope of the class $\mathscr{F}$ for some constants $C_{1}>0, C_{2}>0$. Consider the functions

$$
F_{1}(s, x):= \begin{cases}\frac{\Phi(|s-x|)}{\tilde{\psi}(s)}, & \text { for } x \neq s \\ 0, & \text { otherwise }\end{cases}
$$

and

$$
F_{2}(s, x):= \begin{cases}E_{\mathscr{F}}(x) \frac{s-x}{|s-x|}, & \text { for } x \neq s \\ 0, & \text { otherwise }\end{cases}
$$

Let

$$
\begin{aligned}
\mathscr{T}_{1} & :=\left\{F_{1}(s, \cdot): s \in \mathbb{R}^{d}\right\}, \\
\mathscr{F}_{1, k}: & :=\left\{F_{1}(s, \cdot): s \in \Delta_{k}\right\}=\left\{2^{-k} \Phi(|s-\cdot|): s \in \Delta_{k}\right\}, \\
\mathscr{F}_{2}^{(j)} & :=\left\{F_{2}^{(j)}(s, \cdot): s \in \mathbb{R}^{d}\right\}, \quad 1 \leq j \leq d .
\end{aligned}
$$

Clearly, $\mathscr{F}_{1}=\bigcup_{0}^{\infty} \mathscr{F}_{1, k}$. It's also obvious that $E_{\mathscr{F}}$ is an envelope of $\mathscr{F}_{1}$ and $\mathscr{F}_{2}^{(j)}, 1 \leq j \leq d$. The inequality

$$
\begin{aligned}
&\left|F^{(j)}\left(s_{1}, x\right)-F^{(j)}\left(s_{2}, x\right)\right| \\
&= \left\lvert\,\left(\frac{\Phi\left(\left|s_{1}-x\right|\right)}{\tilde{\psi}\left(s_{1}\right)}-\frac{\Phi\left(\left|s_{2}-x\right|\right)}{\tilde{\psi}\left(s_{2}\right)}\right) \frac{s_{1}^{(j)}-x^{(j)}}{\left|s_{1}-x\right|}\right. \\
& \left.+\frac{\Phi\left(\left|s_{2}-x\right|\right)}{\tilde{\psi}\left(s_{2}\right)}\left(\frac{s_{1}^{(j)}-x^{(j)}}{\left|s_{1}-x\right|}-\frac{s_{2}^{(j)}-x^{(j)}}{\left|s_{2}-x\right|}\right) \right\rvert\, \\
& \leq\left|\frac{\Phi\left(\left|s_{1}-x\right|\right)}{\tilde{\psi}\left(s_{1}\right)}-\frac{\Phi\left(\left|s_{2}-x\right|\right)}{\tilde{\psi}\left(s_{2}\right)}\right|+E_{\mathscr{F}}(x)\left|\frac{s_{1}^{(j)}-x^{(j)}}{\left|s_{1}-x\right|}-\frac{s_{2}^{(j)}-x^{(j)}}{\left|s_{2}-x\right|}\right| \\
& \leq\left|F_{1}\left(s_{1}, x\right)-F_{1}\left(s_{2}, x\right)\right|+\left|F_{2}^{(j)}\left(s_{1}, x\right)-F_{2}^{(j)}\left(s_{2}, x\right)\right|
\end{aligned}
$$

implies that for any $1 \leq r<+\infty$ and $Q \in \mathscr{P}\left(\mathbb{R}^{d}\right)$ with $\int_{\mathbb{R}^{d}} E_{\mathscr{F}}^{r}(x) Q(d x)<+\infty$, the following estimate holds for the metrics $d_{Q, r}$ :

$$
\begin{aligned}
& d_{Q, r}\left(F^{(j)}\left(s_{1}, \cdot\right), F^{(j)}\left(s_{2}, \cdot\right)\right) \\
& \quad \leq d_{Q, r}\left(F_{1}\left(s_{1}, \cdot\right), F_{1}\left(s_{2}, \cdot\right)\right)+d_{Q, r}\left(F_{2}^{(j)}\left(s_{1}, \cdot\right), F_{2}^{(j)}\left(s_{2}, \cdot\right)\right) .
\end{aligned}
$$

It follows that

$$
H_{d_{Q, r}}\left(\mathscr{F}^{(j)} ; \varepsilon\right) \leq H_{d_{Q, r}}\left(\mathscr{F}_{1} ; \varepsilon / 2\right)+H_{d_{Q, r}}\left(\mathscr{F}_{2}^{(j)} ; \varepsilon / 2\right)
$$

and

$$
\begin{equation*}
H_{r}\left(\mathscr{F}^{(j)} ; \varepsilon\right) \leq H_{r}\left(\mathscr{F}_{1} ; \varepsilon / 2\right)+H_{r}\left(\mathscr{F}_{2}^{(j)} ; \varepsilon / 2\right) . \tag{A.17}
\end{equation*}
$$

In order to estimate the entropies, note that for all $k \geq 0 \mathscr{F}_{1, k}$ is a VC-subgraph class with $V\left(\mathscr{F}_{1, k}\right)$ depending only on $d$. The proof follows the lines of Pollard (1984), Chapter 2, Example 26. Indeed, for all $s \in \Delta_{k}$,

$$
\begin{aligned}
\left\{(x, t): 2^{-k} \Phi(|s-x|) \geq t\right\} & =\left\{(x, t):|s-x|^{2} \geq\left(\Phi^{-1}\left(2^{k} t\right)\right)^{2}\right\} \\
& =\left\{(x, t): \sum_{1}^{d}\left(s^{(j)}-x^{(j)}\right)^{2} \geq\left(\Phi^{-1}\left(2^{k} t\right)\right)^{2}\right\} \\
& =\left\{(x, t): P(x)-\left(\Phi^{-1}\left(2^{k} t\right)\right)^{2} \geq 0\right\},
\end{aligned}
$$

where $P(x)$ is a polynomial of $x=\left(x^{(1)}, \ldots, x^{(d)}\right)$ of the second degree. Denote by $\mathscr{H}$ the linear span of the set of all such polynomials and a function $\left(\Phi^{-1}\left(2^{k} t\right)\right)^{2}$. Then for any $s \in \Delta_{k}$ the set $\left\{(x, t): 2^{-k} \Phi(|s-x|) \geq t\right\}$ can be represented as $\{(x, t): h(x, t) \geq 0\}$ for some $h \in \mathscr{H}$. Thus, $\mathscr{F}_{1, k}$ is a VCsubgraph class and $V\left(\mathscr{F}_{1}, k\right)$ depends only on $d$. It follows that for $1 \leq r<+\infty$
$N_{r}\left(\mathscr{T}_{1}, k ; \varepsilon\right) \leq C \varepsilon^{-V}, \varepsilon>0$ with constants $C>0, V>0$ depending only on $d$ and $r$. Since the function $2^{-k} E_{\mathscr{F}}(\cdot)$ is an envelope of $\mathscr{F}_{1, k}$, we have for all $Q \in \mathscr{P}\left(\mathbb{R}^{d}\right)$,

$$
N_{d_{Q, r}}\left(\mathscr{F}_{1, k} ; \varepsilon 2^{-k}\left(\int_{\mathbf{X}} E_{\mathscr{F}}^{r} d Q\right)^{1 / r}\right) \leq C \varepsilon^{-V}, \quad \varepsilon>0
$$

or

$$
N_{d_{Q, r}}\left(\mathscr{F}_{1, k} ; \varepsilon\left(\int_{\mathbf{X}} E_{\mathscr{F}}^{r} d Q\right)^{1 / r}\right) \leq C 2^{-k V} \varepsilon^{-V}, \quad \varepsilon>0 .
$$

Hence, for all $Q \in \mathscr{P}\left(\mathbb{R}^{d}\right)$ and for all $\varepsilon>0$,

$$
\begin{aligned}
N_{d_{Q, r}}\left(\mathscr{F}_{1} ; \varepsilon\left(\int_{\mathbf{X}} E_{\mathscr{F}}^{r} d Q\right)^{1 / r}\right) & \leq \sum_{k=0}^{\infty} N_{d_{Q, r}}\left(\mathscr{F}_{1, k} ; \varepsilon\left(\int_{\mathbf{X}} E_{\mathscr{F}}^{r} d Q\right)^{1 / r}\right) \\
& \leq C \sum_{1}^{\infty} 2^{-k V} \varepsilon^{-V} .
\end{aligned}
$$

Therefore,

$$
\begin{equation*}
H_{r}\left(\mathscr{F}_{1} ; \varepsilon\right)=O\left(\log \frac{1}{\varepsilon}\right) \quad \text { as } \varepsilon \rightarrow 0 \tag{A.18}
\end{equation*}
$$

Now we show that all classes $\mathscr{F}_{2}^{(j)}, 1 \leq j \leq d$ are $V C$-subgraph, too. We have for any $s \in \mathbb{R}^{d}$,

$$
\begin{align*}
&\{(x, t)\left.: F_{2}^{(j)}(s, x) \geq t \geq 0\right\} \\
&=\left\{(x, t): s^{(j)}-x^{(j)} \geq \frac{t|s-x|}{E_{\mathscr{F}}(x)}, x \neq s, t \geq 0\right\} \cup\{(s, 0)\} \\
&=\left\{(x, t):\left(s^{(j)}-x^{(j)}\right)^{2}-\frac{t^{2} \sum_{k=1}^{d}\left(s^{(k)}-x^{(k)}\right)^{2}}{E_{\mathscr{F}}^{2}(x)} \geq 0,\right.  \tag{A.19}\\
&x \neq s, t \geq 0\} \cup\{(s, 0)\} .
\end{align*}
$$

Let $\mathscr{H}$ be the linear span of functions $1, x^{(j)},\left(x^{(j)}\right)^{2}, t^{2}, t^{2}\left(x^{(j)} / E_{\mathscr{F}}^{2}(x)\right)$ and $t^{2}\left(x^{(j)} / E^{\mathscr{F}}(x)\right)^{2}$. Then the class of sets $\measuredangle:=\{\{(x, t): h(x, t) \geq 0\}: h \in \mathscr{H}\}$ is Vapnik-Chervonenkis, and so is the class

$$
\left.\tilde{\mathscr{C}}:=\{C \cap\{(x, t): x \neq s, t \geq 0\}\} \cup\{(s, 0)\}: C \in \mathscr{C}, s \in \mathbb{R}^{d}\right\} .
$$

By (A.19), for all $s \in \mathbb{R}^{d},\left\{(x, t): F_{2}^{(j)}(s, x) \geq t \geq 0\right\} \in \tilde{\boldsymbol{C}}$. Therefore $\{\{(x, t)$ : $\left.\left.F_{2}^{(j)}(s, x) \geq t \geq 0\right\}: s \in \mathbb{R}^{d}\right\}$ is a Vapnik-Chervonenkis class. Quite similarly, the class $\left\{\left\{(x, t): F_{2}^{(j)}(s, x) \leq t \leq 0\right\}: s \in \mathbb{R}^{d}\right\}$ is Vapnik-Chervonenkis, too. It follows that for all $1 \leq j \leq d$, classes of functions $\mathscr{F}_{2}^{(j)}$ are VC-subgraph, and we have

$$
\begin{equation*}
H_{r}\left(\mathscr{F}_{2}^{(j)} ; \varepsilon\right)=O\left(\log \frac{1}{\varepsilon}\right) \quad \text { as } \varepsilon \rightarrow 0 \tag{A.20}
\end{equation*}
$$

It follows from (A.17), (A.18) and (A.20) that for all $1 \leq j \leq d$,

$$
H_{r}\left(\mathscr{T}^{(j)} ; \varepsilon\right)=O\left(\log \frac{1}{\varepsilon}\right) \quad \text { as } \varepsilon \rightarrow 0
$$

which implies

$$
\begin{equation*}
H_{r}(\mathscr{F} ; \varepsilon)=O\left(\log \frac{1}{\varepsilon}\right) \quad \text { as } \varepsilon \rightarrow 0 \tag{A.21}
\end{equation*}
$$

Condition (5.3) implies that $\int_{\mathbb{R}^{d}} E_{\mathscr{T}}(x) P(d x)<+\infty$, which along with (A.21) implies that $\mathscr{F} \in G C(P)$, and (5.4) follows. It, in turn, implies (5.5) and (5.6).

Under condition (5.7) we have $\int_{\mathbb{R}^{d}} E_{\mathscr{F}}^{2}(x) P(d x)<+\infty$, which implies [again, along with (A.21)] that $\mathscr{F} \in \operatorname{CLT}(P)$. It follows that the process

$$
n^{1 / 2} \frac{F_{P_{n}}(s)-F_{P}(s)}{\tilde{\psi}(s)}=n^{1 / 2}\left(P_{n}-P\right)(F(s, \cdot)), \quad s \in \mathbb{R}^{d}
$$

converges weakly in $l^{\infty}\left(\mathbb{R}^{d}\right)$ to

$$
\frac{\xi_{P}(s)}{\tilde{\psi}(s)}=\int_{\mathbb{R}^{d}} F(s, x) W_{P}^{\circ}(d x), s \in \mathbb{R}^{d} .
$$

By the definition of the function $\tilde{\psi}$, we have for all $s \in \mathbb{R}^{d} 1 / 2 \leq(\psi(s) / \tilde{\psi}(s)) \leq$ 1. Therefore, (5.8) also holds. The kernel

$$
G(s, x):=\frac{\Phi(|s-x|)}{\psi(s)} \frac{s-x}{|s-x|} \mathbf{I}_{\{s \neq x\}}
$$

is continuous in $s$ at any point $s_{0} \neq x$. Since $P$ is nonatomic, it is, in fact, continuous in $s \in \mathbb{R}^{d}$ for $P$-almost all $x$. This, by dominated convergence, implies continuity of $G(s, \cdot)$ in $s \in \mathbb{R}^{d}$ with respect to the metric $\rho_{P}$. Since $\xi_{P}(s)=\psi(s) W_{P}^{\circ}(G(s, \cdot)), W_{P}^{\circ}$ is a.s. $\rho_{P}$-continuous on $\left\{G(s ; \cdot): s \in \mathbb{R}^{d}\right\}$, and $\psi$ is continuous in $\mathbb{R}^{d}$, we can conclude that the process $\xi_{P}$ is continuous on $\mathbb{R}^{d}$. This fact and (5.8) imply (5.9) and (5.10).

Proof of Theorem 5.7. We use Theorem 5.4. Condition (i) obviously holds. Condition (ii) with $\psi(s):=\Phi(|s|) \vee 1$ is a consequence of the following fact (which is rather easy to prove): under conditions (5.1) and (5.2), $F_{P}(s)=\Phi(|s|)(s /|s|)+o(\Phi(|s|))$ as $s \rightarrow \infty$. Condition (iii) follows from Theorem 5.6. Next, since

$$
\delta_{f}(s ; x):=\operatorname{Diam}(\partial f(s ; x))= \begin{cases}\rho_{\varphi}, & \text { for } s=x \\ 0, & \text { otherwise }\end{cases}
$$

we have

$$
\sup _{s \in \mathbb{R}^{d}} \int_{\mathbb{R}^{d}} \delta_{f}(s ; x) P_{n}(d x)=\frac{\rho_{\varphi}}{n}=o\left(n^{-1 / 2}\right) .
$$

So, condition (iv) also holds. Finally, condition (v) is obvious, since $B_{P_{n}}=B_{P}=$ $B\left(0 ; R_{\varphi}\right)$, and the result follows from Theorem 5.4.

Proof of Theorem 4.5. We apply Theorem 5.4 to $F:=G$, and $R, R_{n}$ instead of $P, P_{n}$. First note that $G_{R}$ is continuously differentiable in $\mathbb{R}^{m d}$ with nonsingular derivative $G_{R}^{\prime}(S), S \in \mathbb{R}^{m d}$ given by (4.4). This means that condition (i) of Theorem 5.4 holds. Assume that $\psi \equiv 1$, so that (ii) holds. Note that $G$ is $R$-Donsker, since the class of functions $\bigcup_{j=1}^{m+d}\left\{G^{(j)}(S ; \cdot):, S \in \mathbb{R}^{m d}\right\}$ is a Vapnik-Chervonenkis subgraph with a square integrable envelope, and for all $1 \leq j \leq d$, the functions $G^{(j)}(S ; \cdot, \cdot)$ are $\rho_{R^{-}}$-continuous in $S \in \mathbb{R}^{m d}$. Thus (iii) holds. Next we have

$$
\delta_{g}(S ; x, y)= \begin{cases}|x|, & \text { if } S x=y \\ 0, & \text { otherwise }\end{cases}
$$

Therefore,

$$
\sup _{S \in \mathbb{R}^{n d}} \int_{\mathbb{R}^{m+d}} \delta_{g}(S ; x, y) R_{n}(d x ; d y)=\sup _{S \in \mathbb{R}^{m d}} \sum_{i=1}^{n}\left|X_{i}\right| \mathbf{I}_{\left\{S X_{i}=Y_{i}\right\}}
$$

Under conditions (i), (ii), we have $\mathbb{E} \mathbf{I}_{\{S X=Y\}}|X|=0$, and rather standard empirical processes argument shows that

$$
\sup _{S \in \mathbb{R}^{n d}} \sum_{i=1}^{n}\left|X_{i}\right| \mathbf{I}_{\left\{S X_{i}=Y_{i}\right\}}=o_{p}\left(n^{-1 / 2}\right),
$$

so (iv) holds. By the law of large numbers in separable Banach spaces, the sequence of stochastic processes $n^{-1}\left(\left|V X_{1}\right|+\cdots+\left|V X_{n}\right|\right)$ converges to the function $\mathbb{E}|V X|$ uniformly in $V \in S^{d m-1}$. This implies condition (v).

Thus, we can use Theorem 5.4 to get the following asymptotic representation of any $J_{n}: \Omega \times B_{R} \mapsto \mathbb{R}^{m d}$, such that $J_{n}(T) \in \partial g_{R_{n}}^{-1}(T)$ for all $T \in B_{R}$ :

$$
J_{n}=G_{R}^{-1}-\operatorname{inv}\left(G_{R}^{\prime} \circ G_{R}^{-1}\right)\left(G_{R_{n}}-G_{R}\right) \circ G_{R}^{-1}+o_{p}\left(n^{-1 / 2}\right),
$$

which holds locally uniformly in $B_{R}$.
Take $r<1$. Using consistency of $\bar{X}_{n}$, it is easy to show that there exists a compact $K \subset B_{R}$ such that with probability close to $1, t \otimes \bar{X}_{n} \in K$ for all $t$ with $|t| \leq r$ and for all $n$ large enough. It follows that uniformly in $|t| \leq r$,

$$
\begin{gathered}
J_{n}\left(t \otimes \bar{X}_{n}\right)=G_{R}^{-1}\left(t \otimes \bar{X}_{n}\right)-\operatorname{inv}\left(G_{R}^{\prime} \circ G_{R}^{-1}\left(t \otimes \bar{X}_{n}\right)\right)\left(G_{R_{n}}-G_{R}\right) \\
\circ G_{R}^{-1}\left(t \otimes \bar{X}_{n}\right)+o_{p}\left(n^{-1 / 2}\right) .
\end{gathered}
$$

Since $G_{R}^{-1}$ is uniformly continuous in $K, n^{1 / 2}\left(G_{R_{n}}-G_{R}\right)$ is asymptotically equicontinuous and $\bar{X}_{n} \rightarrow \mathbb{E} X$ as $n \rightarrow \infty$, we have uniformly in $|t| \leq r$,

$$
\begin{aligned}
J_{n}\left(t \otimes \bar{X}_{n}\right)= & {\left[G_{R}^{-1}\left(t \otimes \bar{X}_{n}\right)-G_{R}^{-1}(t \otimes \mathbb{E} X)\right]+G_{R}^{-1}(t \otimes \mathbb{E} X) } \\
& -\operatorname{inv}\left(G_{R}^{\prime} \circ G_{R}^{-1}(t \otimes \mathbb{E} X)\right)\left(G_{R_{n}}-G_{R}\right) \circ G_{R}^{-1}(t \otimes \mathbb{E} X)+o_{p}\left(n^{-1 / 2}\right)
\end{aligned}
$$

Note that there exists $J_{n}: \Omega \times B_{R} \mapsto \mathbb{R}^{m d}$, such that $J_{n}(T) \in \partial g_{R_{n}}^{-1}(T)$ for all $T \in B_{R}$ and $Q_{n}(t)=J_{n}\left(t \otimes \bar{X}_{n}\right)$ for all $t \in B_{d}$. Using differentiability of $G_{R}^{-1}$ and the condition $\left|\bar{X}_{n}-\mathbb{E} X\right|=O_{P}\left(n^{-1 / 2}\right)$, we now easily get asymptotic representation (4.5), and complete the proof.

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