# THE BAHADUR-KIEFER REPRESENTATION FOR $U$-QUANTILES ${ }^{1}$ 

By Miguel A. Arcones

University of Texas


#### Abstract

We consider the distributional and the almost sure pointwise Bahadur-Kiefer representation for $U$-quantiles. We show that the order of this representation depends on the order of the local variance of the empirical process of $U$-statistic structure at the $U$-quantile. Our results indicate that $U$-quantiles can be smoother than quantiles. $U$-quantiles can either be as unsmooth as quantiles or can behave as differentiable statistical functionals.


1. Introduction. First, let us recall the Bahadur-Kiefer representation for quantiles. Let $\left\{X_{i}\right\}_{i=1}^{\text {º }^{\infty}}$ be a sequence of i.i.d. r.v.'s, let $F_{n}$ be the empirical distribution function, let $F$ be the cumulative distribution function of $X_{i}$ and let $0<p<1$. Define $\xi_{n}:=\inf \left\{t: F_{n}(t) \geq p\right\}$ and

$$
\begin{equation*}
R_{n}:=\xi_{n}-\xi_{0}+\left(F^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(F_{n}\left(\xi_{0}\right)-F\left(\xi_{0}\right)\right), \tag{1.1}
\end{equation*}
$$

where $F\left(\xi_{0}\right)=p$. Kiefer (1967) showed that if $F$ is second differentiable at $\xi_{0}$ and $F^{\prime}\left(\xi_{0}\right)>0$, then

$$
\begin{equation*}
n^{3 / 4} R_{n} \rightarrow_{d} p^{1 / 4}(1-p)^{1 / 4}\left|g_{1}\right|^{1 / 2} g_{2}, \tag{1.2}
\end{equation*}
$$

where $g_{1}$ and $g_{2}$ are two independent standard normal r.v.'s. He also proved that
(1.3) $\quad \limsup _{n \rightarrow \infty} \pm(n / 2 \log \log n)^{3 / 4} R_{n}=2^{1 / 2} 3^{-3 / 4} p^{1 / 4}(1-p)^{1 / 4} \quad$ a.s.

The purpose of this paper is to present similar results for $U$-quantiles. Let $\left\{X_{i}\right\}_{i=1}^{\infty}$ be a sequence of i.i.d. r.v.'s with values in a measurable space $(S, \mathscr{S})$. Let $h: S^{m} \rightarrow \mathbb{R}$ be a measurable symmetric function. Let $H(t)=$ $\operatorname{Pr}\left\{h\left(X_{1}, \ldots, X_{m}\right) \leq t\right\}$. The empirical distribution of $U$-statistic structure is defined by

$$
\begin{equation*}
H_{n}(t):=\frac{(n-m)!}{n!} \sum_{\left(i_{1}, \ldots, i_{m}\right) \in I_{m}^{n}} I_{h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right) \leq t}, \tag{1.4}
\end{equation*}
$$

[^0]where $I_{m}^{n}=\left\{\left(i_{1}, \ldots, i_{m}\right): 1 \leq i_{j} \leq n\right.$ and $i_{k} \neq i_{j}$ for $\left.k \neq j\right\}$. Let $0<p<1$. Suppose that $H\left(\xi_{0}\right)=p$. The $U$-quantile is defined by
\[

$$
\begin{equation*}
\xi_{n}:=\inf \left\{t: H_{n}(t) \geq p\right\} . \tag{1.5}
\end{equation*}
$$

\]

Several common estimators are $U$-quantiles. For example, one very often used alternative to the median as a center of symmetry is the HodgesLehmann estimator: the median of $2^{-1}\left(X_{i}+X_{j}\right), 1 \leq i<j \leq n$ [see Hodges and Lehmann (1963)]. This is the $U$-quantile (with respect to $p=1 / 2$ ) of the kernel $h\left(x_{1}, x_{2}\right)=2^{-1}\left(x_{1}+x_{2}\right)$. We refer to Lehmann (1975) for different extensions of this estimator and applications to nonparametric statistics. Another interesting example is the $U$-quantile of the kernel $h\left(x_{1}, x_{2}\right)=\mid x_{1}-$ $x_{2} \mid$ with respect to $p=1 / 2$. This $U$-quantile is a measure of the spread of the distribution. It was introduced by Bickel and Lehmann (1979). Choudhury and Serfling (1988) introduced an $U$-quantile which estimates the regression slope. Consider the linear regression model: $Y_{i}=\alpha+\beta X_{i}+\delta_{i} ; \alpha$ and $\beta$ are constants and $\delta_{i}$ is an r.v. independent of $X_{i}$. The $U$-quantile of the kernel $h\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right)\right)=\left(y_{2}-y_{1}\right) /\left(x_{2}-x_{1}\right)$, with respect to $p=1 / 2$, is a natural estimator of the parameter $\beta$. This estimator is the median of the values $\left(Y_{j}-Y_{i}\right) /\left(X_{j}-X_{i}\right), 1 \leq i<j \leq n$. Some references in the study of $U$-quantiles are Serfling (1984), Janssen, Serfling and Veraverbeke (1984), Helmers, Janssen and Serfling (1988) and Choudhury and Serfling (1988).

Here, we will study the distributional and the a.s. behavior of

$$
\begin{equation*}
R_{n}:=\xi_{n}-\xi_{0}+\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)\right) \tag{1.6}
\end{equation*}
$$

using empirical process techniques. Finding the asymptotic behavior of (1.6), we grasp a very good insight into the effect of the influence curve in the asymptotics of $U$-quantiles. The Bahadur-Kiefer representation of a statistical functional measures how close is, asymptotically, the linear expansion of a statistical functional to the statistical functional itself. It is a way to measure the differentiability of the statistical functional. We refer to Serfling (1980), Chapter 6, and Dudley $(1992,1994)$ for other ways to measure differentiability. One interesting application of Bahadur-Kiefer representations is to obtain sequential fixed-width confidence intervals for a parameter [see Chow and Robbins (1965) and Geertsema (1970)].

The leading idea to deal with (1.6) is to do a Hoeffding decomposition, to show that the terms of order 2 and larger vanish and to find the order from the first term of this decomposition. Next, we describe the Hoeffding decomposition. Given a measurable function on $S^{m}$, the $U$-statistic with kernel $h$ is defined by

$$
\begin{equation*}
U_{n}(h):=\frac{(n-m)!}{n!} \sum_{\left(i_{1}, \ldots, i_{m}\right) \in I_{m}^{n}} h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right) . \tag{1.7}
\end{equation*}
$$

We will abbreviate $E h:=E\left[h\left(X_{1}, \ldots, X_{m}\right)\right], P_{n} f=n^{-1} \sum_{j=1}^{n} f\left(X_{j}\right)$ and $P f=$ $E[f(X)]$, where $X$ is a copy of $X_{1}$. We define

$$
\begin{equation*}
\pi_{k, m} f\left(x_{1}, \ldots, x_{k}\right)=\left(\delta_{x_{1}}-P\right) \cdots\left(\delta_{x_{k}}-P\right) P^{m-k} f \tag{1.8}
\end{equation*}
$$

where $Q_{1}, \ldots, Q_{m} f=\int \cdots \int f\left(x_{1}, \ldots, x_{m}\right) d Q_{1}(x) \cdots d Q_{m}\left(x_{m}\right)$. Then, the Hoeffding decomposition of the $U$-statistic $U_{n}(f)$ can be written as

$$
\begin{equation*}
U_{n}(f)=\sum_{k=0}^{m}\binom{m}{k} U_{n}\left(\pi_{k, m} f\right) \tag{1.9}
\end{equation*}
$$

In particular, this expansion applies to the term $\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)\right)$ in (1.6), allowing us to see how close $\xi_{n}-\xi_{0}$ is to a true linear term. Here, we will see that the order of

$$
\begin{equation*}
E\left[\left|g(X, t)-g\left(X, \xi_{0}\right)\right|^{2}\right] \quad \text { as } t \rightarrow \xi_{0} \tag{1.10}
\end{equation*}
$$

determines the order of $R_{n}$, where $g(x, t)=\operatorname{Pr}\left\{h\left(x, X_{2}, \ldots, X_{m}\right) \leq t\right\}$. In particular, we will see that if $E\left[\left|g(X, t)-g\left(X, \xi_{0}\right)\right|^{2}\right]=O\left(\left|t-\xi_{0}\right|^{v}\right)$ as $t \rightarrow$ $\xi_{0}$, for some $v>0$, then

$$
\begin{equation*}
n^{(v+2) / 4} R_{n}=O_{\mathrm{P}}(1) \tag{1.11}
\end{equation*}
$$

and

$$
\begin{equation*}
(n / \log \log n)^{(v+2) / 4} R_{n}=O(1) \quad \text { a.s. } \tag{1.12}
\end{equation*}
$$

Observe that by the Cauchy-Schwarz inequality

$$
\begin{aligned}
\left|H(t)-H\left(\xi_{0}\right)\right|^{2} & =\mid E\left[I_{h\left(X_{1}, \ldots, X_{m}\right) \leq t}-I_{h\left(X_{1}, \ldots, X_{m}\right) \leq \xi_{0}}| |^{2}\right. \\
& \leq E_{1}\left[\left|E_{2, \ldots, m}\left[I_{h\left(X_{1}, \ldots, X_{m}\right) \leq t}-I_{h\left(X_{1}, \ldots, X_{m}\right) \leq \xi_{0}}\right]\right|^{2}\right] \\
& =E\left[\left|g(X, t)-g\left(X, \xi_{0}\right)\right|^{2}\right] \\
& \leq E\left[\left|I_{h\left(X_{1}, \ldots, X_{m}\right) \leq t}-I_{h\left(X_{1}, \ldots, X_{m}\right) \leq \xi_{0}}\right|^{2}\right]=\left|H(t)-H\left(\xi_{0}\right)\right|
\end{aligned}
$$

where by $E_{i_{1}, \ldots, i_{k}}$ we mean integration with respect to the coordinates $i_{1}, \ldots, i_{k}$. So, if $H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and $E[\mid g(X, t)-$ $\left.\left.g\left(X, \xi_{0}\right)\right|^{2}\right]=O\left(\left|t-\xi_{0}\right|^{v}\right)$ as $t \rightarrow \xi_{0}$, for some $v>0$, then $1 \leq v \leq 2$. Finding the exact order of (1.10) may be difficult or impossible, but, by (1.13), (1.11), and (1.12) always hold with $v=1$. For a smooth statistical functional, the term $R_{n}$ is $O_{\mathrm{P}}\left(n^{-1}\right)$ and $O\left(n^{-1}(\log \log n)\right)$ a.s. (case $v=2$ ). These are the orders of all the examples mentioned above. These estimators enjoy a much better differentiability than the median.

We must mention the previous work in this problem. Choudhury and Serfling (1988) showed that, under some mild conditions,

$$
\begin{equation*}
n^{3 / 4}(\log n)^{-3 / 4}\left(H_{n}\left(\xi_{0}\right)-p+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)=O(1) \quad \text { a.s. } \tag{1.14}
\end{equation*}
$$

[see also Lemma 4.2 in Geertsema (1970)]. This result was used in Gijbels, Janssen and Veraverbeke (1988) to find weak and strong representations for trimmed $U$-statistics. We also must mention the work by Shi (1995) in the uniform Bahadur-Kiefer representation for the $U$-quantiles of $h\left(x_{1}, \ldots\right.$, $\left.x_{m}\right)=\max \left(x_{1}, \ldots, x_{m}\right)$. Other papers related to the present one are the ones by Carroll (1978), Jurečková (1980), Jurečková and Sen (1987), Deheuvels
and Mason (1992) and Arcones (1994a) in the Bahadur-Kiefer representation of $M$-estimators.

Our main tools are certain limit theorems which hold uniformly over VC subgraph classes of functions. Given a set $S$ and a collection of subsets $\mathscr{C}$, for $A \subset S$, let $\Delta^{\mathscr{E}}(A)=\operatorname{card}\{A \cap C: C \in \mathscr{C}\}$, let $m^{\mathscr{E}}(n)=\max \left\{\Delta^{\mathscr{C}}(A): \operatorname{card}(A)=\right.$ $n\}$ and let $s(\mathscr{C})=\inf \left\{n: m^{\mathscr{C}}(n)<2^{n}\right\} ; \mathscr{C}$ is said to be a VC class of sets if $s(\mathscr{C})<\infty$. General properties of VC classes of sets can be found in Chapters 9 and 11 in Dudley (1984). Given a function $f: S \rightarrow \mathbb{R}$, the subgraph of $f$ is the set $\{(x, t) \in S \times \mathbb{R}: 0 \leq t \leq f(x)$ or $f(x) \leq t \leq 0\}$. A class of functions $\mathscr{F}$ is a VC subgraph class if the collection of subgraphs of $\mathscr{F}$ is a VC class. The interest of these classes of functions lies in their good properties with respect to covering numbers. Given a pseudometric space $(T, d)$ the $\varepsilon$-covering number $N(\varepsilon, T, d)$ is defined by

$$
\begin{equation*}
N(\varepsilon, T, d)=\min \{m: \text { there exists a covering of } T \tag{1.15}
\end{equation*}
$$

by $m$ balls of radius $\leq \varepsilon\}$.
Given a positive measure $\mu$ on $(S, \mathscr{S})$, we define $N_{2}(\varepsilon, \mathscr{F}, \mu)=N(\varepsilon, \mathscr{F}$, $\left.\|\cdot\|_{L_{2}(\mu)}\right)$. If $\mathscr{F}$ is a VC subgraph class [Pollard (1984), Proposition 2.25], there are finite constants $A$ and $v$ such that, for each probability measure $\mu$ with $\mu F^{2}<\infty$,

$$
\begin{equation*}
N_{2}(\varepsilon, \mathscr{F}, \mu) \leq A\left(\left(\mu F^{2}\right)^{1 / 2} / \varepsilon\right)^{v} \tag{1.16}
\end{equation*}
$$

where $F(x)=\sup _{f \in \mathscr{F}}|f(x)|$ and $A$ and $v$ can be chosen depending only on $s(\mathscr{F})$, that is, uniformly over all the classes of functions with the same number $s(\mathscr{F})$. By the maximal inequality for sub-Gaussian processes [see Theorem 2.3.1 in Marcus and Pisier (1981); see also Theorem 1 in Dudley (1967)], there is a constant $c$ depending only on $A$ and $v$ such that for any class of functions satisfying (1.16),

$$
\begin{equation*}
n^{-1} E\left[\sup _{f \in \mathscr{F}}\left|\sum_{i=1}^{n} \varepsilon_{i} f\left(X_{i}\right)\right|^{2}\right] \leq c E\left[F^{2}(X)\right] \tag{1.17}
\end{equation*}
$$

where $\left\{\varepsilon_{i}\right\}_{i=1}^{\infty}$ is a Rademacher sequence independent of the sequence $\left\{X_{i}\right\}_{i=1}^{\infty}$.
One of the main ingredients to study the distributional Bahadur-Kiefer representation of $U$-quantiles will be the weak convergence of a sequence of stochastic processes. By weak convergence, we mean weak convergence of random elements with values in $l_{\infty}(\mathscr{F})$ as in Hoffmann-Jørgensen (1984). $l_{\infty}(\mathscr{F})$ is the Banach space formed by all the uniformly bounded functions on $\mathscr{F}$ with the norm $\|x\|_{\mathscr{F}}=\sup _{f \in \mathscr{F}}|x(f)|$. Let $\left\{Z_{n}(f): f \in \mathscr{F}\right\}, n \geq 1$, be a sequence of stochastic processes, and let $\{Z(f): f \in \mathscr{F}\}$ be another stochastic process. The sequence of stochastic processes $\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ is said to converge weakly to $\{Z(f): f \in \mathscr{F}\}$ in $l_{\infty}(\mathscr{F})$ if:
(i) $\sup _{f \in \mathscr{F}}\left|Z_{n}(f)\right|<\infty$ a.s. for each $n$ large enough;
(ii) there exists a separable set $S$ of $l_{\infty}(\mathscr{F})$ such that $\operatorname{Pr}^{*}\{Z \in S\}=1$;
(iii) $E^{*}\left[H\left(Z_{n}\right)\right] \rightarrow E[H(Z)]$ for each bounded, continuous function $H$ in $l_{\infty}(\mathscr{F})$.

It is well known [see, e.g., Andersen and Dobric (1987)] that this type of convergence is equivalent to the convergence of the finite-dimensional distributions plus a finite-dimensional approximation, that is, $\left\{Z_{n}(f): f \in \mathscr{F}\right\}$, $n \geq 1$, converges weakly to $\{Z(f): f \in \mathscr{F}\}$ if and only if the finite-dimensional distributions of $\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ converge to those of $\{Z(f): f \in \mathscr{F}\}$; and for each $\eta>0$, there exists a map $\pi: \mathscr{F} \rightarrow \mathscr{F}$ such that $\#\{\pi f: f \in \mathscr{F}\}$ is finite and

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \operatorname{Pr}^{*}\left\{\sup _{f \in \mathscr{F}}\left|Z_{n}(f)-Z_{n}(\pi f)\right| \geq \eta\right\} \leq \eta \tag{1.18}
\end{equation*}
$$

We also will use the fact that if $\{Z(f): f \in \mathscr{F}\}$ is a stochastic process such that there exists a separable set $S$ of $l_{\infty}(\mathscr{F})$ with $\operatorname{Pr}^{*}\{Z \in S\}=1$, then $\left(\mathscr{F}, \rho_{1}\right)$ is totally bounded and $\operatorname{Pr}^{*}\left\{Z \in C_{u}\left(\mathscr{F}, \rho_{1}\right)\right\}=1$, where $C_{u}\left(\mathscr{F}, \rho_{1}\right)$ is the set of all uniformly bounded and $\rho_{1}$-uniformly continuous functions in $\mathscr{F}$ and $\rho_{1}\left(f_{1}, f_{2}\right):=E\left[\min \left(\left|Z\left(f_{1}\right)-Z\left(f_{2}\right)\right|, 1\right)\right]$ [see Arcones (1995)]. In particular, if $\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ converge weakly to $\{Z(f): f \in \mathscr{F}\}$, then

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} \operatorname{Pr}^{*}\left\{\sup _{\substack{\rho_{1}\left(f_{1}, f_{2}\right) \leq \delta \\ f_{1}, f_{2} \in \mathscr{F}}}\left|Z_{n}\left(f_{1}\right)-Z_{n}\left(f_{2}\right)\right| \geq \eta\right\}=0 \tag{1.19}
\end{equation*}
$$

for each $\eta>0$.
To study the almost sure Bahadur-Kiefer representation of $U$-quantiles instead of using weak convergence, we use a property similar to the compact law of the iterated logarithm: $\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ is a sequence of stochastic processes such that there is a subset $K$ of $l_{\infty}(\mathscr{F})$ satisfying that, with probability $1,\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ is relatively compact in $l_{\infty}(\mathscr{F})$ and its limit set is $K$. Given a sequence of stochastic processes $\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ and a subset $K$ of $l_{\infty}(\mathscr{F})$, we have that the following are equivalent:
(a) With probability $1,\left\{Z_{n}(f): f \in \mathscr{F}\right\}$ is relatively compact in $l_{\infty}(\mathscr{F})$ and its limit set is $K$.
(b) For each $f_{1}, \ldots, f_{m} \in \mathscr{F}$, with probability $1,\left\{\left(Z_{n}(f), \ldots, Z_{n}\left(f_{m}\right)\right)\right\}, n \geq 1$, is relatively compact in $\mathbb{R}^{m}$ and its limit set is $\left\{\left(x\left(f_{1}\right), \ldots, x\left(f_{m}\right)\right): x \in K\right\}$; and for each $\eta>0$, there exists a map $\pi: \mathscr{F} \rightarrow \mathscr{F}$ such that $\#\{\pi f: f \in \mathscr{F}\}$ is finite and

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \sup _{f \in \mathscr{F}}\left|Z_{n}(f)-Z_{n}(\pi f)\right| \leq \eta \quad \text { a.s. } \tag{1.20}
\end{equation*}
$$

If either (a) or (b) holds holds, $K$ is a compact set of $l_{\infty}(\mathscr{F})$. If $K$ is a compact set of $l_{\infty}(\mathscr{F})$, then $\left(\mathscr{F}, \rho_{2}\right)$ is totally bounded and $K \subset C_{u}\left(\mathscr{F}, \rho_{2}\right)$, where $\rho_{2}\left(f_{1}, f_{2}\right):=\sup _{x \in K}\left|x\left(f_{1}\right)-x\left(f_{2}\right)\right|$. In particular,

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} \sup _{\substack{\rho_{2}\left(f_{1}, f_{2}\right) \leq \delta \\ f_{1}, f_{2} \in \mathscr{F}}}\left|Z_{n}\left(f_{1}\right)-Z_{n}\left(f_{2}\right)\right|=0 \quad \text { a.s. } \tag{1.21}
\end{equation*}
$$

We refer for all last facts on the compact law of the iterated logarithm to Arcones and Giné (1995). Usually $K$ is the unit ball of the reproducing kernel Hilbert space (r.k.h.s.) of a covariance function on $\mathscr{F}$. A function $R: \mathscr{F} \times \mathscr{F} \rightarrow \mathbb{R}$
is a covariance function on $\mathscr{F}$, if

$$
\sum_{j=1}^{m} \sum_{k=1}^{m} a_{j} a_{k} R\left(f_{j}, f_{k}\right) \geq 0
$$

for each $a_{1}, \ldots, a_{m} \in \mathbb{R}$ and each $f_{1}, \ldots, f_{m} \in \mathscr{F}$. Then there is a mean-zero Gaussian process $\{W(f): F \in \mathscr{F}\}$ such that $E\left[W\left(f_{1}\right) W\left(f_{2}\right)\right]=R\left(f_{1}, f_{2}\right)$ for each $f_{1}, f_{2} \in \mathscr{F}$. Let $\mathscr{L}$ be the linear subspace of $L_{2}$, generated by $\{W(f)$ : $f \in \mathscr{F}\}$. The r.k.h.s. of the covariance function $R(\cdot, \cdot)$ is the following class of functions on $\mathscr{F}$ :
$\left\{(\alpha(f))_{f \in \mathscr{F}}\right.$ : there exists $\xi \in \mathscr{L}$ such that $\alpha(f)=E[W(f) \xi]$ for each $\left.f \in \mathscr{F}\right\}$.
This space is endowed with the inner product

$$
\left\langle\alpha_{1}, \alpha_{2}\right\rangle:=E\left[\xi_{1} \xi_{2}\right]
$$

if $\alpha_{i}(f)=E\left[W(f) \xi_{i}\right]$ for each $f \in \mathscr{F}$ and each $i=1,2$. The unit ball of this r.k.h.s. is

$$
K:=\left\{(E[W(f) \xi])_{f \in \mathscr{F}}: \xi \in \mathscr{L} \text { and } E\left[\xi^{2}\right] \leq 1\right\}
$$

Here, $\rho_{2}\left(f_{1}, f_{2}\right)=\sup _{x \in K}\left|x\left(f_{1}\right)-x\left(f_{2}\right)\right|=\left\|W\left(f_{1}\right)-W\left(f_{2}\right)\right\|_{2}$. A reference in r.k.h.s.'s is Aronszajn (1950).
2. The distributional Bahadur-Kiefer representation for $\boldsymbol{U}$-quan-
tiles. Here, we consider the distributional order of the Bahadur-Kiefer representation for $U$-quantiles. First, we give an upper bound to $R_{n}$.

THEOREM 1. With the above notation, suppose that:
(i) there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}, H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and

$$
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+O\left(\left(t-\xi_{0}\right)^{2}\right)
$$

as $t \rightarrow \xi_{0}$;
(ii) there is a sequence of real numbers $\left\{a_{n}\right\}_{n=1}^{\infty}$ such that

$$
a_{n}^{2} E\left[\left|g\left(X, \xi_{0}+t n^{-1 / 2}\right)-g\left(X, \xi_{0}\right)\right|^{2}\right]=O(1)
$$

for each $t \in \mathbb{R}$.
Then

$$
\begin{equation*}
a_{n} n^{1 / 2}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)=O_{\mathrm{P}}(1) \tag{2.1}
\end{equation*}
$$

Proof. c will design a finite constant which may vary from line to line. By (1.13) and hypotheses (i) and (ii),

$$
\begin{equation*}
a_{n}^{2} n^{-1}=O(1) \tag{2.2}
\end{equation*}
$$

Since, for each $\varepsilon>0, H_{n}\left(\xi_{0} \pm \varepsilon\right) \rightarrow H\left(\xi_{0} \pm \varepsilon\right)$ a.s. and $H\left(\xi_{0}-\varepsilon\right)<p<$ $H\left(\xi_{0}+\varepsilon\right)$, we have that

$$
\begin{equation*}
\xi_{n} \rightarrow \xi_{0} \quad \text { a.s. } \tag{2.3}
\end{equation*}
$$

Since the class $\left\{I_{h\left(X_{1}, \ldots, X_{m}\right) \leq t}: t \in \mathbb{R}\right\}$ is a VC subgraph class of functions, by Theorem 4.10 in Arcones and Giné (1993), $\left\{n^{1 / 2}\left(H_{n}(t)-H(t)\right): t \in \mathbb{R}\right\}$ converges weakly to a Gaussian process. From this, (2.3) and the fact that $E\left[\left|g(X, t)-g\left(X, \varepsilon_{0}\right)\right|^{2}\right] \rightarrow 0$, as $t \rightarrow \xi_{0}$, it follows that

$$
\begin{equation*}
n^{1 / 2}\left(H_{n}\left(\xi_{n}\right)-H_{n}\left(\xi_{0}\right)-H\left(\xi_{n}\right)+H\left(\xi_{0}\right)\right) \rightarrow_{\mathrm{P}} 0 \tag{2.4}
\end{equation*}
$$

By hypothesis (i), there is a $\delta>0$ such that $H(t)$ is continuous and increasing in $\left[\xi_{0}-\delta, \xi_{0}+\delta\right]$. Hence, for $\left\{i_{1}, \ldots, i_{m}\right\} \cap\left\{j_{1}, \ldots, j_{m}\right\}=\varnothing$ and $i_{1}<\cdots$ $<i_{m}$ and $j_{1}<\cdots<j_{m}$,

$$
\operatorname{Pr}\left\{h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right)=h\left(X_{j_{1}}, \ldots, X_{j_{m}}\right) \in\left[\xi_{0}-\delta, \xi_{0}+\delta\right]\right\}=0
$$

This implies that, for all $\left|s-\xi_{0}\right|<\delta$,

$$
\left|H_{n}(s)-H_{n}(s-)\right| \leq\binom{ n}{m}^{-1}\left|\binom{n}{m}-\binom{n-m}{m}\right| \leq c n^{-1} \quad \text { a.s. }
$$

Therefore, eventually

$$
\begin{equation*}
\left|H_{n}\left(\xi_{n}\right)-p\right| \leq c n^{-1} \quad \text { a.s. } \tag{2.5}
\end{equation*}
$$

By (2.4) and (2.5),

$$
\begin{equation*}
n^{1 / 2}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H\left(\xi_{n}\right)-H\left(\xi_{0}\right)\right) \rightarrow_{\mathrm{P}} 0 \tag{2.6}
\end{equation*}
$$

By hypothesis (i), there exists a positive constant $\eta$ such that if $\left|t-\xi_{0}\right| \leq \eta$, then $\left|H(t)-H\left(\xi_{0}\right)\right| \geq 2^{-1} H^{\prime}\left(\xi_{0}\right)\left|t-\xi_{0}\right|$. So, if $\left|\xi_{n}-\xi_{0}\right| \leq \eta$, then

$$
\begin{aligned}
2^{-1} H^{\prime}\left(\xi_{0}\right) n^{1 / 2}\left|\xi_{n}-\xi_{0}\right| \leq & n^{1 / 2}\left|H\left(\xi_{n}\right)-H\left(\xi_{0}\right)\right| \\
\leq & n^{1 / 2}\left|H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H\left(\xi_{n}\right)-H\left(\xi_{0}\right)\right| \\
& +n^{1 / 2}\left|H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)\right|=O_{\mathrm{P}}(1) .
\end{aligned}
$$

From this and (2.3), $n^{1 / 2}\left|\xi_{n}-\xi_{0}\right|=O_{\mathrm{Pr}}$ (1). The last estimation, (2.6) and hypothesis (i) imply that

$$
n^{1 / 2}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right) \rightarrow_{\mathrm{P}} 0
$$

Next, we show that, for each $M<\infty$,

$$
\begin{align*}
\sup _{|t| \leq M} a_{n} n^{1 / 2} \mid H_{n}\left(\xi_{0}+t n^{-1 / 2}\right) & -H\left(\xi_{0}+t n^{-1 / 2}\right)  \tag{2.7}\\
& -H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right) \mid=O_{\mathrm{P}}(1)
\end{align*}
$$

By the Hoeffding decomposition, it suffices to show that

$$
\begin{equation*}
\sup _{|t| \leq M} a_{n} n^{1 / 2}\left|\left(P_{n}-P\right)\left(g\left(\cdot, \xi_{0}+t n^{-1 / 2}\right)-g\left(\cdot, \xi_{0}\right)\right)\right|=O_{\mathrm{P}}(1) \tag{2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\sup _{|t| \leq M} a_{n} n^{1 / 2}\left|U_{n} \pi_{k, m}\left(I_{h \leq \xi_{0}+t n^{-1 / 2}}-I_{h \leq \xi_{0}}\right)\right| \rightarrow_{\mathrm{P}} 0 \tag{2.9}
\end{equation*}
$$

for $k \leq 2 \leq m$. Since the class $\{g(x, t): t \in \mathbb{R}\}$ is increasing in $\mathbb{R}$, it is a VC subgraph class. Hence, by (1.17),

$$
\begin{aligned}
& E\left[a_{n}^{2} n \sup _{|t| \leq M}\left|\left(P_{n}-P\right)\left(g\left(\cdot, \xi_{0}+t n^{-1 / 2}\right)-g\left(\cdot, \xi_{0}\right)\right)\right|^{2}\right] \\
& \leq c a_{n}^{2} E\left[\sup _{|t| \leq M}\left|g\left(X, \xi_{0}+t n^{-1 / 2}\right)-g\left(X, \xi_{0}\right)\right|^{2}\right] \\
& \leq c a_{n}^{2} E\left[\left|g\left(X, \xi_{0}+M n^{-1 / 2}\right)-g\left(X, \xi_{0}\right)\right|^{2}\right] \\
& +a_{n}^{2} E\left[\left|g\left(X, \xi_{0}-M n^{-1 / 2}\right)-g\left(X, \xi_{0}\right)\right|^{2}\right]=O(1) .
\end{aligned}
$$

[Observe that the classes $\mathscr{F}_{n}:=\left\{g\left(\cdot, \xi_{0}+t n^{-1 / 2}\right)-g\left(\cdot, \xi_{0}\right):|t| \leq M\right\}$ are all VC subgraph classes and $s\left(\mathscr{F}_{n}\right) \leq s\left(\mathscr{F}_{1}\right)$ for each $n$. So, we may choose a finite constant $c$ in (1.17) uniformly on $n$.] Therefore, (2.8) follows. (2.9) follows from (2.2) and Corollary 5.7 in Arcones and Giné (1993). Hence, (2.7) holds. Composing the process in (2.7) with $n^{1 / 2}\left(\xi_{n}-\xi_{0}\right)$, we get that

$$
a_{n} n^{1 / 2}\left(H_{n}\left(\xi_{n}\right)-H\left(\xi_{n}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right)=O_{\mathrm{P}}(1)
$$

By this, (2.2) and (2.5),

$$
a_{n} n^{1 / 2}\left(H\left(\xi_{0}\right)-H\left(\xi_{n}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right)=O_{\mathrm{P}}(1)
$$

So, the result follows.
From a previous theorem, it follows that if condition (i) holds and there exists a $1 \leq v \leq 2$ such that $E\left[\left|g(X, t)-g\left(X, \xi_{0}\right)\right|^{2}\right]=O\left(\left|t-\xi_{0}\right|^{v}\right)$ as $t \rightarrow \xi_{0}$, then

$$
n^{(v+2) / 4} R_{n}=O_{\mathrm{P}}(1)
$$

Next, we will find the exact order of this representation under some extra conditions. The exact order of this representation is determined by the order of

$$
\begin{equation*}
E\left[\left|g\left(X, \xi_{0}+t n^{-1 / 2}\right)-g\left(X, \xi_{0}\right)\right|^{2}\right] \tag{2.10}
\end{equation*}
$$

as $n \rightarrow \infty$. Finding the order of (2.10) could be difficult. By (1.13),

$$
\begin{aligned}
& n^{1 / 2} E\left[\left|g\left(X, \xi_{0}+t n^{-1 / 2}\right)-g\left(X, \xi_{0}\right)\right|^{2}\right] \\
& \quad \leq n^{1 / 2}\left|H\left(\xi_{0}+t n^{-1 / 2}\right)-H\left(\xi_{0}\right)\right|=O(1)
\end{aligned}
$$

that is, condition (ii) holds with $a_{n}=n^{1 / 4}$. So, under the easy-to-verify condition (i) in Theorem 1, we have that

$$
\begin{equation*}
n^{3 / 4}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)=O_{\mathrm{P}}(1) \tag{2.11}
\end{equation*}
$$

We will need the following CLT for triangular arrays indexed by VC classes. It follows from Theorem 2.6 in Alexander (1987).

Theorem 2. Let $\Theta$ be a subset of $\mathbb{R}^{d}$ and let $\theta_{0}$ be a point in the interior of $\Theta$. Let $g: S \times \Theta \rightarrow \mathbb{R}$ be a function such that $g(\cdot, \theta)$ is a measurable function for each $\theta \in \Theta$. Let $M<\infty$. Let $G_{R}(x)=\sup _{\left|\theta-\theta_{0}\right| \leq R} \mid g(x, \theta)-$ $g\left(x, \theta_{0}\right) \mid$, where $|\cdot|$ is the Euclidean distance $\mathbb{R}^{d}$. Suppose that:
(i) there is a $\delta_{0}>0$ such that $\left\{g(x, \theta)-g\left(x, \theta_{0}\right):\left|\theta-\theta_{0}\right| \leq \delta_{0}\right\}$ is a VC subgraph class;
(ii) there are sequences $\left\{a_{n}\right\}_{n=1}^{\infty}$ and $\left\{b_{n}\right\}_{n=1}^{\infty}$ such that $a_{n} \rightarrow \infty, b_{n} \rightarrow 0$ and

$$
\lim _{n \rightarrow \infty} a_{n}^{2} \operatorname{Var}\left(g\left(X, \theta_{0}+t b_{n}\right)-g\left(X, \theta_{0}+s b_{n}\right)\right)
$$

exists for each $|s|,|t| \leq M$;
(iii) $a_{n}^{2} E\left[G_{M b_{n}}^{2}(X)\right]=O(1)$;
(iv) $a_{n}^{2} E\left[G_{M b_{n}}^{2}(X) I_{G_{M b_{n}}(X) \geq \tau a_{n}^{-1} n^{1 / 2}}\right] \rightarrow 0$ for each $\tau>0$;
(v) $\lim _{\delta \rightarrow 0} \lim _{\sup _{n \rightarrow \infty}} \sup _{|t-s| \leq \delta,|s|,|t| \leq M} a_{n} \| g\left(X, \theta_{0}+t b_{n}\right)-g\left(X, \theta_{0}+\right.$ $\left.s b_{n}\right) \|_{2}=0$.

Then

$$
\begin{array}{r}
\left\{a _ { n } n ^ { - 1 / 2 } \sum _ { i = 1 } ^ { n } \left(g\left(X_{i}, \theta_{0}+t b_{n}\right)-g\left(X_{i}, \theta_{0}\right)-E\left[g\left(X, \theta_{0}+t b_{n}\right)\right]\right.\right. \\
\left.\left.+E\left[g\left(X, \theta_{0}\right)\right]\right):|t| \leq M\right\}
\end{array}
$$

converges weakly to the centered Gaussian process $\{Z(t):|t| \leq M\}$ determined by $Z(0)=0$ and $\|Z(t)-Z(s)\|_{2}=\rho(t, s)$.

Next, we see that under some conditions the order $a_{n}=n^{1 / 4}$ is attained.
Theorem 3. Suppose that:
(i) there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}, H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and there exists a finite constant $b$ such that

$$
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+b\left(t-\xi_{0}\right)^{2}+o\left(\left(t-\xi_{0}\right)^{2}\right)
$$

as $t \rightarrow \xi_{0}$;
(ii) there is a real number $\beta$ such that

$$
\lim _{\varepsilon \rightarrow 0+} \varepsilon^{-1} E\left[\left|g\left(X, \xi_{0}+\varepsilon t\right)-g\left(X, \xi_{0}+\varepsilon s\right)\right|^{2}\right]=\beta^{2}|t-s|
$$

for each $t, s \in \mathbb{R}$.
Then

$$
\begin{aligned}
& n^{3 / 4}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right) \\
& \quad \rightarrow_{d} m^{3 / 2} \beta\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1 / 2}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 4}\left|g_{1}\right|^{1 / 2} g_{2},
\end{aligned}
$$

where $g_{1}$ and $g_{2}$ are two independent standard normal r.v.'s.

Proof. We claim that by Theorem 2,

$$
\begin{equation*}
\left.\left\{n^{3 / 4} m\left(P_{n}-P\right)\left(g\left(\cdot, \xi_{0}+t n^{-1 / 2}\right)-g\left(\cdot, \xi_{0}\right)\right):|t| \leq M\right)\right\} \tag{2.12}
\end{equation*}
$$

converges weakly to $\{Z(t):|t| \leq M\}$, where $Z$ is a mean-zero Gaussian process with covariance given by

$$
E[Z(t) Z(s)]=2^{-1} m^{2} \beta(|t|+|s|-|t-s|)
$$

for each $s, t \in \mathbb{R}$. Observe that $g(x, t)$ is nondecreasing in $t$ for each fixed $x$. So, by a standard approximation argument,

$$
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} \sup _{\substack{|t-s| \leq \delta \\|s|| | t \mid \leq M}} n^{1 / 2} E\left[\left|g\left(X, \xi_{0}+t n^{-1 / 2}\right)-g\left(X, \xi_{0}+s n^{-1 / 2}\right)\right|^{2}\right]=0
$$

for each $M<\infty$. Observe also that $|g(x, t)| \leq 1$.
By Corollary 5.7 in Arcones and Giné (1993),
(2.13) $\sup _{|t| \leq M} n^{3 / 4}\left|U_{n} \pi_{k, m}\left(I_{h \leq \xi_{0}+t n^{-1 / 2}}-I_{h \leq \xi_{0}}\right)\right| \rightarrow_{\mathrm{P}} 0 \quad$ for $2 \leq k \leq m$.

From (2.12) and (2.13), the process

$$
\begin{array}{r}
\left\{Z_{n}(t):=n^{3 / 4}\left(H_{n}\left(\xi_{0}+t n^{-1 / 2}\right)-H\left(\xi_{0}+t n^{-1 / 2}\right)\right.\right. \\
\left.\left.-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right):|t| \leq M\right\}
\end{array}
$$

converges weakly to the process $\{Z(t)$ : $|t| \leq M\}$. Let $T_{M}^{*}=\{t:|t| \leq M\} \cup\{\infty\}$, let $Z_{n}(\infty):=n^{1 / 2}\left(\xi_{n}-\xi_{0}\right)$, let $Z(\infty)$ be a Gaussian r.v. with mean zero, covariance $m^{2}\left(H^{\prime}\left(\xi_{0}\right)\right)^{-2} \operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)$ and independent of the process $\{Z(t):|t| \leq$ $M\}$. Ву (2.11)

$$
\begin{equation*}
n^{1 / 2}\left(\xi_{n}-\xi_{0}\right)+n^{1 / 2}\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(N_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)\right) \rightarrow_{\mathrm{P}} 0 \tag{2.14}
\end{equation*}
$$

So, from this and the central limit theorem for triangular arrays, the finitedimensional distributions of $\left\{Z_{n}(t): t \in T_{M}^{*}\right\}$ converge to those of $\left\{Z(t): t \in T_{M}^{*}\right\}$. Condition (1.18) holds for $\left\{Z_{n}(t): t \in T_{M}^{*}\right\}$, because it holds for $\left\{Z_{n}(t):|t| \leq M\right\}$. Therefore, $\left\{Z_{n}(t): t \in T_{M}^{*}\right\}$ converges weakly to $\left\{Z(t): t \in T_{M}^{*}\right\}$.

By composing $Z_{n}(t)$ with $Z_{n}(\infty)$, we get that

$$
n^{3 / 4}\left(H_{n}\left(\xi_{n}\right)-H\left(\xi_{n}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right) \rightarrow_{\mathrm{d}} Z(Z(\infty))
$$

We have that $Z(\infty)$ has the distribution of $m\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 2} g_{1}$, where $g_{1}$ is a standard normal r.v. For each $t \in \mathbb{R}$, the distribution of

$$
Z\left(m\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 2} t\right)
$$

is that of

$$
m \beta m^{1 / 2}\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1 / 2}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 4}|t|^{1 / 2} g_{2}
$$

where $g_{2}$ is a standard normal r.v. So, by conditioning on $g_{1}$, we obtain that the distribution of

$$
Z\left(m\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 2} g_{1}\right)
$$

is that of

$$
m^{3 / 2} \beta\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1 / 2}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 4}\left|g_{1}\right|^{1 / 2} g_{2}
$$

where $g_{1}$ and $g_{2}$ are independent standard normal r.v.'s.
Theorem 3 applies to $h\left(x_{1}, \ldots, x_{m}\right)=\max _{1 \leq i \leq m} x_{i}$. Let $F(t)$ be the distribution function of $X_{i}$. Suppose that $\left(F\left(\xi_{0}\right)\right)^{m}=p, F$ is second differentiable at $\xi_{0}$ and $F^{\prime}\left(\xi_{0}\right)>0$. Then

$$
\begin{aligned}
n^{3 / 4} & \left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right) \\
& \rightarrow_{d} m p^{(4 m-3) / 4 m}\left(1-p^{1 / m}\right)^{1 / 4}\left|g_{1}\right|^{1 / 2} g_{2},
\end{aligned}
$$

where $g_{1}$ and $g_{2}$ are independent standard normal r.v.'s. Observe that $H(t)=(F(t))^{m}$ and

$$
g(x, t)=\operatorname{Pr}\left\{\max \left(x, X_{2}, \ldots, X_{m}\right) \leq t\right\}=I_{x \leq t}(F(t))^{m-1} .
$$

So, $H^{\prime}\left(\xi_{0}\right)=m p^{(m-1) / m} F^{\prime}\left(\xi_{0}\right)$ and $\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)=p^{(2 m-1) / m}\left(1-p^{1 / m}\right)$. For $s<t$

$$
\begin{aligned}
\epsilon^{-1} E & {\left[\left|g\left(X, \xi_{0}+t \epsilon\right)-g\left(X, \xi_{0}+s \epsilon\right)\right|^{2}\right] } \\
= & \left.\epsilon^{-1}\left(\left(F\left(\xi_{0}+t \epsilon\right)\right)^{m-1}-F\left(\xi_{0}+s \epsilon\right)\right)^{m-1}\right)^{2} F\left(\xi_{0}+s \epsilon\right) \\
& \quad+\epsilon^{-1}\left(F\left(\xi_{0}+t \epsilon\right)\right)^{2 m-2}\left(F\left(\xi_{0}+t \epsilon\right)-F\left(\xi_{0}+\epsilon s\right)\right),
\end{aligned}
$$

which converges to $p^{2(m-1) / m}(t-s) F^{\prime}\left(\xi_{0}\right)$. We also have that

$$
\epsilon^{-1 / 2} E\left[\left(g\left(X, \xi_{0}+t \epsilon\right)-g\left(X, \xi_{0}\right)\right) g\left(X, \xi_{0}\right)\right] \rightarrow 0
$$

Next, we see how the order $n$ in the Bahadur-Kiefer representation of $U$-quantiles can be attained.

Theorem 4. Suppose that:
(i) there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and there exists a finite constant $b$ such that

$$
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+b\left(t-\xi_{0}\right)^{2}+o\left(\left(t-\xi_{0}\right)^{2}\right)
$$

as $t \rightarrow \xi_{0}$ for some $b \in \mathbb{R}$;
(ii) there is a real number $\beta$ such that

$$
\lim _{\varepsilon \rightarrow 0+} \varepsilon^{-2} \operatorname{Var}\left(g\left(X, \xi_{0}+\varepsilon t\right)-g\left(X, \xi_{0}+\varepsilon s\right)\right)=\beta(t-s)^{2}
$$

for each $t, s \in \mathbb{R}$;
(iii) $\lim _{\varepsilon \rightarrow 0+} \varepsilon^{-2} E\left[\left|g\left(X, \xi_{0}+\varepsilon t\right)-g\left(X, \xi_{0}\right)\right|^{2} I_{\left|g\left(X, \xi_{0}+\varepsilon t\right)-g\left(X, \xi_{0}\right)\right| \geq \tau}\right]=0$
for each $t \in \mathbb{R}$ and each $\tau>0$;
(iv) there is a real number $\alpha$ such that

$$
\lim _{\epsilon \rightarrow 0+} \epsilon^{-1} \operatorname{Cov}\left(g\left(X, \xi_{0}+\epsilon t\right)-g\left(X, \xi_{0}\right), g\left(X, \xi_{0}\right)\right)=\alpha t
$$

for each $t \in \mathbb{R}$.
(v) There is a $\delta>0$ such that for any $\left|t-\xi_{0}\right|<\delta$ and for any combinations $i_{1}<\cdots<i_{m}$ and $j_{1}<\cdots<j_{m}$ such that $\left\{i_{1}, \ldots, i_{m}\right\} \neq\left\{j_{1}, \ldots, j_{m}\right\}$, we have that

$$
\operatorname{Pr}\left\{h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right)=h\left(X_{j_{1}}, \ldots, X_{j_{m}}\right)=t\right\}=0 .
$$

Then

$$
n\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)
$$

converges in distribution to $\lambda_{1} g_{1}^{2}+\lambda_{2} g_{2}^{2}$, where $g_{1}$ and $g_{2}$ are independent standard normal random variables,

$$
\begin{align*}
\lambda_{1} & =-2^{-1}\left(c_{1,2}+b c_{2,2}\right)-2^{-1}\left(c_{1,1} c_{2,2}+2 b c_{1,2} c_{2,2}+b^{2} c_{2,2}^{2}\right)^{1 / 2},  \tag{2.15}\\
\lambda_{2} & =-2^{-1}\left(c_{1,2}+b c_{2,2}\right)+2^{-1}\left(c_{1,1} c_{2,2}+2 b c_{1,2} c_{2,2}+b^{2} c_{2,2}^{2}\right)^{1 / 2},  \tag{2.16}\\
c_{1,1} & =m^{2} \beta^{2}, c_{2,2}=m^{2}\left(H^{\prime}\left(\xi_{0}\right)\right)^{-2} \operatorname{Var}\left(g\left(X, \xi_{0}\right)\right) \text { and }  \tag{2.17}\\
c_{1,2} & =-m^{2}\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1} \alpha .
\end{align*}
$$

Proof. Observe that, by hypotheses (i) and (ii), $m \geq 2$. By the method in the proof of Theorem 3,

$$
\left\{n\left(H_{n}\left(\xi_{0}+t n^{-1 / 2}\right)-H\left(\xi_{0}+t n^{-1 / 2}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right):|t| \leq M\right\}
$$

and $n^{1 / 2}\left(\xi_{n}-\xi_{0}\right)$ converge jointly to $\left\{t Y_{1}:|t| \leq M\right\}$ and $Y_{2}$, where $Y_{1}$ and $Y_{2}$ are jointly normal random variables with mean zero and

$$
E\left[Y_{1}^{2}\right]=c_{1,1}, E\left[Y_{2}^{2}\right]=c_{2,2} \quad \text { and } \quad E\left[Y_{1} Y_{2}\right]=c_{1,2},
$$

where $c_{1,1}, c_{2,2}$ and $c_{1,2}$ are as in (2.17). Hence

$$
n\left(H_{n}\left(\xi_{n}\right)-H\left(\xi_{n}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)+b\left(\xi_{n}-\xi_{0}\right)^{2}\right)
$$

converges in distribution to $Y_{1}, Y_{2}+b Y_{2}^{2}$. We have that all $h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right)$ in $\left[\xi_{0}-t, \xi_{0}+t\right]$ are different. So, $\left|H_{n}\left(\xi_{n}\right)-p\right| \leq\binom{ n}{m}^{-1}$. Therefore,

$$
n\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)
$$

converges in distribution to $-Y_{1} Y_{2}-b Y_{2}^{2}$. If $c_{1,1} c_{2,2}-c_{1,2}^{2}=0, Y_{1}$ and $Y_{2}$ are linearly dependent, and the distribution of $-Y_{1} Y_{2}-b Y_{2}^{2}$ is that of $\lambda g^{2}$, where $g$ is a standard normal random variable and

$$
\lambda=E\left[-Y_{1} Y_{2}-b Y_{2}^{2}\right]=-\left(c_{1,2}+b c_{2,2}\right) .
$$

If $c_{1,1} c_{2,2}-c_{1,2}^{2}=0$, then

$$
\lambda_{1} g_{1}^{2}+\lambda_{2} g_{2}^{2}=-\left(c_{1,2}+b c_{2,2}\right)^{+} g_{1}^{2}-\left(c_{1,2}+b c_{2,2}\right)^{-} g_{2}^{2},
$$

where $x^{-}=\max (-x, 0), x^{+}=\max (x, 0), \lambda_{1}$ and $\lambda_{2}$ are as in (2.15) and (2.16) and $g_{1}$ and $g_{2}$ are two independent standard normal random variables. So, $\lambda_{1} g_{1}^{2}+\lambda_{2} g_{2}^{2}$ has the same distribution as $\lambda g^{2}$. Assume now that $c_{1,1} c_{2,2}-$ $c_{1,2}^{2}>0$. Let $Z_{1}=c_{2,2}^{-1 / 2}\left(c_{1,1} c_{2,2}-c_{1,2}^{2}\right)^{-1 / 2}\left(c_{2,2} Y_{2}-c_{1,2} Y_{2}\right)$ and $Z_{2}=$ $c_{2,2}^{-1 / 2} Y_{2}$. We have that $Z_{1}$ and $Z_{2}$ are two independent standard normal random variables and $-Y_{1} Y_{2}-b Y_{2}^{2}=-a_{1} Z_{1} Z_{2}-a_{2} Z_{2}^{2}$, where $a_{1}=$ $\left(c_{1,1} c_{2,2}-c_{1,2}^{2}\right)^{1 / 2}$ and $a_{2}=c_{1,2}+b c_{2,2}$. Let $g_{1}=q_{1}^{-1} a_{1} Z_{1}+q_{1}^{-1}\left(a_{2}+\left(a_{1}^{2}+\right.\right.$ $\left.\left.a_{2}^{2}\right)^{1 / 2}\right) Z_{2}$ and $g_{2}=q_{2}^{-1} a_{1} Z_{1}+q_{2}^{-1}\left(a_{2}-\left(a_{1}^{2}+a_{2}^{2}\right)^{1 / 2}\right) Z_{2}$, where

$$
\begin{align*}
& q_{1}=2^{1 / 2}\left(a_{1}^{2}+a_{2}^{2}+a_{2}\left(a_{1}^{2}+a_{2}^{2}\right)^{1 / 2}\right)^{1 / 2} \text { and } \\
& q_{2}=2^{1 / 2}\left(a_{1}^{2}+a_{2}^{2}-a_{2}\left(a_{1}^{2}+a_{2}^{2}\right)^{1 / 2}\right)^{1 / 2} \tag{2.18}
\end{align*}
$$

Then, $g_{1}$ and $g_{2}$ are two independent standard normal random variables and

$$
-a_{1} Z_{1} Z_{2}-a_{2} Z_{2}^{2}=\lambda_{1} g_{1}^{2}+\lambda_{2} g_{2}^{2},
$$

where $\lambda_{1}$ and $\lambda_{2}$ are as in (2.15) and (2.16). So, the result follows.
Theorem 4 applies to $h\left(x_{1}, x_{2}\right)=x_{1}+x_{2}$. Suppose that $\operatorname{Pr}\left\{X_{1}+X_{2} \leq\right.$ $\left.\xi_{0}\right\}=p$, the distribution function $F$ of $X$ is second differentiable and its first and second derivatives are uniformly bounded and $\int_{-\infty}^{\infty} F^{\prime}\left(\xi_{0}-x\right) F^{\prime}(x) d x>$ 0 . Then the thesis of Theorem 4 holds for $h\left(x_{1}, x_{2}\right)=x_{1}+x_{2}$ with

$$
\begin{aligned}
H^{\prime}\left(\xi_{0}\right)= & \int_{-\infty}^{\infty} F^{\prime}\left(\xi_{0}-x\right) F^{\prime}(x) d x \\
b= & 2^{-1} \int_{-\infty}^{\infty} F^{\prime \prime}\left(\xi_{0}-x\right) F^{\prime}(x) d x \\
\beta^{2}= & \int_{-\infty}^{\infty}\left(F^{\prime}\left(\xi_{0}-x\right)\right)^{2} F^{\prime}(x) d x-\left(\int_{-\infty}^{\infty} F^{\prime}\left(\xi_{0}-x\right) F^{\prime}(x) d x\right)^{2} \\
\alpha= & \int_{-\infty}^{\infty} F^{\prime}\left(\xi_{0}-x\right) F\left(\xi_{0}-x\right) F^{\prime}(x) d x \\
& -\int_{-\infty}^{\infty} F^{\prime}\left(\xi_{0}-x\right) F^{\prime}(x) d x \int_{-\infty}^{\infty} F\left(\xi_{0}-x\right) F^{\prime}(x) d x
\end{aligned}
$$

and $g(x, t)=F(t-x)$.
Theorem 4 also applies to the kernel $h\left(x_{1}, x_{2}\right)=\left|x_{1}-x_{2}\right|^{r}$, where $r>0$. We omit the details. As mentioned in the Introduction, the $U$-quantile over this kernel, in the case $r=1$ and $p=1 / 2$, was considered by Bickel and Lehmann (1979). It is an estimator of the spread of the distribution.

Theorem 4 also applies to the kernel $h\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right)\right)=\left(y_{2}-y_{1}\right) /$ ( $x_{2}-x_{1}$ ) (under some regularity conditions). Consider the linear regression model: $Y_{i}=\alpha+\beta X_{i}+\delta_{i} ; \alpha$ and $\beta$ are constants and $\delta_{i}$ is an r.v. independent of $X_{i}$. The $U$-quantile over the kernel $h$, with respect to $p=1 / 2$, is an estimator of the regression slope $\beta$. Here

$$
g\left(\left(x_{1}, y_{1}\right), t\right)=\operatorname{Pr}\left\{(\beta-t) X_{2}+\delta_{2} \leq y_{1}-t x_{1}-\alpha\right\} .
$$

3. The a.s. Bahadur-Kiefer representation of $\boldsymbol{U}$-quantiles. First, we present some results on the law of the iterated logarithm for processes. The following lemma is similar in spirit to Theorem 3.1 in Kuelbs (1976).

Lemma 5. Let $\left\{Z_{n}(t): t \in T\right\}$ be a stochastic process indexed by a parameter set $T$. Let $\rho$ be a pseudometric on $T$. Let $K$ be a compact subset of $C_{u}(T, \rho)$. Assume that the following conditions are satisfied:
(i) $(T, \rho)$ is totally bounded;
(ii) $\lim _{\delta \rightarrow 0} \lim \sup _{n \rightarrow \infty} \sup _{\rho\left(t_{1}, t_{2}\right) \leq \delta}\left|Z_{n}\left(t_{1}\right)-Z_{n}\left(t_{2}\right)\right|=0$ a.s.;
(iii) for each $m \in \mathbb{N}$ and each $t_{1}, \ldots, t_{m} \in T$, with probability 1 , the sequence $\left\{\left(Z_{n}\left(t_{1}\right), \ldots, Z_{n}\left(t_{m}\right)\right)_{n=1}^{\infty}\right.$ is relatively compact in $\mathbb{R}^{m}$ and its limit set is contained in $\left\{\left(x\left(t_{1}\right), \ldots, x\left(t_{m}\right)\right): x \in K\right\}$.

Then, with probability 1 , the sequence $\left\{Z_{n}(t): t \in T\right\}$ is relatively compact in $l_{\infty}(T)$ and its limit set is contained in $K$.

Proof. By the Arzelá-Ascoli theorem, with probability $1\left\{Z_{n}(t): t \in T\right\}$ is relatively compact in $l_{\infty}(T)$. Let $\left\{t_{p}\right\}_{p=1}^{\infty}$ be a countable dense subset of ( $T, \rho$ ). Let $A$ be a measurable set having probability 1 , such that in $A$,

$$
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} \sup _{\rho\left(t_{1}, t_{2}\right) \leq \delta}\left|Z_{n}\left(t_{1}\right)-Z_{n}\left(t_{2}\right)\right|=0
$$

and for each $m \in \mathbb{N}$ the sequence $\left\{\left(Z_{n}\left(t_{1}\right), \ldots, Z_{n}\left(t_{m}\right)\right)\right\}_{n=1}^{\infty}$ is relatively compact in $\mathbb{R}^{m}$ and its limit set is contained in $\left\{\left(x\left(t_{1}\right), \ldots, x\left(t_{m}\right)\right): x \in K\right\}$. Suppose that $x \in l_{\infty}(T)$ is a limit point of a sequence $\left\{Z_{n}(t): t \in T\right\}$ satisfying the previous two conditions. By the first condition, $x \in C_{u}(T, \rho)$. Let $n_{j}$ be a subsequence such that $Z_{n_{i}} \rightarrow x$. By the second condition, for each $m \geq 1$, there is an $x^{(m)} \in K$ such that

$$
\left(x\left(t_{1}\right), \ldots, x\left(t_{m}\right)\right)=\left(x^{(m)}\left(t_{1}\right), \ldots, x^{(m)}\left(t_{m}\right)\right) .
$$

Since $K$ is compact, the sequence $\left\{x^{(m)}\right\}_{m=1}^{\infty}$ has a limit point $y \in K$. Since $x, y \in C_{u}(T, \rho)$ and $x\left(t_{p}\right)=y\left(t_{p}\right)$ for each $p \geq 1, x=y$.

We also need the following law of the iterated logarithm.
Theorem 6. Let $\left\{X_{j}\right\}_{j=1}^{\infty}$ be a sequence of i.i.d. r.v.'s with values in a measurable space $(S, \mathscr{S})$. Let $g: S \times T \rightarrow \mathbb{R}$ be a function such that $g(\cdot, t): S \rightarrow \mathbb{R}$ is a measurable function for each $t \in T$. Let $R(\cdot, \cdot)$ be a covariance function on T. Let $\left\{b_{n}\right\}$ be a sequence of real numbers from the interval $(0,1]$ and let $\left\{a_{n}\right\}$ be a sequence of positive real numbers. Suppose that:
(i) there is a scalar product defined for each $t \in T$ and each $0 \leq u \leq 1$, so that $u t \in T$;
(ii) $\limsup _{n \rightarrow \infty} a_{n}^{2} n^{-1} \operatorname{Var}\left(\sum_{j=1}^{p} \lambda_{j} g\left(X, b_{n} t_{j}\right)\right) \leq \sum_{j, k=1}^{p} \lambda_{j} \lambda_{k} R\left(t_{j}, t_{k}\right)$ for each $t_{1}, \ldots, t_{p} \in T$ and each $\lambda_{1}, \ldots, \lambda_{p} \in \mathbb{R} ;$
(iii) $\left\{b_{n}\right\}$ and $\left\{a_{n} n^{-1}(\log \log n)^{-1 / 2}\right\}$ are nonincreasing sequences;
(iv) $\quad \lim _{\gamma \rightarrow 1+} \lim \sup _{n \rightarrow \infty} \sup _{m: n \leq m \leq \gamma n} a_{n}^{-1}\left|a_{m}-a_{n}\right|=0$
and $\lim _{\gamma \rightarrow 1+} \lim \sup _{n \rightarrow \infty} \sup _{m: n \leq m \leq \gamma n} b_{n}^{-1}\left|b_{m}-b_{n}\right|=0$;
(v) for each $t \in T, \lim _{u \rightarrow 1_{-}} \rho(u t, t)=0$, where

$$
\rho^{2}(s, t):=R(s, s)+R(t, t)-2 R(s, t) .
$$

(vi) $\quad \lim _{\tau \rightarrow 0+} \limsup _{n \rightarrow \infty} \sup _{t \in T} n^{-1} a_{n}^{2} E\left[\left(g\left(X, b_{n} t\right)-E\left[g\left(X, b_{n} t\right)\right]\right)^{2}\right]$

$$
\left.\times I_{G\left(X, b_{n}\right) \geq \tau(\log \log n)^{-1 / 2} a_{n}^{-1} n}\right)=0,
$$

where $G(x):=\sup _{t \in T} \mid g\left(x, b_{n} t\right)-E\left[g\left(X, b_{n} t\right)\right] ;$
(vii) there are positive constants $r_{1}$ and $r_{2}$, such that

$$
\sum_{j=2}^{\infty} \sup _{r_{1}(\log j)^{-1} \leq r \leq r_{2}} e^{-r_{1} r^{-1}} 2^{j} \operatorname{Pr}\left\{G\left(X, b_{2^{j}}\right) \geq r 2^{j}(\log j)^{1 / 2} a_{2_{j}}^{-1}\right\}<\infty ;
$$

(viii) ( $T, \rho$ ) is totally bounded;
(ix) $a_{n}(2 \log \log n)^{-1 / 2} \sup _{t \in T}\left|\left(P_{n}-P\right) g\left(\cdot, b_{n} t\right)\right| \rightarrow{ }_{P} 0$;
(x) $\lim \sup _{\delta \rightarrow 0} \lim \sup _{n \rightarrow \infty} \sup _{\rho(s, t) \leq \delta} a_{n}^{2} n^{-1} \operatorname{Var}\left(g\left(X, b_{n} s\right)-\right.$ $\left.g\left(X, b_{n} t\right)\right)=0$.

Then, with probability 1,

$$
\begin{equation*}
\left\{(2 \log \log n)^{-1 / 2} a_{n}\left(P_{n}-P\right) g\left(\cdot, b_{n} t\right): t \in T\right\}, \quad n \geq 1, \tag{3.1}
\end{equation*}
$$

is relatively compact in $l_{\infty}(T)$ and its limit set is contained in the unit ball of the r.k.h.s. of the covariance function $R(\cdot, \cdot)$.

Proof. It follows from Lemma 5, using the method in the proof of Theorem 3.1 in Arcones (1994b). So, we omit the proof.

Since the first element of the Hoeffding decomposition of $I_{h \leq t}$ can be difficult to find, we first give a sharp upper bound.

Theorem 7. Let $h: S^{m} \rightarrow \mathbb{R}$ be a symmetric measurable function and let $0<p<1$. Suppose that there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}$, $H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and there exists a finite constant $b$ such that

$$
\begin{equation*}
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+O\left(\left(t-\xi_{0}\right)^{2}\right) \tag{3.2}
\end{equation*}
$$

as $t \rightarrow \xi_{0}$. Then
$\lim \sup (n / 2 \log \log n)^{3 / 4}\left|H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right| \leq l \quad$ a.s., $n \rightarrow \infty$
where $l=2^{1 / 2} 3^{-3 / 4} m^{3 / 2}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 4}$.
Proof. Take

$$
M>m\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 2}
$$

Let $b_{n}=(2 \log \log n / n)^{1 / 2}$ and let
$Z_{n}(t):=(n / 2 \log \log n)^{3 / 4}\left(H_{n}\left(\xi_{0}+t b_{n}\right)-H\left(\xi_{0}+t b_{n}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right)$.
By Theorem 4.7 in Arcones and Giné (1995),

$$
\begin{equation*}
\sup _{|t| \leq M}(n / 2 \log \log n)^{3 / 4}\left|U_{n} \pi_{k, m}\left(I_{h \leq \xi_{0}+t b_{n}}-I_{h \leq \xi_{0}}\right)\right| \rightarrow 0 \quad \text { a.s. } \tag{3.3}
\end{equation*}
$$

We claim that by Theorem 6 , with probability 1 ,

$$
\begin{equation*}
\left\{(n / 2 \log \log n)^{3 / 4} m\left(P_{n}-P\right)\left(g\left(\cdot, \xi_{0}+b_{n} t\right)-g\left(\cdot, \xi_{0}\right)\right):|t| \leq M\right\} \tag{3.4}
\end{equation*}
$$

is relatively compact and its limit set is contained in the unit ball of the r.k.h.s. of the mean-zero Gaussian process $\{Z(t):|t| \leq M\}$ having covariance $E[Z(t) Z(s)]=2^{-1} m^{2} H^{\prime}\left(\xi_{0}\right)(|s|+|t|-|s-t|)$ for each $s, t \in \mathbb{R}$; that is, the limit set is contained in

$$
\begin{equation*}
K_{0}:=\left\{(\gamma(t))_{|t| \leq M}: \gamma(0)=0 \text { and } \int_{-M}^{M}\left(\gamma^{\prime}(t)\right)^{2} d t \leq m^{2} H^{\prime}\left(\xi_{0}\right)\right\} . \tag{3.5}
\end{equation*}
$$

We are applying Theorem 6 with $T=[-M, M], R(s, t)=E[Z(t) Z(s)]$ and $a_{n}=n^{3 / 4}(2 \log \log n)^{-1 / 4}$. We have that

$$
\begin{aligned}
& \limsup _{n \rightarrow \infty} b_{n}^{-1} m^{2} \operatorname{Var}\left(\sum_{j=1}^{p} \lambda_{j}\left(g\left(X, \xi_{0}+b_{n} t_{j}\right)-g\left(X, \xi_{0}\right)\right)\right) \\
& \quad \leq \sum_{j, k=1}^{p} \lambda_{j} \lambda_{k} R\left(t_{j}, t_{k}\right),
\end{aligned}
$$

because

$$
b_{n}^{-1} m^{2}\left(E\left[\sum_{j=1}^{p} \lambda_{j}\left(g\left(X, \xi_{0}+b_{n} t_{j}\right)-g\left(X, \xi_{0}\right)\right)\right]\right)^{2} \rightarrow 0
$$

and

$$
\begin{aligned}
& b_{n}^{-1} m^{2} E\left[\left(\sum_{j=1}^{p} \lambda_{j}\left(g\left(X, \xi_{0}+b_{n} t_{j}\right)-g\left(X, \xi_{0}\right)\right)\right)^{2}\right] \\
& \quad \leq b_{n}^{-1} m^{2} E\left[\left(\sum_{j=1}^{p} \lambda_{j}\left(I_{h \leq \xi_{0}+b_{n} t_{j}}-I_{h \leq \xi_{0}}\right)\right)^{2}\right] \rightarrow \sum_{j, k=1}^{p} \lambda_{j} \lambda_{k} R\left(t_{j}, t_{k}\right) .
\end{aligned}
$$

The rest of the hypotheses in Theorem 6 are either trivial or can be checked similarly to the conditions already checked. So, (3.4) and (3.5) follow. By (3.3)-(3.5),

$$
\begin{equation*}
\left\{Z_{n}(t):|t| \leq M\right\} \tag{3.6}
\end{equation*}
$$

is a.s. relatively compact and its limit set is contained in $K_{0}$. Let

$$
\begin{equation*}
Z_{n}(\infty):=(n / 2 \log \log n)^{1 / 2}\left(\xi_{n}-\xi_{0}\right) \tag{3.7}
\end{equation*}
$$

and let $Z(\infty)$ be a Gaussian r.v. with mean 0 and variance

$$
m^{2}\left(H^{\prime}\left(\xi_{0}\right)\right)^{-2} \operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)
$$

and independent of the process $\{Z(t):|t| \leq M\}$, let $T_{M}^{*}=[-M, M] \cup\{\infty\}$ and let $\rho(s, t)=\|Z(t)-Z(s)\|_{2}$, where $s, t \in T_{M}^{*}$. Using the method in Remark 10 in Arcones (1994a) (applying Kolmogorov's exponential inequalities), it is easy to see that for any $t_{1}, \ldots, t_{p} \in T_{M}^{*},\left\{\left(Z_{n}\left(t_{1}\right), \ldots, Z_{n}\left(t_{p}\right)\right)\right\}$ is a.s. relatively compact and its limit set is contained in the unit ball of $\left(Z\left(t_{1}\right), \ldots, Z\left(t_{p}\right)\right)$. Since $\left\{Z_{n}(t):|t| \leq M\right\}$ satisfies a type of compact law of the iterated logarithm,

$$
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} \sup _{\substack{\left|t_{1}\right|| | t_{2} \mid \leq M \\ \rho\left(t_{1}, t_{2}\right) \leq \delta}}\left|Z_{n}\left(t_{1}\right)-Z_{n}\left(t_{2}\right)\right|=0 \quad \text { a.s. }
$$

and ( $[-M, M], \rho$ ) is totally bounded. So, it is also true that

$$
\lim _{\delta \rightarrow 0} \limsup _{n \rightarrow \infty} \sup _{\substack{t_{1}, t_{2} \in T_{\mathcal{N}}^{*} \\ \rho\left(t_{1}, t_{2}\right) \leq \delta}}\left|Z_{n}\left(t_{1}\right)-Z_{n}\left(t_{2}\right)\right|=0 \quad \text { a.s. }
$$

and $\left(T_{m}^{*}, \rho\right)$ is totally bounded. So, by Lemma 5 the process $\left\{Z_{n}(t): t \in T_{M}^{*}\right\}$ is a.s. relatively compact and its limit set is contained in the reproducing kernel Hilbert space of $\left\{Z(t): t \in T_{M}^{*}\right\}$; that is, it is contained in

$$
\begin{aligned}
& K_{1}:=\left\{(\gamma(t))_{t \in T_{M}^{*}}: \gamma(0)=0\right. \text { and } \\
& \qquad \begin{aligned}
&\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1} \int_{-M}^{M}\left(\gamma^{\prime}(t)\right)^{2} d t \\
&\left.+\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{-1}\left(H^{\prime}\left(\xi_{0}\right)\right)^{2}(\gamma(\infty))^{2} \leq m^{2}\right\} .
\end{aligned}
\end{aligned}
$$

By Theorem 4.1 in Arcones (1993),

$$
\begin{equation*}
(n / 2 \log \log n)^{1 / 2}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right) \rightarrow 0 \quad \text { a.s. } \tag{3.8}
\end{equation*}
$$

From this, the law of the iterated logarithm of $\left\{Z_{n}(t): t \in T_{M}^{*}\right\}$ and composition

$$
\left\{(n / 2 \log \log n)^{3 / 4}\left(H_{n}\left(\xi_{n}\right)-H\left(\xi_{n}\right)+H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)\right)\right\}
$$

is a.s. relatively compact and its limit set is contained in

$$
\begin{align*}
\{\gamma(v): \gamma(0) & =0 \quad \text { and } \quad H^{\prime}\left(\xi_{0}\right)^{-1} \int_{-M}^{M}\left(\gamma^{\prime}(t)\right)^{2} d t  \tag{3.9}\\
& \left.+\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{-1}\left(H^{\prime}\left(\xi_{0}\right)\right)^{2} v^{2} \leq m^{2}\right\}
\end{align*}
$$

By the argument in Proposition 1.1 in Deheuvels and Mason (1992), the set in (3.9) is $[-l, l]$, where $l=2^{1 / 2} 3^{-3 / 4} m^{3 / 2}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 4}$.

Next, we see how the order $(n / 2 \log \log n)^{3 / 4}$ can be attained in the Bahadur-Kiefer representation of $U$-quantiles.

Theorem 8. Let $h: S^{m} \rightarrow \mathbb{R}$ be a symmetric measurable function. Let $0<p<1$. Suppose that:
(i) there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}, H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and there exists a finite constant $b$ such that

$$
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+b\left(t-\xi_{0}\right)^{2}+o\left(\left(t-\xi_{0}\right)^{2}\right)
$$

as $t \rightarrow \xi_{0}$;
(ii) there is a real number $\beta$ such that

$$
\lim _{\varepsilon \rightarrow 0+} \varepsilon^{-1} E\left[\left|g\left(X, \xi_{0}+\varepsilon t\right)-g\left(X, \xi_{0}+\varepsilon s\right)\right|^{2}\right]=\beta^{2}|t-s|
$$

for each $t, s \in \mathbb{R}$.
(iii) $\lim _{\epsilon \rightarrow 0+} \epsilon^{-1 / 2} E\left[\left(g\left(X, \xi_{0}+\epsilon t\right)-g\left(X, \xi_{0}\right)\right) g\left(X, \xi_{0}\right)\right]=0$
for each $t \in \mathbb{R}$.
Then, with probability 1,

$$
\left\{(n / 2 \log \log n)^{3 / 4}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)\right\}
$$

is relatively compact and its limit set is $[-l, l]$, where

$$
l=2^{1 / 2} 3^{-3 / 4} m^{3 / 2} \beta\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1 / 2}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 4}
$$

In particular,

$$
\limsup _{n \rightarrow \infty}(n / 2 \log \log n)^{3 / 4}\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)=l \quad \text { a.s. }
$$

Proof. Take $M>m\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 2}$. The proof is similar to the proof of Theorem 7. We have that by Theorem 3.1 in Arcones (1994b)

$$
\begin{aligned}
&\left\{( n / 2 \operatorname { l o g } \operatorname { l o g } n ) ^ { 3 / 4 } m ( P _ { n } - P ) \left(g \left(\cdot, \xi_{0}+\right.\right.\right.\left.t(2 \log \log n / n)^{1 / 2}\right) \\
&\left.\left.-g\left(\cdot, \xi_{0}\right)\right):|t| \leq M\right\}
\end{aligned}
$$

is a.s. relatively compact and its limit set is

$$
\left\{(\gamma(t))_{|t| \leq M}: \gamma(0)=0 \quad \text { and } \quad \int_{-M}^{M}\left(\gamma^{\prime}(t)\right)^{2} d t \leq m^{2} \beta^{2}\right\} .
$$

So, the arguments in Theorem 7 imply the result.
Theorem 8 applies to $h\left(x_{1}, \ldots, x_{m}\right)=\max _{1 \leq i \leq m} x_{i}$, if there exists $\xi_{0}$ with $\left(F\left(\xi_{0}\right)\right)^{m}=p$, where $F(t)$ is the distribution function of $X_{i}, F$ is second differentiable at $\xi_{0}$ and $F^{\prime}\left(\xi_{0}\right)>0$. In this case

$$
l=2^{1 / 2} 3^{-3 / 4} m p^{(4 m-3) / 4 m}\left(1-p^{1 / m}\right)^{1 / 4} .
$$

Theorem 9. Let $h: S^{m} \rightarrow \mathbb{R}$ be a symmetric measurable function. Let $0<p<1$. Suppose that:
(i) there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}, H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and there exists a finite constant b such that

$$
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+b\left(t-\xi_{0}\right)^{2}+o\left(\left(t-\xi_{0}\right)^{2}\right)
$$

as $t \rightarrow \xi_{0}$;
(ii) there is a real number $\beta$ such that

$$
\lim _{\varepsilon \rightarrow 0+} \varepsilon^{-2} \operatorname{Var}\left(g\left(X, \xi_{0}+\varepsilon t\right)-g\left(X, \xi_{0}+\varepsilon s\right)\right)=\beta^{2}(t-s)^{2}
$$

for each $t, s \in \mathbb{R}$;
(iii) $\lim _{\varepsilon \rightarrow 0+} \varepsilon^{-2} E\left[\lg \left(X, \xi_{0}+\varepsilon M\right)\right.$

$$
\left.-\left.g\left(X, \xi_{0}-\varepsilon M\right)\right|^{2} I_{\left|g\left(X, \xi_{0}+\varepsilon M\right)-g\left(X, \xi_{0}-\varepsilon M\right)\right| \geq \tau}\right]=0
$$

for each $M, \tau>0$;
(iv) there is a real number $\alpha$ such that

$$
\lim _{\epsilon \rightarrow 0+} \epsilon^{-1} \operatorname{Cov}\left(g\left(X, \xi_{0}+\epsilon t\right)-g\left(X, \xi_{0}\right), g\left(X, \xi_{0}\right)\right)=\alpha t
$$

for each $t \in \mathbb{R}$;
(v) there is a $\delta>0$ such that for any $\left|t-\xi_{0}\right|<\delta$ and for any combinations $i_{1}<\cdots<i_{m}$ and $j_{1}<\cdots<j_{m}$ such that $\left\{i_{1}, \ldots, i_{m}\right\} \neq\left\{j_{1}, \ldots, j_{m}\right\}$, we have that

$$
\operatorname{Pr}\left\{h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right)=h\left(X_{j_{1}}, \ldots, X_{j_{m}}\right)=t\right\}=0 .
$$

Then, with probability 1,

$$
\left\{(n / 2 \log \log n)\left(H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right)\right\}, n \geq 1,
$$

is relatively compact and its limit set is $\left\{\lambda_{1} x_{1}^{2}+\lambda_{2} x_{2}^{2}: x_{1}^{2}+x_{2}^{2} \leq 1\right\}$, where $\lambda_{1}$ and $\lambda_{2}$ are in (2.15) and (2.16).

Proof. Take $M>m\left(H^{\prime}\left(\xi_{0}\right)\right)^{-1}\left(\operatorname{Var}\left(g\left(X, \xi_{0}\right)\right)\right)^{1 / 2}$. Let

$$
\begin{aligned}
Z_{n}(t):=(n / 2 \log \log n)( & H_{n}\left(\xi_{0}+t(2 \log \log n / n)^{1 / 2}\right) \\
& \left.-H\left(\xi_{0}+t(2 \log \log n / n)^{1 / 2}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)\right)
\end{aligned}
$$

for $|t| \leq M$ and let $Z_{n}(\infty)$ as in (3.7). By the method in Theorem 8, we get that $\left\{Z_{n}(t): t \in T_{M}^{*}\right\}$ is a.s. relatively compact and its limit set is

$$
\begin{aligned}
\left\{(\gamma(t))_{t \in T_{M}^{*}}: \gamma(t)\right. & =t u_{1} \text { for }|t| \leq M \text { and } \\
\gamma(\infty) & \left.=u_{2}, \text { where }\left(u_{1}, u_{2}\right) \in K\right\},
\end{aligned}
$$

where $K$ is the unit ball of the reproducing kernel Hilbert space of the random vector ( $Y_{1}, Y_{2}$ ), where $Y_{1}$ and $Y_{2}$ are jointly normal random variables with mean zero and

$$
E\left[Y_{1}^{2}\right]=c_{1,1}, E\left[Y_{2}^{2}\right]=c_{2,2} \quad \text { and } \quad E\left[Y_{1} Y_{2}\right]=c_{1,2}
$$

where $c_{1,1}, c_{2,2}$ and $c_{1,2}$ are as in (2.17). By composition

$$
\left\{(n / 2 \log \log n)\left(H_{n}\left(\xi_{n}\right)-H\left(\xi_{n}\right)-H_{n}\left(\xi_{0}\right)+H\left(\xi_{0}\right)+b\left(\xi_{n}-\xi_{0}\right)^{2}\right)\right\}
$$

$$
n \geq 1,
$$

is a.s. relatively compact and its limit set is

$$
\left\{u_{1} u_{2}+b u_{2}^{2}:\left(u_{1}, u_{2}\right) \in K\right\} .
$$

Therefore,

$$
\left\{(n / 2 \log \log n)\left(H_{0}^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)+H\left(\xi_{n}\right)-H\left(\xi_{0}\right)\right)\right\}, n \geq 1,
$$

is a.s. relatively compact and its limit set is

$$
L:=\left\{-u_{1} u_{2}-b u_{2}^{2}:\left(u_{1}, u_{2}\right) \in K\right\} .
$$

Now,

$$
K=\left\{\left(E\left[Y_{1}\left(b_{1} Y_{1}+b_{2} Y_{2}\right)\right], E\left[Y_{2}\left(b_{1} Y_{1}+b_{2} Y_{2}\right)\right]\right):\left\|b_{1} Y_{1}+b_{2} Y_{2}\right\|_{2} \leq 1\right\} .
$$

If $c_{1,1} c_{2,2}-c_{1,2}^{2}=0$, then $Y_{1}=c_{1,1}^{1 / 2} g$ and $Y_{2}=\operatorname{sign}\left(c_{1,2}\right) c_{2,2}^{1 / 2} g$, where $g$ is a standard normal random variable. Thus,

$$
K=\left\{\left(E\left[Y_{1} x g\right], E\left[Y_{2} x g\right]\right)=\left\{\left(c_{1,1}^{1 / 2} x, \operatorname{sign}\left(c_{1,2}\right) c_{2,2}^{1 / 2} x\right): x^{2} \leq 1\right\}\right.
$$

and

$$
L=\left\{-\left(c_{1,2}+b c_{2,2}\right) x^{2}: x^{2} \leq 1\right\} .
$$

We have that

$$
\begin{aligned}
& \left\{\lambda_{1} x_{1}^{2}+\lambda_{2} x_{2}^{2}: x_{1}^{2}+x_{2}^{2} \leq 1\right\} \\
& \quad=\left\{\left(c_{1,2}+b c_{2,2}\right)^{-} x_{1}^{2}-\left(c_{1,2}+b c_{2,2}\right)^{+} x_{2}^{2}: x_{1}^{2}+x_{2}^{2} \leq 1\right\}=L
\end{aligned}
$$

where $\lambda_{1}$ and $\lambda_{2}$ are as in (2.15) and (2.16). So, the claim follows in this case. Assume that $c_{1,1} c_{2,2}-c_{1,2}^{2}>0$. Letting $Z_{1}=c_{2,2}^{-1 / 2}\left(c_{1,1} c_{2,2}-c_{1,2}^{2}\right)^{-1 / 2}\left(c_{2,2} \times\right.$ $Y_{1}-c_{1,2} Y_{2}$ ) and $Z_{2}=c_{2,2}^{-1 / 2} Y_{2}$, we have that $Z_{1}$ and $Z_{2}$ are independent standard normal random variables and

$$
\begin{aligned}
K & =\left\{\left(E\left[Y_{1}\left(d_{1} Z_{1}+d_{2} Z_{2}\right)\right], E\left[Y_{2}\left(d_{1} Z_{1}+d_{2} Z_{2}\right)\right]\right):\left\|d_{1} Z_{1}+d_{2} Z_{2}\right\|_{2} \leq 1\right\} . \\
& =\left\{\left(c_{2,2}^{-1 / 2}\left(c_{1,1} c_{2,2}-c_{1,2}^{2}\right)^{1 / 2} d_{1}+c_{2,2}^{-1 / 2} c_{1,2} d_{2} c_{2,2}^{1 / 2} d_{2}\right): d_{1}^{2}+d_{2}^{2} \leq 1\right\} .
\end{aligned}
$$

From this and the fact that

$$
\begin{aligned}
& -\left(c_{2,2}^{-1 / 2}\left(c_{1,1} c_{2,2}-c_{1,2}^{2}\right)^{1 / 2} d_{1}+c_{2,2}^{-1 / 2} c_{1,2} d_{2}\right) c_{2,2}^{1 / 2} d_{2}-b\left(c_{2,2}^{1 / 2} d_{2}\right)^{2} \\
& \quad=-a_{1} d_{1} d_{2}-a_{2} d_{2}^{2},
\end{aligned}
$$

where $a_{1}=\left(c_{1,1} c_{2,2}-c_{1,2}^{2}\right)^{1 / 2}$ and $a_{2}=c_{1,2}+b c_{2,2}$, we get that

$$
L=\left\{-a_{1} d_{1} d_{2}-a_{2} d_{2}^{2}: d_{1}^{2}+d_{2}^{2} \leq 1\right\} .
$$

If $x_{1}=q_{1}^{-1} a_{1} d_{1}+q_{1}^{-1}\left(a_{2}+\left(a_{1}^{2}+a_{2}^{2}\right)^{1 / 2}\right) d_{2}$ and $x_{2}=q_{2}^{-1} a_{1} d_{1}+q_{2}^{-1}\left(a_{2}-\right.$ $\left.\left(a_{1}^{2}+a_{2}^{2}\right)^{1 / 2}\right) d_{2}$, where $q_{1}$ and $q_{2}$ are as in (2.18), then

$$
-a_{1} d_{1} d_{2}-a_{2} d_{2}^{2}=\lambda_{1} x_{1}^{2}+\lambda_{2} x_{2}^{2}
$$

and

$$
d_{1}^{2}+d_{2}^{2}=x_{1}^{2}+x_{2}^{2} .
$$

Therefore,

$$
L=\left\{\lambda_{1} x_{1}^{2}+\lambda_{2} x_{2}^{2}: x_{1}^{2}+x_{2}^{2} \leq 1\right\}
$$

and the result follows.
Theorem 9 applies to kernels like $h\left(x_{1}, x_{2}\right)=x_{1}+x_{2}, h\left(x_{1}, x_{2}\right)=\mid x_{1}-$ $\left.x_{2}\right|^{r}, r>0$, and $h\left(\left(x_{1}, y_{1}\right),\left(x_{2}, y_{2}\right)\right)=\left(y_{2}-y_{1}\right) /\left(x_{2}-x_{1}\right)$, under some regularity conditions. We omit the details.

Using arguments similar to those used above, it is easy to obtain the following theorems.

Theorem 10. Let $\left\{X_{j}\right\}_{j=1}^{\infty}$ be a sequence of i.i.d. r.v.'s with values in a measurable space $(S, \mathscr{S})$. Let $(T, \rho)$ be a pseudometric space. Let $g: S \times T \rightarrow \mathbb{R}$ be a function such that $g(\cdot, t): S \rightarrow \mathbb{R}$ is a measurable function for each $t \in T$. Let $\left\{b_{n}\right\}$ be a sequence of real numbers from the interval $(0,1]$ and let $\left\{a_{n}\right\}$ be a sequence of positive real numbers. Suppose that:
(i) there is a scalar product defined for each $t \in T$ and each $0 \leq u \leq 1$, so that $u t \in T$;
(ii) $\lim \sup _{n \rightarrow \infty} \sup _{t \in T} a_{n}^{2} n^{-1} \operatorname{Var}\left(g\left(X, b_{n} t\right)\right)<\infty$;
(iii) $\left\{b_{n}\right\}$ and $\left\{a_{n} n^{-1}(\log \log n)^{-1 / 2}\right\}$ are nonincreasing sequences;
(iv) $\lim \sup _{n \rightarrow \infty} \sup _{n \leq m \leq 2 n} a_{n}^{-1} a_{m}<\infty$;
(v) there are positive constants $r_{1}$ and $r_{2}$, such that

$$
\sum_{j=2}^{\infty} \sup _{r_{1}(\log j)^{-1} \leq r \leq r_{2}} e^{-r_{1} r^{-1}} 2^{j} \operatorname{Pr}\left\{G\left(x, b_{2^{j}}\right) \geq r 2^{j}(\log j)^{1 / 2} a_{2_{j}}^{-1}\right\}<\infty,
$$

where $G(x):=\sup _{t \in T} \mid g\left(x, b_{n} t\right)-E\left[g\left(X, b_{n} t\right)\right]$;
(vi) $(2 \log \log n)^{-1 / 2} a_{n} \sup _{t \in T}\left|\left(P_{n}-P\right) g\left(\cdot, b_{n} t\right)\right|=O_{\mathrm{P}}(1)$.

Then there is a finite constant $c$ such that

$$
\limsup _{n \rightarrow \infty} \sup _{|t| \leq M}(2 \log \log n)^{-1 / 2} a_{n}\left|\left(P_{n}-P\right) g\left(\cdot, b_{n} t\right)\right| \leq c \quad \text { a.s. }
$$

Theorem 11. Let $h: S^{m} \rightarrow \mathbb{R}$ be a symmetric measurable function. Let $0<p<1$. Suppose that:
(i) there is a real number $\xi_{0}$ such that $H\left(\xi_{0}\right)=p, H$ is continuous in a neighborhood of $\xi_{0}, H$ is differentiable at $\xi_{0}, H^{\prime}\left(\xi_{0}\right)>0$ and there exists a constant b such that

$$
H(t)=H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(t-\xi_{0}\right)+b\left(t-\xi_{0}\right)^{2}+o\left(\left(t-\xi_{0}\right)^{2}\right)
$$

as $t \rightarrow \xi_{0}$;
(ii) there is a $1 \leq v \leq 2$ such that

$$
\limsup _{\varepsilon \rightarrow 0}|\varepsilon|^{-v} E\left[\left|g\left(X, \xi_{0}+\varepsilon\right)-g\left(X, \xi_{0}\right)\right|^{2}\right]<\infty .
$$

(iii) there is a $\delta>0$ such that for any $\left|t-\xi_{0}\right|<\delta$ and for any combinations $i_{1}<\cdots<i_{m}$ and $j_{1}<\cdots<j_{m}$ such that $\left\{i_{1}, \ldots, i_{m}\right\} \neq\left\{j_{1}, \ldots, j_{m}\right\}$, we have that

$$
\operatorname{Pr}\left\{h\left(X_{i_{1}}, \ldots, X_{i_{m}}\right)=h\left(X_{j_{1}}, \ldots, X_{j_{m}}\right)=t\right\}=0 .
$$

Then there is a finite constant $c$ such that
$\lim \sup (n / 2 \log \log n)^{(v+2) / 4}\left|H_{n}\left(\xi_{0}\right)-H\left(\xi_{0}\right)+H^{\prime}\left(\xi_{0}\right)\left(\xi_{n}-\xi_{0}\right)\right| \leq c \quad$ a.s. $n \rightarrow \infty$

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Department of Mathematics
University of Texas
Austin, Texas 78712 -1082
E-mail: arcones@math.utexas.edu


[^0]:    Received December 1993; revised July 1995
    ${ }^{1}$ Research supported in part by NSF Grant DMS-93-02583 and carried out at the University of Utah.

    AMS 1991 subject classifications. Primary 62E20; secondary 60F05, 60F15.
    Key words and phrases. Quantiles, Bahadur-Kiefer representations, $U$-statistics, empirical processes.

