# UNIVERSITY OF CALIFORNIA, SAN DIEGO 

## DEPARTMENT OF ECONOMICS

A SUBSAMPLING APPROACH TO ESTIMATING THE DISTRIBUTION OF DIVERGING STATISTICS WITH APPLICATIONS TO ASSESSING FINANCIAL MARKET RISKS

BY<br>PATRICE BERTAIL<br>CHRISTIAN HAEFKE<br>DIMITRIS N. POLITIS<br>AND<br>HALBERT WHITE

# A subsampling approach to estimating the distribution of diverging statistics with applications to assessing financial market risks 

Patrice Bertail ${ }^{*} \quad$ Christian Haefke ${ }^{\dagger} \quad$ Dimitris N. Politis ${ }^{\ddagger}$<br>Halbert White ${ }^{\dagger}$

January 2000


#### Abstract

In this paper we propose a subsampling estimator for the distribution of statistics diverging at either known or unknown rates when the underlying time series is strictly stationary and strong mixing. Based on our results we provide a detailed discussion how to estimate extreme order statistics with dependent data and present two applications to assessing financial market risk. Our method performs well in estimating Value at Risk and provides a superior alternative to Hill's estimator in operationalizing Safety First portfolio selection.


JEL Classification: C14, C49, G11
Keywords: Resampling Methods, Extreme Value Statistics, Value at Risk, Portfolio Selection.

[^0]
## 1 Introduction

Politis and Romano (1994) introduced the methodology of "subsampling" and showed that it leads to valid large-sample statistical inferences in general estimation situations and data structures (including i.i.d. and stationary data), provided that the employed estimator converges to the unknown parameter at a known rate $\tau_{n}$, and possesses (when normalized by the rate $\tau_{n}$ ) a nondegenerate asymptotic distribution. Later, Bertail et al. (1999) dispensed with the assumption that the convergence rate $\tau_{n}$ is known, by using a preliminary round of subsampling to consistently estimate the rate $\tau_{n}$.

In the present paper, we show how subsampling can be used to approximate the sampling distributions of diverging (as opposed to converging) statistics that are particularly useful in the context of financial time series and in particular in assessing financial risk. We consider several such statistics in detail as examples. In addition, we deal with the issues of unknown divergence rate and/or unknown rate of "escaping means".

To set up the context for our exposition, let $X^{n}=\left(X_{1}, \ldots, X_{n}\right)$ be an observable stretch of a (strictly) stationary time series $\left\{X_{t}, t \in Z\right\}$ which - for simplicity - will be assumed real-valued. The random variables $\left\{X_{t}, t \in Z\right\}$ are defined on a probability space denoted by $(P, \Omega, \mathcal{A})$, and the first marginal distribution of the $\left\{X_{t}, t \in Z\right\}$ sequence is denoted by $F$; in other words, $P\left(X_{1} \leq x\right)=F(x)$.

We assume that the time series satisfies a weak dependence condition. In particular, we impose the strong mixing condition, namely that $\alpha_{X}(k) \rightarrow 0$ as $k \rightarrow \infty$; the Rosenblatt strong mixing coefficients are defined as usual by

$$
\alpha_{X}(k)=\sup _{A, B}|P(A \cap B)-P(A) P(B)|,
$$

where $A$ and $B$ are events in the $\sigma$-algebras generated by $\left\{X_{t}, t<0\right\}$ and $\left\{X_{t}, t \geq k\right\}$ respectively; the case where $X_{1}, \ldots, X_{n}$ are independent, identically distributed (i.i.d.) is an important special case of the general scenario in which $\alpha_{X}(k)=0$ for all $k>0$.

A statistic $T_{n}=T_{n}\left(X^{n}\right)$ is computed from the data. If $T_{n}$ estimates an unknown parameter $\theta=\theta(P)$, and if $T_{n}$ is consistent for $\theta$, then the subsampling methodology of Politis and Romano (1994) immediately applies to provide a consistent approximation to the sampling distribution of $T_{n}$ under minimal assumptions. In addition, if the rate of convergence of $T_{n}$ to $\theta$ is a priori unknown, then it can be estimated via a preliminary round of subsampling and used in the sampling distribution approximation; see Bertail et al. (1999).

In contrast, we are concerned here with diverging (as opposed to converging) statistics, i.e., statistics for which ${ }^{1} T_{n} \xrightarrow{P} \infty$ as $n \rightarrow \infty$. In particular, we will

[^1]work under the assumption that there exists a (nondegenerate) distribution $K(x, P)$, continuous in $x$, such that
\[

$$
\begin{equation*}
K_{n}(x, P) \equiv P\left\{\tau_{n} T_{n} \leq x\right\} \rightarrow K(x, P) \tag{1}
\end{equation*}
$$

\]

as $n \rightarrow \infty$ for any real number $x$, where $\tau_{n}$ is a decreasing function of $n$. For example, it may be the case that $\tau_{n}=n^{-\gamma}$ where $\gamma$ is a positive constant. Examples of such diverging statistics with well-defined (after normalization) asymptotic distributions are given by the extreme order statistics, e.g. the maximum or the $k$ th largest value, from heavy-tailed observations; see Leadbetter et al. (1983) where an excellent treatment of extreme order statistics from stationary observations is given.

Our interest in such statistics is directly motivated by certain quantities relevant for the measurement of exposure to risk in asset markets. To describe these, let $\left\{P_{t}\right\}$ be a sequence of asset prices, such as the daily closing price for the S\&P 500 Index, and let $X_{t} \equiv \ln \left(P_{t} / P_{t-1}\right)$ be the one-period return.

Example 1a. The return on the asset from time period 0 to time period $n$ can be measured by

$$
\ln \left(P_{n} / P_{0}\right)=\sum_{t=1}^{n} X_{t}
$$

A simple form of the efficient markets theory implies that $X_{t}-\mu$ is a martingale difference sequence with respect to publicly available information, where $\mu$ is some constant; note that this is compatible with the assumption that $\left\{X_{t}\right\}$ is a stationary mixing process. In fact, many continuous-time stochastic models of asset price evolution used in the finance literature directly imply that $\left\{X_{t}\right\}$ is a stationary, strong mixing process - examples are the simple geometric Brownian motion a la Black-Scholes (1973) and the stochastic volatility models of Hull and White (1987). Taking $T_{n}=\ln \left(P_{n} / P_{0}\right)-n \mu$ then gives a simple example of a statistic for which (1) holds: standard central limit results ensure that $\tau_{n}=1 / \sqrt{n}$. Note, however, that our interest is directed toward the distribution of $T_{n}$, i.e., $K_{n}(x, P)$, rather than the limiting distribution $K(x, P)$. Significantly, this simple case plays a central role in financial risk assessment, through the notion of "value at risk" (VaR) (see for example Hendricks (1996) or Jorion (1997)). One leading measure of VaR is defined as the quantity $q$ such that the probability of the asset value falling below $q$ at the end of a specified time period (of duration $n$ ) is a given value $p$. In other words, the value at risk is simply the quantile

$$
q=K_{n}^{-1}(p, P)
$$

The Basle Committee on Banking Supervision has recently proposed the use of value at risk for assessing capital adequacy, while the Securities and Exchange
e.g. $O_{P}($.$) and o_{P}($.$) will also be used in the sequel -$ see Brockwell and Davis (1987) for definitions.

Commission (SEC) has recently required that all corporate treasurers report estimates of value at risk for assets held in corporate portfolios. Standard methods based on assumptions of normality may deliver rather poor estimates of the desired quantiles, due to the well-known excess kurtosis typical of financial market returns. The methods discussed here provide a convenient approach potentially more accurate than current methods (see e.g. Hendricks (1996)).

Example 1b. Another measure of value at risk that is computable without any knowledge of the centering factor $n \mu$ is based on the quantiles of $U_{n}=$ $\ln \left(P_{n} / P_{0}\right)$. A simple rescaling generally will not ensure that equation (1) holds with $U_{n}$ instead of $T_{n}$. We therefore have to take into account a centering factor which in general may take the form of a sequence $\mu_{n}$; in other words, starting with a general statistic $U_{n}$, we define $T_{n}=U_{n}-\mu_{n}$, where $\mu_{n}$ is an appropriately chosen sequence such that (1) holds as stated. This is actually a most general set-up, and our statistic $U_{n}$ will diverge to infinity if the sequence $\tau_{n}$ is decreasing (the "inflating variance case"), and/or if the product $\tau_{n} \mu_{n}$ is increasing (the "escaping mean case"). Our section 3 is devoted to this general set-up. ${ }^{2}$

Example 2. Next, consider the worst cumulative return that an asset yields at any point during a fixed time horizon of duration $n$,

$$
T_{n}=\min _{1 \leq t \leq n} R_{t}, \quad R_{t}=\sum_{\tau=1}^{t} X_{\tau}
$$

This quantity determines whether "stop" orders are hit or bankruptcy occurs during the horizon and also determines whether "knock-out" or "knock-in" provisions of certain exotic but increasingly popular options are activated. The distribution of $T_{n}$ gives the probability of occurrence of such events. Further, the quantiles of $T_{n}$ provide another useful measure of risk exposure, which we shall call "extreme value at risk", (XVaR) as these give the quantities $q$ such that a loss will exceed $q$ at any time during the specified period with specified probability. Analytic analysis of the distribution of $T_{n}$ is not simple. Nevertheless, our methods provide a simple and direct way to consistently estimate the distribution or quantiles for this choice of $T_{n}$.

Example 3. Our third example comes from the literature of "safety first" portfolio selection (Roy, 1952; Arzac and Bawa, 1977; Jansen, et al., 1998).

[^2]In this theory, it is of interest to know the distribution of the worst possible one-period return over a given time horizon,

$$
T_{n}=\max \left(-X_{1}, \ldots,-X_{n}\right)
$$

This quantity is also the focus of "worst case scenario" analysis, an alternative to value at risk discussed by Boudoukh, Richardson and Whitelaw (1995).

A standard regularity condition on the tail of the distribution $F$ of the $X_{t} \mathrm{~s}$ is that, for some $\alpha>0$ measuring tail thickness and a finite constant $\Gamma$, we have

$$
x^{\alpha}(1-F(x)) \rightarrow \Gamma \quad \text { as } x \rightarrow \infty .
$$

Under an assumption of strong mixing, it can be shown in this case that $\tau_{n}$ is proportional to $n^{-1 / \alpha}$; Leadbetter et al. (1983). Our methods will permit us to estimate not only the distribution $K_{n}$ and quantiles $K_{n}^{-1}$, but also the tail parameters $\alpha$ and $\Gamma$.

Without limiting ourselves to a particular example, we work out in Section 2 consistent approximations to the decreasing sequence $\tau_{n}$ and to $K_{n}(x, P)$ and $K(x, P)$ via subsampling in the case where the centering sequence $\mu_{n}$ is negligible, i.e. where $\mu_{n}=0$ or at least $\tau_{n} \mu_{n} \rightarrow 0$. We further show that the quantiles of $K_{n}(x, P)$ and $K(x, P)$ can also be estimated consistently; as a result, prediction intervals for $T_{n}$ can be formulated. We also establish that, if the objective is estimation of the sampling distribution of $T_{b}$ (with $b$ fixed), then a simple subsampling estimator may be applied, and a reference to the large-sample distribution $K(x, P)$ may be unnecessary even if the form of $K(x, P)$ is explicitly known. In section 3 , we consider the general case when the centering sequence $\mu_{n}$ is unknown and nonnegligible, and we give convergent estimators of $\tau_{n}$ and $\mu_{n}$. Our results are illustrated by the case of the maximum when the domain of attraction of the underlying distribution is totally unknown. In Section 4 we focus attention on $T_{n}$ being the maximum of heavy-tailed stationary observations. In that case, it is well-known that $K(x, P)$ has a specific form (the type II extreme distribution - see Leadbetter et al. (1983)) which depends on $P$ only through $F$, and in particular only on the tail of $F$. We show how subsampling can be successfully used to give a valid estimate of the 'thickness' parameter $\alpha$ of the tail of $F$; consequently, the subsampling estimate of the distribution of $T_{n}$ and/or the limit distribution $K(x, P)$ (with estimated 'thickness' parameter) can be inverted for use in the problem of safety-first portfolio selection. In Section 5 we illustrate applications of our subsampling methodology by obtaining estimates of various quantities relevant for measuring exposure to financial risk, including value at risk and safety first portfolio selection. All technical proofs are placed in the Appendix.

## 2 Subsampling for diverging statistics: the inflating variance case

### 2.1 Some definitions and a first result

Define $Y_{i}$ to be the subsequence $\left(X_{i}, X_{i+1}, \ldots, X_{i+b_{n}-1}\right)$, for $i=1, \ldots, q$, and $q=n-b_{n}+1$; note that $Y_{i}$ consists of $b_{n}$ consecutive observations from the $X_{1}, \ldots, X_{n}$ sequence, and the order of the observations is preserved. Now let $T_{b_{n}, i}$ be the value of the statistic $T_{b}$ applied to the subsample $Y_{i}$. The subsampling distribution of $\tau_{n} T_{n}$, based on a subsample size $b_{n}$, is defined by

$$
\begin{equation*}
K_{b_{n}}\left(x \mid X^{n}, \tau_{.}\right) \equiv q^{-1} \sum_{i=1}^{q} 1\left\{\tau_{b_{n}} T_{b_{n}, i} \leq x\right\} \tag{2}
\end{equation*}
$$

where $1\{A\}$ is the indicator of set $A$; note that the subsample size $b_{n}$ is allowed to depend on the actual sample size $n$.

The following theorem shows that $K_{b_{n}}\left(x \mid X^{n}, \tau\right.$.) is a consistent distribution estimator; it complements the results of Politis and Romano (1994) by dealing with the case of possibly diverging statistics, i.e., where $\tau_{n}$ may be decreasing.

Theorem 1 Let the $X$-sequence be stationary and strong mixing. Let $\tau_{n}$ be a given known function of $n$.
(a) If $b_{n}=b$ (a constant), then, as $n \rightarrow \infty$,

$$
K_{b_{n}}\left(x \mid X^{n}, \tau_{.}\right)=K_{b}(x, P)+o_{P}(1)
$$

for all $x$.
(b) Suppose the convergence (1) is true for $\tau_{n}$ a known function of $n$ and that $b_{n} \rightarrow \infty$ as $n \rightarrow \infty$, but with $b_{n} / n \rightarrow 0$; then

$$
\begin{equation*}
\sup _{x}\left|K_{b_{n}}\left(x \mid X^{n}, \tau\right)-K(x, P)\right|=o_{P}(1) \tag{3}
\end{equation*}
$$

and

$$
\begin{equation*}
\sup _{x}\left|K_{b_{n}}\left(x \mid X^{n}, \tau_{.}\right)-K_{d_{n}}(x, P)\right|=o_{P}(1) \tag{4}
\end{equation*}
$$

where $d_{n}$ is any sequence such that $d_{n} \rightarrow \infty$.
Part (a) of the Theorem indicates that, if our objective is to estimate $K_{b}(x, P)$ for some fixed $b$, then this can be achieved in a consistent manner using $K_{b_{n}}\left(x \mid X^{n}, \tau\right.$.) with $b_{n}=b$. If our objective is to estimate $K_{d_{n}}(x, P)$ for some $d_{n}$ such that $d_{n} \rightarrow \infty$ (e.g. $d_{n}=n$ ), then this can be accomplished using $K_{b_{n}}\left(x \mid X^{n}, \tau_{\text {. }}\right)$ with $b_{n} \rightarrow \infty$ (but always such that $\left.b_{n} / n \rightarrow 0\right)$.

Remark 1.1 Note that assumption (1) of part (b) the Theorem concerns the existence of a (continuous) limit distribution; the shape of this limit distribution may well be unknown, and can be estimated by (3) above.

Remark 1.2 Looking at the proof of part (a) of the Theorem, note first that for the asymptotic approximation to be useful in a finite-sample situation, $b$ not only has to be fixed, but it has to be small as compared to $n$ (so that $b / n$ is small). Also note that for part (a), existence of a limit distribution (as in (1)) is not required; thus, part (a) is true under very few assumptions. Furthermore, it turns out that we do not even need to consider a divergence rate $\tau_{n}$ in order to use a result similar to part (a); this is addressed by the following corollary for which the following definitions are required. Let

$$
\bar{K}_{b_{n}}\left(x \mid X^{n}\right) \equiv K_{b_{n}}\left(x \mid X^{n}, 1\right)
$$

and

$$
\bar{K}_{b}(x, P) \equiv P\left\{T_{b} \leq x\right\}
$$

Corollary 1 Let the $X$-sequence be stationary and strong mixing. If $b_{n}=b$ ( $a$ constant), then, as $n \rightarrow \infty$,

$$
\bar{K}_{b}\left(x \mid X^{n}\right)=\bar{K}_{b}(x, P)+o_{P}(1)
$$

for all $x$. In addition, if $\bar{K}_{b}(x, P)$ is continuous in $x$, we have

$$
\sup _{x}\left|\bar{K}_{b}\left(x \mid X^{n}\right)-\bar{K}_{b}(x, P)\right|=o_{P}(1)
$$

Another immediate corollary (in the fixed $b$ case) pertains to estimation of the quantiles of $\bar{K}_{b}(x, P)$ relevant for setting prediction intervals for 'future' observations of $T_{b}$. Given a distribution $G$ on the real line and a number $t \in$ $(0,1)$, we let $G^{-1}(t)$ denote the quantile transformation, i.e., $G^{-1}(t)=\inf \{x$ : $G(x) \geq t\}$, which reduces to the regular inverse of the function $G$ if $G$ happens to be invertible.

Corollary 2 Let the $X$-sequence be stationary and strong mixing. If $b_{n}=b$ ( $a$ constant), and if $\bar{K}_{b}(x, P)$ is continuous in $x$, then, as $n \rightarrow \infty$,

$$
\begin{equation*}
\bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=\bar{K}_{b}^{-1}(t, P)+o_{P}(1) \tag{5}
\end{equation*}
$$

for any $t \in(0,1)$.

### 2.2 Case of unknown divergence rate $\tau_{n}$

In Theorem 1 we treated the case of known divergence rate $\tau_{n}$. Nevertheless, in many cases $\tau_{n}$ is unknown, its form depending on the (unknown) probability measure $P$. It is often the case that the assumption $\tau_{n}=n^{-\gamma}$ is reasonable, but then typically the positive constant $\gamma$ is unknown as it may depend crucially on
$P$. In the case of converging statistics (i.e., $\tau_{n}=n^{+\gamma}$ ), estimation of the rate was treated in Bertail et al. (1999). We will now show how similar ideas are applicable also in the case of diverging statistics (i.e., $\tau_{n}=n^{-\gamma}$ ).

The following Lemma relates the quantiles of $\bar{K}_{b_{n}}\left(x \mid X^{n}\right)$ to those of $K(x, P)$ through the (unknown) rate $\tau_{n}$; it will be the first step towards constructing an estimator of the rate as in Bertail et al. (1999).

Lemma 1 Let $k_{0}=\sup \{x: K(x, P)=0\}$ and $k_{1}=\inf \{x: K(x, P)=1\}$, and assume that $K(x, P)$ is continuous and strictly increasing on $\left(k_{0}, k_{1}\right)$ as a function of $x$. If (3) is true as $n$ tends to infinity, then

$$
\begin{equation*}
\tau_{b_{n}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=K^{-1}(t, P)+o_{P}(1) \tag{6}
\end{equation*}
$$

for any $t \in(0,1)$.
For example, if $t=1 / 2$, then (6) simply says that $b_{n}^{-\gamma} \operatorname{median}\left\{T_{b_{n}, i}, i=\right.$ $1, \ldots, q\} \xrightarrow{P}$ median of $K(x, P)$, where median $\left\{T_{b_{n}, i}, i=1, \ldots, q\right\}$ is the sample median of the collection $\left\{T_{b_{n}, i}, i=1, \ldots, q\right\}$, and $\gamma$ is suitably chosen.

Therefore, if we look at the sample median of the $\left\{T_{b_{n}, i}, i=1, \ldots, q\right\}$ and/or other sample quantiles (i.e. different choices of $t$ in (6)) for different choices of $b_{n}$ (with all choices satisfying $b_{n} \rightarrow \infty$ and $b_{n} / n \rightarrow 0$ ), then we can estimate the constant $\gamma$ by a simple regression as discussed next. Having estimated $\gamma$, an estimate of the limit quantiles $K^{-1}(t, P)$ are also immediately available through (6).

To see this, note that by taking logarithms in (6) and assuming that $\tau_{n}=n^{-\gamma}$ we get

$$
\begin{equation*}
\log \left(\left|\bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)\right|\right)=\log \left(\left|K^{-1}(t, P)\right|\right)+\gamma \log b_{n}+o_{P}(1) \tag{7}
\end{equation*}
$$

so that $\gamma$ is simply the 'slope' in regressing $\log \left(\left|\bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)\right|\right)$ on $\log b_{n}$. So if we consider different subsample sizes $b_{i, n}, i=1, \ldots, I>1$, we can use the least squares estimator of slope, namely

$$
\begin{equation*}
\gamma_{I} \equiv \frac{\sum_{i=1}^{I}\left(y_{i}-\bar{y}\right)\left(\log \left(b_{i, n}\right)-\overline{\log }\right)}{\sum_{i=1}^{I}\left(\log \left(b_{i, n}\right)-\overline{\log }\right)^{2}} \tag{8}
\end{equation*}
$$

where $y_{i}=\log \left(\left|\bar{K}_{b_{i, n}}^{-1}\left(t \mid X^{n}\right)\right|\right), \bar{y}=I^{-1} \sum_{i=1}^{I} y_{i}$ and $\overline{\log }=I^{-1} \sum_{i=1}^{I} \log \left(b_{i, n}\right)$. Lemma 2 below establishes the consistency of $\gamma_{I}$ following similar results of Bertail et al. (1999).

Lemma 2 Let the $X$-sequence be stationary and strong mixing. Suppose (1) is true, and $\tau_{n}=n^{-\gamma}$ where $\gamma$ is a positive unknown constant. Let $K(x, P)$ be continuous and strictly increasing on $\left(k_{0}, k_{1}\right)$ as a function of $x$, where $k_{0}=$ $\sup \{x: K(x, P)=0\}$ and $k_{1}=\inf \{x: K(x, P)=1\}$. Let $b_{i, n}=n^{\beta_{i}}, 1>\beta_{1}>$ $\cdots>\beta_{I}>0$, and let $t$ be a point in $(0,1)$. Then $\gamma_{I}=\gamma+o_{P}\left((\log n)^{-1}\right)$.

As a remark, note that instead of looking at individual quantiles we could use differences of quantiles, e.g. the inter-quartile range, etc., in order to estimate $\gamma$; this is because (6) implies that, if $t_{1}<t_{2}$,

$$
\tau_{b_{n}}\left(\bar{K}_{b_{n}}^{-1}\left(t_{2} \mid X^{n}\right)-\bar{K}_{b_{n}}^{-1}\left(t_{1} \mid X^{n}\right)\right)=K^{-1}\left(t_{2}, P\right)-K^{-1}\left(t_{1}, P\right)+o_{P}(1)
$$

and thus
$\log \left(\bar{K}_{b_{n}}^{-1}\left(t_{2} \mid X^{n}\right)-\bar{K}_{b_{n}}^{-1}\left(t_{1} \mid X^{n}\right)\right)=\log \left(K^{-1}\left(t_{2}, P\right)-K^{-1}\left(t_{1}, P\right)\right)+\gamma \log b_{n}+o_{P}(1)$
as well. Bertail et al. (1999) originally developed this idea of looking at ranges, and we will take it up again in section 3 .

Lemma 3 below is an immediate corollary of Lemmas 1 and 2; its essence is that the quantiles $K^{-1}(x, P)$ (that represent the 'constant' term in the regression (7)) can be consistently estimated as well, using our estimator of $\gamma$ in place of the true one, or -for that matter- using any other (at least as accurate) estimator of $\gamma$.

Lemma 3 Let $\hat{\gamma}$ be a statistic calculated from the data $X^{n}$ such that $\hat{\gamma}=\gamma+$ $o_{P}\left((\log n)^{-1}\right)$. Under the assumptions of Lemma 1, and assuming that $\tau_{n}=n^{-\gamma}$ where $\gamma$ is a positive unknown constant, we have that

$$
\hat{\tau}_{b_{n}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=K^{-1}(t, P)+o_{P}(1)
$$

where $\hat{\tau}_{b_{n}}=b_{n}^{-\hat{\gamma}}$.
Lemma 2 can be generalized to take into account many different quantiles (i.e., $t$-points) at the same time. So let some points $t_{j} \in(0,1)$ be given, where $j=1, \ldots, J \geq 1$, together with subsample sizes $b_{i, n}, i=1, . ., I>1$. Equation (7) yields

$$
\begin{equation*}
y_{i, j} \equiv \log \left(\left|\bar{K}_{b_{i, n}}^{-1}\left(t_{j} \mid X^{n}\right)\right|\right)=a_{j}+\gamma \log \left(b_{i, n}\right)+u_{i, j} \tag{9}
\end{equation*}
$$

where $a_{j} \equiv \log \left(\left|K^{-1}\left(t_{j}, P\right)\right|\right)$ and $u_{i, j}=o_{P}(1), i=1, \ldots, I$ and $j=1, \ldots, J$.
Following Bertail et al. (1999) we suggest the following ANOVA-type estimator of $\gamma$ :

$$
\begin{equation*}
\gamma_{I, J} \equiv \frac{\sum_{i=1}^{I}\left(y_{i, .}-\bar{y}\right)\left(\log \left(b_{i, n}\right)-\overline{\log }\right)}{\sum_{i=1}^{I}\left(\log \left(b_{i, n}\right)-\overline{\log }\right)^{2}} \tag{10}
\end{equation*}
$$

where $y_{i, .}=J^{-1} \sum_{j=1}^{J} y_{i, j}=J^{-1} \sum_{j=1}^{J} \log \left(\left|\bar{K}_{b_{i, n}}^{-1}\left(t_{j} \mid X^{n}\right)\right|\right)$,
$\bar{y}=(I J)^{-1} \sum_{i=1}^{I} \sum_{j=1}^{J} y_{i, j}$ and $\overline{\log }=I^{-1} \sum_{i=1}^{I} \log \left(b_{i, n}\right)$. Theorem 2 now offers a generalization of Lemma 2.

Theorem 2 Let the $X$-sequence be stationary and strong mixing. Suppose (1) is true, and $\tau_{n}=n^{-\gamma}$ where $\gamma$ is a positive unknown constant. Let $K(x, P)$
be continuous and strictly increasing on $\left(k_{0}, k_{1}\right)$ as a function of $x$, where $k_{0}=$ $\sup \{x: K(x, P)=0\}$ and $k_{1}=\inf \{x: K(x, P)=1\}$. Let $b_{i, n}=n^{\beta_{i}}, 1>$ $\beta_{1}>\cdots>\beta_{I}>0$, and let points $t_{j} \in(0,1), j=1, \ldots J \geq 1$ be given. Then $\gamma_{I, J}=\gamma+o_{P}\left((\log n)^{-1}\right)$.

What is perhaps more important than rate estimation per se, is that the estimated divergence rate can be used in turn in order to achieve consistent estimation of $K(x, P)$ and its quantiles. As a matter of fact, any estimator $\hat{\gamma}$ such that $\hat{\gamma}=\gamma+o_{P}\left((\log n)^{-1}\right)$ is accurate enough to be used in estimation of $K(x, P)$.

Theorem 3 Let the $X$-sequence be stationary and strong mixing. Suppose (1) is true, and $\tau_{n}=n^{-\gamma}$ where $\gamma$ is a positive unknown constant. Let $K(x, P)$ be continuous and strictly increasing on $\left(k_{0}, k_{1}\right)$ as a function of $x$, where $k_{0}=$ $\sup \{x: K(x, P)=0\}$ and $k_{1}=\inf \{x: K(x, P)=1\}$. Let $\hat{\gamma}$ be a statistic calculated from the data $X^{n}$ and such that $\hat{\gamma}=\gamma+o_{P}\left((\log n)^{-1}\right)$; also let $\widehat{\tau}_{n}=n^{-\hat{\gamma}}$. Let $b_{n}$ be such that $b_{n} \rightarrow \infty$ as $n \rightarrow \infty$, but $b_{n} / n \rightarrow 0$. Then

$$
\begin{equation*}
\sup _{x}\left|K_{b_{n}}\left(x \mid X^{n}, \widehat{\tau}\right)-K(x, P)\right|=o_{P}(1) . \tag{11}
\end{equation*}
$$

Let $t \in(0,1)$, and let $c_{n}(t)=K_{b_{n}}^{-1}\left(t \mid X^{n}, \widehat{\tau}.\right)$ be the $t^{t h}$ quantile of the subsampling distribution $K_{b_{n}}\left(x \mid X^{n}, \widehat{\tau}.\right)$. Then

$$
\begin{equation*}
P\left\{\hat{\tau}_{d_{n}} T_{d_{n}} \leq c_{n}(t)\right\} \longrightarrow t \tag{12}
\end{equation*}
$$

where $d_{n}$ is any sequence such that $d_{n} \rightarrow \infty$, i.e., $c_{n}(t)$ consistently estimates the $t$-quantile of $\hat{\tau}_{d_{n}} T_{d_{n}}$.

Remark 3.1. Theorem 3 is a generalization of part (b) of Theorem 1 in the case where the divergence rate $\tau_{n}$ is unknown and must be estimated. Note in particular that equation (12) may be used for the construction of prediction intervals for future observations of $T_{d_{n}}$.

## 3 Subsampling for diverging statistics: general case

As mentioned in Example 1b of the Introduction, the divergence of a general statistic $U_{n}$ may be due either to an inflating variance and/or an escaping mean; this very general case is addressed in the present section. We assume that $U_{n}$ is a computable statistic such that, as $n \rightarrow \infty$ for any real number $x$,

$$
\begin{equation*}
K_{n}(x, P) \equiv P\left\{\tau_{n}\left(U_{n}-\mu_{n}\right) \leq x\right\} \longrightarrow K(x, P) \tag{13}
\end{equation*}
$$

note that equation (13) is identical to (1) with the definition $T_{n}=U_{n}-\mu_{n}$.

Estimating simultaneously the divergence rate $\tau_{n}$ and the centering factor $\mu_{n}$ (or at least its dominant part) is the subject of Section 3.2 below; we also achieve the construction of consistent subsampling distribution estimators. Note that stated in full generality, the problem of estimating $\mu_{n}$ and $\tau_{n}$ is not identifiable: indeed if $\tau_{n} \mu_{n}$ converges to some constant $C \neq 0$, then $P\left(\tau_{n} U_{n} \leq x\right) \longrightarrow K^{\prime}(x, P)=K(x-C, P)$, i.e., we also have an acceptable limiting distribution with the choice $\mu_{n}^{\prime}=0$. To make the problem identifiable, we may -without loss of generality- assume that $\tau_{n} \mu_{n}$ is either a strictly monotone function or identically zero (that is $\mu_{n}=0$ ); otherwise, we may assume that the median of $K(x, P)$ is known up to some constant. In the same way the rate of convergence can not be identified up to a multiplicative constant. Note however that these constants have no importance in the construction of confidence intervals because of the invariance properties of subsampling distributions.

### 3.1 Centering and normalization of the maximum

Before proceeding with our theoretical results we now elaborate on our Example 1 of the Introduction for further motivation; the elaboration consists in allowing for arbitrary distributions generating the data.

To fix ideas we now take $U_{n}=\max \left\{X_{1}, \ldots, X_{n}\right\}$, with $X_{i}$ being i.i.d. from distribution $F$; the assumption of independence is not crucial and can be replaced with an assumption of stationarity and weak dependence (mixing) -see Leadbetter et al. (1983). Then there are only three possible shapes for the limiting distribution $K(x, P)$ of the appropriately centered and standardized extreme value statistic $U_{n}$; the three cases are:

$$
\begin{gathered}
\text { Type } I: \quad K(x, P)=\Lambda(x) \equiv \exp (-\exp (-x)) \\
\text { Type } I I: \quad K(x, P)=\Phi_{\alpha}(x) \equiv \begin{cases}0 & \text { if } x \leq 0 \\
\exp \left(-x^{-\alpha}\right) & \text { if } x>0\end{cases} \\
\text { Type } I I I: \quad K(x, P)=\Psi_{\alpha}(x) \equiv \begin{cases}\exp \left(-(-x)^{\alpha}\right) & \text { if } x \leq 0 \\
1 & \text { if } x>0\end{cases}
\end{gathered}
$$

where $\alpha$ is some positive constant less than two. Which of the three cases obtains depends on the generating distribution $F$. We say that $F$ "falls in the domain of attraction" of distribution $G$ (denoted by $F \in D(G)$ ) if the limiting distribution $K(x, P)$ of the appropriately centered and standardized extreme value statistic $U_{n}$ turns out to be $G$.

Characterizing the domain of attraction of $F$ according to its tail as well as the norming and centering constants $\tau_{n}$ and $\mu_{n}$ has received a lot of attention; let us recall the main facts. Denote by $R_{\alpha}$ the set of regularly varying functions
at $\infty$ with index $\alpha$, that is, the set of measurable positive real functions $L(\cdot)$ such that

$$
\lim _{x \rightarrow \infty} L(\lambda x) / L(x)=\lambda^{\alpha},
$$

and define the "quantile" inverse of function $H(\cdot)$ by $H^{-1}(x)=\inf \{y: H(y) \geqslant$ $x\}$. We have the following equivalence (see e.g. Bingham et al. (1987, chap. 8)):

- $F \in D(\Lambda) \Leftrightarrow H^{-1}(x+y)-H(x) \sim y l\left(e^{x}\right)$ when $x \rightarrow \infty$, where $l \in R_{0}$ and we put $H(x) \equiv-\log (1-F(x))$. In that case, one may choose $\mu_{n}=$ $H^{-1}(\log n)$ and $\tau_{n}=\frac{1}{H^{-1}(\log n+1)-H^{-1}(\log n)}$. It is clear in this case that both the centering factor $\mu_{n}$ and the rate $\tau_{n}$ contribute to the divergence of $U_{n}$, and they both depend on the underlying distribution $F$.
- $F \in D\left(\Phi_{\alpha}\right) \Leftrightarrow 1-F(x) \in R_{-\alpha}$, ; in this case one may choose $\mu_{n}=$ 0 and $\tau_{n}=\frac{1}{(1-F)^{-1}\left(\frac{1}{n}\right)}$. This is the heavy-tailed case of Example 1 of the Introduction, which will be revisited in section 4.2 as well.
- $F \in D\left(\Psi_{\alpha}\right) \Leftrightarrow$ The upper endpoint (say $x_{0}$ ) of $F$ is finite and $1-F\left(x_{0}-\right.$ $\left.\frac{1}{x}\right) \in R_{-\alpha}$; in this case one may choose $\mu_{n}=x_{0}$ and $\tau_{n}=\frac{1}{F_{1}^{-1}\left(\frac{1}{n}\right)}$, where $F_{1}(\cdot)=1-F\left(x_{0}-\cdot\right)$. Note that here $T_{n}$ is actually a converging statistic and that $\mu_{n}$ is not escaping. Thus the results of Bertail et al. (1999) may be applied directly in this case. Recentering the subsampling distribution at $T_{n}$ would yield a convergent estimator of the true distribution of the maximum. Nevertheless, we show below how this case may be treated simultaneously with the others in a single unifying framework.

Typically, practitioners make some assumptions on the behavior of the tail of $F$ so that the domain of attraction is known and the problem mainly reduces to the estimation of the index $\alpha$ in case (i) and (ii) and to the characterization of the norming and centering constants in case (iii). A general theory that allows the estimation of the distribution of the maximum and the estimation of the norming and centering sequences without this kind of assumption is of interest from a robustness point of view. Moreover, even in the case (iii), estimating $\tau_{n}$ and $\mu_{n}$ simultaneously is a difficult problem. In the following we show how the subsampling ideas presented in section 2 may be successfully applied in this general case as well.

### 3.2 Estimation of the centering and norming sequences

It is clear that by putting $T_{n}=U_{n}-\mu_{n}$, results obtained in Theorems 1 and 2 and the associated lemmas also hold for $U_{n}$. In particular we have the following straightforward result relating the quantile of the subsampling distribution of $U_{n}$ to the norming and centering factors; its proof is similar to the proof of Lemma 1.

Lemma 4 Assume that the conditions of Lemma 1 hold for $T_{n}=U_{n}-\mu_{n}$ and let

$$
\bar{K}_{U, b_{n}}\left(x \mid X^{n}\right)=q^{-1} \sum_{i=1}^{q} 1\left\{U_{b_{n, i}} \leq x\right\}
$$

denote the unnormalized subsampling distribution of $U_{n}$, where $U_{b_{n, i}}$ is the value of the statistic $U_{b_{n}}$ as computed from subsample $Y_{i}$. Then we have

$$
\begin{equation*}
\bar{K}_{U, b_{n}}^{-1}\left(t \mid X^{n}\right)=\mu_{b_{n}}+\tau_{b_{n}}^{-1} K^{-1}(t, P)+o_{P}\left(\tau_{b_{n}}^{-1}\right) \tag{14}
\end{equation*}
$$

It follows from this lemma that for any point $0<t_{2}<t_{1}<1$ we have

$$
\begin{equation*}
\left|\bar{K}_{U, b_{n}}^{-1}\left(t_{2} \mid X^{n}\right)-\bar{K}_{U, b_{n}}^{-1}\left(t_{1} \mid X^{n}\right)\right| /\left|K^{-1}\left(t_{2}, P\right)-K^{-1}\left(t_{1}, P\right)\right|=\tau_{b_{n}}\left(1+o_{P}(1)\right) \tag{15}
\end{equation*}
$$

Suppose first that we are just interested in estimating the shape of $\tau_{n}$ without assuming that $\tau_{n}=n^{\gamma}$ as in the previous section. Then, since we can not hope to estimate $\tau_{n}$ but up to a multiplicative constant, we can simply choose $\widehat{\tau}_{b_{n}}=\left|\bar{K}_{U, b_{n}}^{-1}\left(t_{2} \mid X^{n}\right)-\bar{K}_{U, b_{n}}^{-1}\left(t_{1} \mid X^{n}\right)\right|$ as our rate estimator. In other words, any range of the subsampling distribution may be used to normalize the subsampling distribution which is a remarkable fact.

Because $\mu_{n}$ is a centering factor, information on the median of the asymptotic distribution may be of some interest. For instance, in case (iii) of section 3.1, the limiting distribution is parameter free with a median equal to $-\log \log 2$. Evaluating equation (14) at the median yields the equation

$$
\tau_{b_{n}} \bar{K}_{U, b_{n}}^{-1}\left(1 / 2 \mid X^{n}\right)=\tau_{b_{n}} \mu_{b_{n}}-\log \log 2+o_{P}(1)
$$

Since an obvious estimator of $\tau_{b_{n}}$ is available, it is not difficult to see that a twostep procedure (i.e. plugging the previous estimator $\widehat{\tau}_{b_{n}}$ in the above equation) yields a convergent estimator of the function $\mu_{n}$. Without loss of generality, we might have imposed the assumption that the median of the limiting distribution is zero. This would simply have shifted both the centering factor and the limiting distribution. With this identifying condition, the median of the subsampling distribution is actually the adequate recentering factor. The above considerations suggest the following general result.

Theorem 4 Assume that there exist sequences $\tau_{n}$ and $\mu_{n}$ such that

$$
K_{n}(x, P) \equiv P\left\{\tau_{n}\left(U_{n}-\mu_{n}\right) \leq x\right\} \longrightarrow K(x, P)
$$

where the asymptotic distribution $K(x, P)$ may be of unknown shape but is without loss of generality- assumed to be normalized such that $K^{-1}(1 / 2, P)=$ 0 and $\left|K^{-1}\left(t_{2}, P\right)-K^{-1}\left(t_{1}, P\right)\right|=1$ at some given points $0<t_{1}<t_{2}<1$.

Let $b_{n}$ be a subsampling size satisfying the assumptions of Theorem 1(b), and consider the median of the undersampling distribution

$$
\widehat{\mu}_{b_{n}}=\bar{K}_{U, b_{n}}^{-1}\left(1 / 2 \mid X^{n}\right)
$$

Then we have $\widehat{\mu}_{b_{n}}=\mu_{b_{n}}+o_{P}\left(\tau_{b_{n}}^{-1}\right)$ as $n \rightarrow \infty$.
Moreover

$$
\widehat{\tau}_{b_{n}}=\left|\bar{K}_{U, b_{n}}^{-1}\left(t_{2} \mid X^{n}\right)-\bar{K}_{U, b_{n}}^{-1}\left(t_{1} \mid X^{n}\right)\right|=\tau_{b_{n}}\left(1+o_{P}(1)\right)
$$

and we have

$$
q^{-1} \sum_{i=1}^{q} 1\left\{\widehat{\tau}_{b_{n}}\left(U_{b_{n, i}}-\widehat{\mu}_{b_{n}}\right) \leq x\right\} \rightarrow K(x, P)
$$

The proof of the theorem follows directly from (14) and (15) and is omitted. Of course, other identifiability conditions may be used to obtain the form of the centering and norming constant. In some cases (for instance in case (iii) of section 3.1), the asymptotic distribution is already "normalized" and entirely known. In that case equation (14) gives straightforward estimators of the normalizing and standardizing constants.

Theorem 5 Assume that $\tau_{n} \rightarrow 0$ as $n \rightarrow \infty$, that $b_{n}$ is a subsampling size satisfying the assumptions of Theorem 1(b), and that the asymptotic distribution $K(t, P)=K(t)$ is known. For any points $\left(t_{1, i}, t_{2, i}\right) i=1, \ldots I$ such that $\left|K^{-1}\left(t_{2, i}, P\right)-K^{-1}\left(t_{1, i}, P\right)\right| \neq 0$, define
$\widehat{\tau}_{b_{n}}=I^{-1} \sum_{j=1}^{I}\left|\bar{K}_{U, b_{n}}^{-1}\left(t_{2, i} \mid X^{n}\right)-\bar{K}_{U, b_{n}}^{-1}\left(t_{1, i} \mid X^{n}\right)\right| /\left|K^{-1}\left(t_{2, i}, P\right)-K^{-1}\left(t_{1, i}, P\right)\right|$.
Then $\widehat{\tau}_{b_{n}}=\tau_{b_{n}}\left(1+o_{P}(1)\right)$. Define

$$
\widehat{\mu}_{b_{n}, I}=I^{-1} \sum_{i=1}^{I}\left(K_{U, b_{n}}^{-1}\left(t_{i} \mid X^{n}\right)\right)-\widehat{\tau}_{b_{n}}^{-1} K^{-1}\left(t_{i} \mid X^{n}\right)
$$

Then

$$
\widehat{\mu}_{b_{n}, I}=\mu_{b_{n}}+o_{P}\left(\tau_{b_{n}}^{-1}\right)
$$

Moreover,

$$
q^{-1} \sum_{i=1}^{q} 1\left\{\widehat{\tau}_{b_{n}}\left(U_{b_{n, i}}-\widehat{\mu}_{b_{n}}\right) \leq x\right\} \rightarrow K(x)
$$

Theorems 4 and 5 are useful in standardizing the subsampling distribution to obtain -for instance- predictive confidence intervals. Theorem 4 can be used to obtain the shape of the normalizing sequence and of the asymptotic distribution
when everything is practically unknown. We emphasize however, that the rate $\tau_{n}$ (and the centering product $\tau_{n} \mu_{n}$ ) are only estimable up to a multiplicative constant because of the identifiability issues we discussed previously.

We now follow up on the "difference of quantiles" idea that originated in Bertail et al. (1999). When the functional form of the rate $\tau_{n}$ is known up to some finite dimensional parameter, it is quite easy to see how the ideas of Section 2 may be implemented to obtain a convergent estimator of the norming sequences. Using a difference of quantiles for some fixed $t_{1}>t_{2}$ we define

$$
\begin{align*}
y\left(t_{1}, t_{2}, b_{n}\right) & \equiv \log \left(\left|\bar{K}_{U, b_{n}}^{-1}\left(t_{2} \mid X^{n}\right)-\bar{K}_{U, b_{n}}^{-1}\left(t_{1} \mid X^{n}\right)\right|\right)  \tag{16}\\
& =y\left(t_{1}, t_{2}, \infty\right)-\log \tau_{b_{n}}\left(1+o_{P}(1)\right) \\
\text { with } y\left(t_{1}, t_{2}, \infty\right) & =\log \left(\left|K^{-1}\left(t_{2}, P\right)-K^{-1}\left(t_{1}, P\right)\right|\right)
\end{align*}
$$

If we ensure that $\tau_{n}$ is of the form $n^{-\gamma}$ or $h(n)^{-\gamma}$ for some $\gamma \neq 0$ and some (known) real function $h$ then a straightforward regression gives the value of the index $\gamma$ as before. This method is referred as the "range" method in Bertail et al.(1999) and is illustrated by the following result.

Theorem 6 Assume that $\tau_{n}=n^{-\gamma}$ for some $\gamma \neq 0$, , let $\left(t_{1, j}, t_{2, j}\right) j=1, \ldots J$ be some points in $(0,1)^{2}$, and $b_{i, n}=n^{\beta_{i}}, 1<i<I$, be some subsampling sizes satisfying the assumptions of Theorem 1(b) ; then the least-squares estimator of $\gamma\left(s a y \widehat{\gamma}_{I, J, n}\right)$ in the model

$$
y\left(t_{1, j,}, t_{2, j}, b_{i, n}\right)=y\left(t_{1, j,} t_{2, j}, \infty\right)+\gamma \log b_{n, i}+u_{j, i}
$$

is such that

$$
\widehat{\gamma}_{I, J, n}=\gamma+o_{P}\left((\log n)^{-1}\right)
$$

The proof of this result is similar to the proof of Theorem 2 and is skipped.

### 3.3 General rates and slowly varying functions

In many interesting situations, $\tau_{n}$ is a regular varying function of index $\gamma$, say $\tau_{n}=n^{\gamma} L(n)$ where $L$ is a slowly varying function that is such that, for any $\lambda>0, \lim _{x \rightarrow \infty} \frac{L(\lambda x)}{L(x)}=1$ (see Bingham, Goldie and Teugels (1987)). Since the behavior of $L$ is only interesting for $x$ tending to $\infty$, the problem of estimating $L$ amounts in fact to finding an "estimated" equivalent of $L$ for large $x$, up to a multiplicative constant term. The Karamata representation theorem states that there exists a measurable bounded function $\varepsilon($.$) , with \varepsilon(u) \longrightarrow 0$ as $u \rightarrow \infty$ such that $L(n)=\exp \left(\int_{1}^{n} u^{-1} \varepsilon(u) d u\right)=\exp \left(h_{1}(\log n)\right)$, where $h_{1}$ is a $C^{\infty}$ function with all derivatives vanishing at $\infty$ (see Bingham et al. (1987), Th. 1.3.1 and

Th 1.3.3) related to $\varepsilon$ by $\varepsilon(t)=h_{1}^{(1)}(\log t)$, where $h_{1}^{(1)}$ is the first derivative of $h_{1}$. Therefore (16) may be rewritten

$$
\begin{align*}
y\left(t_{1}, t_{2}, b_{n}\right) & =y\left(t_{1}, t_{2}, \infty\right)-\gamma \log \left(b_{n}\right)+\log L\left(b_{n}\right)+o_{P}(1)  \tag{17}\\
& =y\left(t_{1}, t_{2}, \infty\right)-\gamma \log \left(b_{n}\right)+h_{1}\left(\log b_{n}\right)+o_{P}(1)
\end{align*}
$$

With $L\left(b_{n}\right) \neq 1$, the preceeding method of regressing $y\left(t_{1}, t_{2}, b_{n}\right)$ on $\log b_{n}$ at some different points $\left(t_{1, j}, t_{2, j}\right)$ and different $b_{n, i}$ still yields a convergent estimator of $\gamma$ but at a rate slower than $\log (n)^{-1}$ (the convergence is due to the fact that $h_{1}\left(\log b_{n}\right)=o\left(\log b_{n}\right)$. Try for instance $h_{1}(x)=\log x$ (in that case $L(n)=\log n)$ and $b_{n, i}=n^{\beta_{i}}$ as in Theorem 5; then $\widehat{\gamma}_{I, J, n}=\gamma+O\left(\frac{\log \log n}{\log n}\right)$. Unfortunately this convergence rate is not sufficient to ensure that $n^{\widehat{\gamma}_{I, J, n}}=$ $n^{\gamma}(1+o(1))$ and thus can not be used as a normalizing factor. In that sense, the results in the preceding sections are not robust to the presence of a slowly varying function. It is clear from this situation that $\gamma$ and $h_{1}$ should be estimated simultaneously. As noticed in Bertail et al. (1999), model (17) is known as a partial spline regression model in the econometric literature (see Engle et al. (1986)) and may be estimated using semi-parametric methods. As a particular case, the same method could be applied in the three cases considered in section 3.1 regardless of the domain of attraction of the distribution. However it seems to us that this is quite overoptimistic because this method would require the computation of a very large number of subsampling distributions and in any case would require huge sample sizes.

For this reason, we now present a more attractive construction under a mild restriction on the slowly varying function. A quite simple and reasonable hypothesis on $L$ will indeed allow us to give more straightforward estimators of $\gamma$ and $L$. The restriction amounts to assuming that

$$
\begin{equation*}
\lim _{x \rightarrow \infty} L\left(x^{2}\right) / L(x)=\exp \left(-C_{0}\right) \tag{18}
\end{equation*}
$$

where $C_{0}$ is some constant. Recall that by the definition of a slowly varying function we have that $L(\lambda x) / L(x)$ converges to 1 for any $\lambda>0$; that is, in (18) we allow $\lambda$ to be of the same order as $x$. This hypothesis clearly holds for the "usual" slowly varying functions of the form $L(x)=\Pi_{i}\left(\log ^{c_{i}} x\right)^{a_{i}}$ where $a_{i}$ are some constants and $\log ^{c_{i}}$ denotes the $c_{i}$-iterated logarithm.

Then for $\lambda_{n}=b_{n}$ such that $b_{n}^{2} / n \rightarrow 0$ we have

$$
\begin{equation*}
y\left(t_{1}, t_{2}, b_{n}\right)-y\left(t_{1}, t_{2}, \lambda_{n} b_{n}\right)=\gamma \log \left(b_{n}\right)+C_{0}+o_{P}(1) \tag{19}
\end{equation*}
$$

It follows from the same arguments as in Section 2 that, using different subsampling sizes, $b_{n, j}, j=1, \ldots, J \geq 2$, and points $\left(t_{1, i}, t_{2, i}\right) i=1, \ldots I$, the least squares estimator of $\gamma$ in (19) gives a convergent estimator $\widehat{\hat{\gamma}}_{I, J}$ of $\gamma$ up to $o_{P}\left(\log (n)^{-1}\right)$. Now we can plug this estimator in (17) to get (for two fixed points $t_{1,0}$ and $\left.t_{2,0}\right)$

Thus the form $h_{1}\left(\log b_{n}\right)$ can be estimated up to the constant term $y\left(t_{1,0}, t_{2,0}, \infty\right)$. Since this constant is of no importance in the construction of the subsampling distribution, we can choose

$$
\widehat{h}_{1}\left(\log \left(b_{n}\right)\right)=y\left(t_{1,0}, t_{2,0}, b_{n}\right)+\widehat{\hat{\gamma}}_{I, J} \log \left(b_{n}\right)
$$

as a convergent estimator of $h_{1}\left(\log b_{n}\right)$. This amounts to standardizing the asymptotic distribution by imposing $y\left(t_{1,0}, t_{2,0}, \infty\right)=1$.

In the following we put

$$
\widehat{L}\left(b_{n}\right)=\exp \left(\widehat{h}_{1}\left(\log \left(b_{n}\right)\right)\right.
$$

and

$$
s(t)=\frac{p(t)}{\int_{0}^{1} p(u) d u}
$$

with

$$
p(t)=\exp \left(-\frac{1}{x}-\frac{1}{1-x}\right)
$$

We can use the extrapolation scheme of Adamovic (see for instance Bingham et al. (1987), Th. 1.3.4) to construct an asymptotic equivalent of $L(x)$ when $x$ $\rightarrow \infty$. Let $b_{n, 0}^{*}=0<b_{n, 1}^{*}<b_{n, 2}^{*}<\ldots \ldots<b_{n, K_{n}}^{*}$, be a sequence of subsampling sizes. Define for any $x>A>0, x \in\left[b_{n, i}^{*}, b_{n, i+1}^{*}\right]$, the extrapolation

$$
\begin{equation*}
\widehat{L}(x)=\widehat{L}\left(b_{n, i}^{*}\right)+\left(\widehat{L}\left(b_{n, i+1}^{*}\right)-\widehat{L}\left(b_{n, i}^{*}\right)\right) \int_{0}^{\left(x-b_{n, i}^{*}\right) /\left(b_{n, i+1}^{*}-b_{n, i}^{*}\right)} s(t) d t \tag{20}
\end{equation*}
$$

Now it is apparent that $\widehat{L}$ is a $C^{\infty}$ slowly varying function that coincides with $\widehat{L}\left(b_{n, i}^{*}\right)=\exp \left(\widehat{h}_{1}\left(\log \left(b_{n, i}^{*}\right)\right)\right.$ at each $b_{n, i}^{*}$. The following theorem states that under some regularity assumptions on the subsampling sizes, the suggested construction gives consistent estimators of $\gamma$ and $L$; note also that the theorem addresses the general case where $\gamma$ may be any real number (positive or negative).

Theorem 7 Assume that we have normalized the asymptotic distribution such that $\left|K^{-1}\left(t_{2,0}, P\right)-K^{-1}\left(t_{1,0}, P\right)\right|=1$ at some fixed points $t_{1,0}$ and $t_{2,0}$. Assume that $\tau_{n}=n^{\gamma} L(n)$ and that $L(x)$ is a normalized slowly varying function satisfying (18); let $\left(t_{1, i}, t_{2, i}\right) i=1, \ldots I$ be some real numbers such that $t_{1, i} \neq t_{2, i}$ for all $i=1, \ldots I$. Choose $b_{n, j}, j=1, \ldots J \geq 2$, such that $b_{n, j}^{2}$ satisfy the hypotheses of Theorem 1(b). Then the least-squares estimators of $\gamma$ and $C_{0}$ in (19) satisfy

$$
\widehat{\hat{\gamma}}_{I, J}=\gamma+o_{P}\left((\log n)^{-1}\right)
$$

and

$$
\widehat{C}_{o}=C_{0}+o(1)
$$

Moreover, for $b_{n, 0}^{*}=0<b_{n, 1}^{*}=n^{1 / 2}<b_{n, 2}^{*}<\ldots \ldots<b_{n, K_{n}}^{*}$ such that $b_{n, ; j+1}^{*} / b_{n, j}^{*}<e, \quad j=1, \ldots M_{n}$, with $M_{n}=o(\log n)$ but $M_{n} \rightarrow \infty$ then we have that the estimator in (20) satisfies

$$
\frac{L(x)}{\widehat{L}(x)}=O_{P}(1), \frac{\widehat{L}(x)}{L(x)}=O_{P}(1), \text { for } x \rightarrow \infty
$$

Finally, we have

This theorem states that it is possible both to obtain simple estimates of the index and the slowly varying function (provided that the behavior of $L(x)$ is similar to that of $L\left(x^{2}\right)$ ). Of course, we need a lot of subsampling distributions (in fact an infinity) to estimate an asymptotic equivalent of $L$ for large $x$. However -and this makes our method more attractive than the non-parametric method suggested before- we need only one subsampling distribution at $b_{n, 1}^{*}=n^{1 / 2}$ to obtain an asymptotic equivalent of $L(n)$. This is clearly due to the fact that in the smooth case considered here, $L(n)$ is asymptotically equivalent to $L\left(n^{1 / 2}\right)$. Because we have imposed $\left|K^{-1}\left(t_{2,0}, P\right)-K^{-1}\left(t_{1,0}, P\right)\right|=1$ at some points $t_{2,0}$ and $t_{1,0}$, the convergence rate is identifiable: this explains why we need to take into account $\widehat{C}_{0}$ in the definition of $\widehat{\hat{\tau}}_{n}$ although constants do not play a real role in our analysis. Of course, the choice of $\mathrm{t}_{2,0}$ and $\mathrm{t}_{1,0}$ is completely arbitrary: we suggest using $t_{2,0}=1-t_{1,0}=0.6915$ so as to have some elements of comparison with the reduced normal distribution.

Finally, note that estimation of the centering factor $\mu_{n}$ may be treated in the same fashion under similar assumptions on its shape.

## 4 The extreme value statistic revisited

### 4.1 The Type I and Type III Cases

To illustrate the results of sections 3.2 and 3.3 we again consider the extreme value statistic of section $3.1 T_{n}=\max \left(X_{1}, \ldots, X_{n}\right)$. If we know that $F$ belongs to $D(\Lambda)$ then (16) may be used to obtain a parametric or nonparametric estimator of $\tau_{n}$. Since in that case the asymptotic distribution is entirely known and equal to $K=\Lambda$ there is no reason to standardize, and an estimator of the rate of convergence is given under the hypotheses of Theorem 7 by $\widehat{\hat{\tau}}_{n}\left|\Lambda^{-1}\left(t_{2,0}\right)-\Lambda^{-1}\left(t_{1,0}\right)\right|$. Other estimators may be proposed in that case since one may also take advantage of the special form of $\tau_{n}$ and $\mu_{n}$ as a function of $H^{-1}$. In the domain of attraction of type (iii) it is common that $\tau_{n}$ and $\mu_{n}$
are respectively of the form $\log n^{-a}$ and $\log n^{b}$ with $b>a$. Standard linear regressions may then be used to obtain convergent estimators of $a$ and $b$.

If one does not know anything about the tail of the distribution, then nonparametric estimators or our Theorem 7 will first give convergent estimators of the rate. Notice that we can get a lot of information directly from the estimated rate. If the estimated rate is increasing with $n$ then this means that $F \in D\left(\Psi_{\gamma}\right)$. In that case we may indifferently use $T_{n}$ as a recentering factor; alternatively, we may simply use the median of the subsampling distribution. If the estimated rate appears to be decreasing, then the median at several $b_{n}$ gives the shape of the recentering factor. In any case, it is possible to construct an approximation of the sampling distribution of the maximum without any information on the tail of the distribution. Since type II is of particular importance in financial applications we will review that case in the next section.

### 4.2 The Type II Case

Assume a regularity condition on the tail of $F$, namely that

$$
\begin{equation*}
x^{\alpha}(1-F(x)) \rightarrow \Gamma \tag{21}
\end{equation*}
$$

as $x \rightarrow \infty$, for some $0<\alpha<2$ that measures the 'thickness' of the tail of $F$; here $\Gamma>0$ is a scaling constant. Under the strong mixing condition and (21) it is well-known that, if $K_{n}(x, P)$ has a large-sample limit, then it must be of a specific form (the type II extreme distribution of our Section 3.1); see Leadbetter et al. (1983) where sufficient conditions for the existence of a limit for $K_{n}(x, P)$ are also found by means of the notion of 'distributional mixing'. Furthermore, the divergence rate $\tau_{n}$ is also specified in this case: $\tau_{n}$ is proportional to $n^{-1 / \alpha}$.

To fix ideas, we now set $\tau_{n}=n^{-1 / \alpha}$; it now follows that under some mixing conditions and (21) we have ${ }^{3}$

$$
K_{n}(x, P) \equiv P\left\{n^{-1 / \alpha} T_{n} \leq x\right\} \rightarrow K(x, P)= \begin{cases}\exp \left(-x^{-\alpha} / \Gamma\right) & \text { if } x>0  \tag{22}\\ 0 & \text { if } x \leq 0\end{cases}
$$

as $n \rightarrow \infty$.
Therefore, in this specific case, not only do we have existence of a limit distribution $K(x, P)$, but we know the form of $K(x, P)$ as well (up to the two unknown parameters $\Gamma$ and $\alpha$ ). Since $\alpha=1 / \gamma$ (where $\gamma$ is as in Theorem 6), an

[^3]estimator of $\alpha$ could be obtained via our Lemma 2 or Theorem 2; these however do not take our explicit knowledge of $K(x, P)$ into account.

Nevertheless, we can modify the subsampling methodology of Section 2 to take knowledge of $K(x, P)$ into account. Note that since $K(x, P)$ is continuous, (22) and Polya's theorem imply that

$$
\sup _{x>0}\left|P\left\{n^{-1 / \alpha} T_{n} \leq x\right\}-K(x, P)\right| \rightarrow 0
$$

which in turn implies that

$$
\sup _{y>0}\left|P\left\{T_{n} \leq y\right\}-\exp \left(-\Gamma^{-1} n y^{-\alpha}\right)\right| \rightarrow 0 .
$$

As in Section 2, we may approximate $P\left\{T_{n} \leq y\right\}$ by $\bar{K}_{b_{n}}\left(y \mid X^{n}\right)$, which suggests that

$$
\log \left(-\log \left|\bar{K}_{b_{n}}\left(y \mid X^{n}\right)\right|\right)=-\log \Gamma+\log b_{n}-\alpha \log y+o_{P}(1)
$$

treating the above equation as a regression of $\log \left(-\log \left|\bar{K}_{b_{n}}\left(y \mid X^{n}\right)\right|\right)$ on $\log y$ we can use different values of $y$ to get least squares estimates of both $\Gamma$ and $\alpha$.

Alternatively, we can look at quantiles again as in Section 2. Note that due to (22) we have an explicit expression for $K^{-1}(t, P)$ as well, namely

$$
K^{-1}(t, P)=(-\Gamma \log t)^{-1 / \alpha}
$$

The above equation affords us the opportunity to use (6) to simultaneously estimate $\Gamma$ and $\alpha$ by least squares; alternatively, we could use an estimate $\hat{\alpha}$ of $\alpha$ (say one obtained by Lemma 2 or Theorem 2 to obtain an estimate of $\Gamma$. To elaborate on the latter, note that by Lemma 3 we have

$$
b_{n}^{-1 / \hat{\alpha}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=(-\Gamma \log t)^{-1 / \hat{\alpha}}+o_{P}(1)
$$

which can be solved for $\Gamma$. To improve upon this estimate, we can look at different $t$-points and different $b_{n} \mathrm{~s}$ and take an average of the corresponding $\Gamma$-solutions.

Another possibility for the case of interest here, i.e. where $T_{n}=\max \left\{X_{1}, \ldots, X_{n}\right\}$, is to use the simple and popular Hill's estimator of $\alpha$ instead of our proposals; see de Haan (1994), de Haan et al. (1994), and Jansen et al. (1998) for a description of Hill's estimator. We note a growing disenchantment in recent literature regarding Hill's estimator; for example, Kearns and Pagan (1997) note that Hill's estimator can be misleading if there is some dependence in the data, whereas Resnick (1997) points out that Hill's estimator can exhibit some occasional strange behavior, and that it certainly can not account for the possible presence of a slowly varying function in the rate.

Recall that to implement Hill's estimator the number $k$ of extreme order statistics that are to be taken into account must be chosen by the practitioner. In this case, $k$ is like a 'smoothing parameter', and its proper choice is crucial
for good results; see e.g. Hall (1990) for a discussion. To implement our estimators of $\alpha$ using Lemma 2 or Theorem 2, a 'smoothing parameter' -namely the subsample sizes used- must also be chosen. Although the procedures were shown to be valid for $b_{i, n}$ proportional to $n^{\beta_{i}}$, the exact optimal choice of the $\beta_{i} \mathrm{~S}$ (and of the proportionality constants) is an open problem at the moment.

## 5 Econometric applications

### 5.1 Data

For our empirical applications we study time series of daily asset returns. For the Value at Risk example, we examine daily returns of the S\&P 500 Index, $X_{t}=\ln \left(P_{t} / P_{t-1}\right)$ where $P_{t}$ is the daily closing $\mathrm{S} \& \mathrm{P} 500$ index. Our observations start January 3, 1985 and end December 31, 1996, a total of $n=3033$ daily observations. For our safety first portfolio allocation example, we consider portfolios formed from the S\&P 500 (an index of "large cap" stocks), the Russel 2000 (an index of "small cap" stocks) and the 30-day Treasury Bill ("TBill"). Our observations for the Russel 2000 and the TBill cover the same period as for the S\&P 500. Summary statistics for the daily returns of our three assets are reported in table 1.

| Daily Returns in Percent |  |  |  |
| :--- | ---: | ---: | ---: |
|  | Russel2000 | S\&P 500 | TBill |
| Mean | 0.05535 | 0.05926 | 0.02257 |
| Std. | 0.7718 | 0.9645 | 0.0062 |
| Min | -12.52000 | -20.33500 | 0.01068 |
| Max | 7.72300 | 9.09940 | 0.03640 |
| n | 3033 | 3033 | 3033 |

Table 1: Descriptive statistics for the data.

### 5.2 Estimating Value-at-Risk

In computing our various VaR estimates we simulate the manner in which these estimates would be deployed in a "live" application. Specifically, we start with a sample of $n_{1}=2275$ observations (1/3/1994) and compute the 1-Day, 10-Day, or 23 -Day VaR estimate. The next day we have a sample of $n_{2}=2276$ observations and we repeat the exercise, continuing in this way until we accumulate an "evaluation" sample of 3 years or 758 observations on estimated VaR. We evaluate the performance of the various estimation methods by computing the observed exceedance rate for the evaluation period. An unbiased method should
have exceedance equal to the VaR level. For example, the $5 \%$ VaR estimates should have a $5 \%$ exceedance for unbiasedness.

We assume $\mu=0$ and use the subsampling methodology to estimate $1 \%$ and $5 \% \mathrm{VaR}$ for 1,10 , and 23 trading days. The subsampling estimates of the $n$-day VaR were arrived at by using Lemma 1 and the assumption that the divergence rate $\tau$ is proportional to the block length:

$$
\begin{aligned}
\tau_{n} \bar{K}_{n}^{-1}\left(p, X^{i}\right) & =K^{-1}(p, P)+o_{p}(1) \\
\tau_{\kappa} \bar{K}_{\kappa}^{-1}\left(p, X^{i}\right) & =K^{-1}(p, P)+o_{p}(1)
\end{aligned}
$$

Dividing the two equations yields:

$$
\begin{align*}
\bar{K}_{n}^{-1}\left(p, X^{i}\right) & =\frac{\tau_{\kappa}}{\tau_{n}} \bar{K}_{\kappa}^{-1}\left(p, X^{i}\right)+o_{p}(1) \\
\bar{K}_{n}^{-1}\left(p, X^{i}\right) & =\left(\frac{\kappa}{n}\right)^{-\gamma} \bar{K}_{\kappa}^{-1}\left(p, X^{i}\right)+o_{p}(1) \tag{23}
\end{align*}
$$

The choice of the largest block length is based on the following procedure: At every point $t$ in time, split the available data into two subsets. For an initial value of the largest block length estimate the parameter $\gamma$ based on data up to period $t-200$ and compute predicted and actual exceedances for the interval $(t-200, t)$. Then adapt the largest block length if the actual and desired exceedance do not match, reestimate and reevaluate until a choice for the largest block length is found that minimizes the deviation from the actual exceedance over the hold-out period of the past 200 days. Use this optimal block length to determine the VaR for the next consecutive five days. Then again reoptimize over the past 200 observations. The optimal block lengths for the first day of our out-of-sample evaluation, January 3, 1994, are reported in table 2. Figures 1 and 2 report the optimal block lengths over the complete out-of-sample period and figures 3 and 4 show the corresponding hold-out sample exceedances. For strictly stationary data it should not be necessary to reoptimize the block size, however, we reoptimize the block size every five trading days in order to take possible nonstationarities in the data into account. This reoptimization leads to time-varying estimates of $\gamma$, the evolution of which is shown in figures 5 and 6.

|  | 1 Day | 10 Days | 23 Days |
| :--- | :---: | :---: | :---: |
| $1 \%$ VaR | 114 | 128 | 239 |
| $5 \%$ VaR | 114 | 136 | 227 |

Table 2: Optimal Block Lengths for the various VaR specifications based on hold-out samples.

|  | Benchmark |  | Subsampling |  |
| :--- | :---: | :---: | :---: | :---: |
|  | $1 \% \mathrm{VaR}$ | $5 \% \mathrm{VaR}$ | $1 \% \mathrm{VaR}$ | $5 \% \mathrm{VaR}$ |
| 1 Day | 0.3 | 3.1 | 1.8 | 6.5 |
| (-QLL) | $(.025)$ | $(.082)$ | $(.024)$ | $(.081)$ |
| 10 Days | 0 | 0.1 | 1.8 | 7.6 |
|  | $(.054)$ | $(.151)$ | $(.051)$ | $(.182)$ |
| 23 Days | 0 | 0 | 3.9 | 8.6 |
|  | $(.125)$ | $(.257)$ | $(.102)$ | $(.349)$ |

Table 3: Results for VaR estimation: Actual exceedance (in percent) of the estimated VaR over a period of 758 trading days, $1 / 3 / 94$ to $12 / 31 / 1996$. Negative Quasi Log Likelihood ${ }^{4}$ in parentheses (smaller values indicate a better goodness of fit).

Results are reported in table 3. The benchmark is computed by simply taking the $1^{s t}$ or $5^{t h}$ percentile of the distribution of historical returns which is also a valid subsampling estimator. Figures $7-9$ show the 1 day, 10 day and 23 day out-of-sample Value-at-Risk estimates, respectively.

### 5.3 Safety first portfolio selection

The safety first method of portfolio selection introduced by Roy (1952) provides an interesting alternative to the traditional mean-variance approach of Markowitz (1952). Recently, Jansen et al. (1998) showed how extreme value theory can be used in a multi-period setting to make operational the version of safety first portfolio selection formulated by Arzac and Bawa (1977). Here we show how our subsampling approach can be used in implementing multi-period safety first portfolio selection.

We consider a portfolio of $k$ risky assets, whose one period net random returns are given by $r_{j}, j=1, \ldots, k$. When a proportion $w_{j}$ of the initial portfolio value is held in asset $j$, the one period random return on the portfolio is

$$
r(w)=\sum_{j=1}^{k} w_{j} r_{j}
$$

where $w$ is the $k \times 1$ vector of weights, which by convention sum to unity $\left(i^{\prime} w=1\right.$, where $i$ is the $k \times 1$ vector of 1 s$)$. Let the $d^{t h}$ quantile of $r(w)$ be given by $q_{d}(w)$, so that

$$
P\left[r(w) \leq q_{d}(w)\right]=d
$$

[^4]If the investor can borrow or lend risklessly at the fixed interest rate $r$ over the period, the safety first separation theorem of Arzac and Bawa (1977) specifies that the investor can obtain optimal portfolio weights $w^{*}$ by solving the problem

$$
\max _{w}[E(r(w))-r] /\left[r-q_{d}(w)\right] .
$$

Once optimal weights $w^{*}$ are given, the amount to be invested is determined as $W+b^{*}$, where $W$ is the investor's initial wealth and $b^{*}$ is optimal borrowing (lending if $b^{*}<0$ ), determined as

$$
\frac{b^{*}}{W}=\left[s-q_{d}\left(w^{*}\right)\right] /\left[q_{d}\left(w^{*}\right)-r\right] .
$$

The parameters $s$ and $d$ summarize the investor's preferences in that, as Arzac and Bawa (1977) prove, the solution just described ensures that the investor maximizes expected return subject to the constraint that wealth at the end of the period $\left(=\left(1+r\left(w^{*}\right)\right) W+b^{*}\left(r\left(w^{*}\right)-r\right)\right.$ falls below a critical level $(1+s) W$ (typically $s<0$ ) with probability $d$.

Jansen et al. (1998) note that investors typically invest for prolonged periods of time, not just for a single period. Nevertheless, a single period of unusually poor performance can be disastrous, especially for fund managers whose clientele tend to be quite mobile. Consequently, Jansen et al. (1998) replace the single period return quantile $q_{d}(w)$ in the safety first optimization above with the quantile for the worst ever single period return $\tilde{q}_{d}(w)$. They show how extreme value theory permits a useful approximation to these quantiles, making use of Hill's (1975) estimator of the tail index $\alpha$.

Alternatively, we can use our subsampling estimator (based on Theorem 2) to approximate the desired quantiles. We apply this decision rule to determine portfolios composed of the S\&P 500, the Russel 2000, and leverage using TBills. Before we discuss the results we briefly discuss estimation - related issues.

### 5.3.1 Hill's estimator for the tail index

Let $\left\{\xi_{i}\right\}_{i=1}^{n}$ be a sequence of random numbers obeying some probability law $F$ which satisfies the necessary regularity conditions. Let $X_{(1)} \leq X_{(2)} \leq \ldots \leq X_{(n)}$ denote the associated order statistics. Furthermore define $\gamma$ to be $1 / \alpha$.

Probably the currently most prominent estimator for the tail index was proposed by Hill (1975) which was derived in a maximum likelihood framework. The estimator is given by

$$
\begin{equation*}
\hat{\gamma}=k^{-1} \sum_{i=1}^{k}\left(\ln X_{(n+1-i)}-\ln X_{(n-k)}\right) \tag{24}
\end{equation*}
$$

Asymptotic normality of the Hill estimator has been shown by e.g. Goldie and Smith (1987) who proved that

$$
\begin{equation*}
k^{1 / 2}(\hat{\gamma}-\gamma) \xrightarrow{d} N\left(0, \gamma^{2}\right) \tag{25}
\end{equation*}
$$

if $k$ increases suitably rapidly as $n \rightarrow \infty$.

### 5.3.2 Nonparametric Subsampling

When we perform the nonparametric subsampling approach we assume a divergence rate $\tau_{n}=n^{-\gamma}$ and estimate $\gamma$ using a "nonparametric" estimator obtained from the regressions of sections 2 and 3 . The recentering factor $\mu_{n}$ is assumed to be zero which is the case for distributions in the domain of Type II extreme value distributions. For the problem at hand we can exploit Geluk and De Haan's (1987) proposition 1.7 which states that the density with the fatter tails always dominates in a mixture of two densities. Hence we estimated $\gamma$ for portfolios with different compositions and kept increasing the block length until the estimates for $\gamma$ were roughly the same. Based on this method we arrive at estimates for $\gamma$ of about .42 for the Russel 2000 and .45 for the S\&P 500 based on block lengths from 3 to 400 .

We simulate the decision process for a critical loss $s$ of $-0.1 \%,-0.5 \%$, and $-1 \%$. Instead of the targeted $1 \%$ of times that we exceeded the critical loss of $s \%$, table 4 reports that $s$ is never undercut by the subsampled portfolio and $13 \%$ of the time by the Hill portfolio. We note, furthermore, that all Hill - portfolios are far too highly leveraged which comes from the fact that the estimated quantiles are too high (i.e. the estimated distribution has too little mass in the tails). Figures $10-12$ show how the value of the portfolios evolves over the out-of-sample period. The Hill safety first portfolio performs more like a "safety last" portfolio. In contrast, the subsampling safety first portfolio exhibits the expected conservative behavior. As expected for such risk averse values of $s$ the subsampling safety first portfolio provides lower returns than either of the stock indexes. However, as $s$ goes from $0.1 \%$ to $1 \%$ the leverage in the portfolio increases and it starts clearly distinguishing itself from a pure TBill position.

## 6 Conclusions

In this paper we extend the results of Bertail et al. (1999) and show how subsampling can be used to approximate the sampling distributions of diverging statistics We provide consistent estimators both in the case of known and unknown divergence rates and illustrate how to estimate measures of tail fatness for dependent data.

Our method performs well in estimating value at risk and our application to safety first portfolio selection indicates that for time series our subsampling estimator constitutes a clear improvement over the conventional Hill estimator which has been devised for iid data.

## 7 *Acknowldegements

We thank Bob Trippi for the provision of computer resources to perform the numerically intensive tasks of section 5 .

|  |  | $\mathrm{S}=.1 \%$ |  | $\mathrm{~s}=.5 \%$ |  | $\mathrm{~S}=1 \%$ |  |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | T-Bill | SubS | Hill | SubS | Hill | SubS | Hill |
| Wealth after 3 years: | 1.160 | 1.166 | 1.239 | 1.182 | 1.518 | 1.202 | 1.918 |
| average annualized returns: | 5.067 | 5.250 | 7.395 | 5.728 | 14.918 | 6.325 | 24.211 |
| \% of times below 0.1\% | 0.000 | 0.000 | 13.852 | - | - | - | - |
| \% of times below 0.5\% | 0.000 | - | - | 0.000 | 13.852 | - | - |
| \% of times below 1.0\% | 0.000 | - | - | - | - | 0.000 | 13.852 |
| Returns (in \%): |  |  |  |  |  |  |  |
| Mean: | 0.020 | 0.020 | 0.028 | 0.022 | 0.057 | 0.024 | 0.092 |
| Standard Deviation: | 0.003 | 0.009 | 0.131 | 0.038 | 0.569 | 0.074 | 1.118 |
| Minimum: | 0.012 | -0.023 | -0.571 | -0.167 | -2.530 | -0.348 | -4.979 |
| Maximum: | 0.023 | 0.052 | 0.593 | 0.178 | 2.623 | 0.337 | 5.161 |
| Leverage (b*): |  |  |  |  |  |  |  |
| Mean: | - | -0.985 | -0.778 | -0.935 | -0.036 | -0.872 | 0.891 |
| Standard Deviation: | - | 0.001 | 0.011 | 0.002 | 0.052 | 0.004 | 0.104 |
| Minimum: | - | -0.987 | -0.802 | -0.946 | -0.145 | -0.894 | 0.676 |
| Maximum: | - | -0.983 | -0.749 | -0.928 | 0.113 | -0.859 | 1.199 |
| Portfolio weight for RU2000: |  |  |  |  |  |  |  |
| Mean: | - | 0.622 | 0.467 | 0.622 | 0.467 | 0.622 | 0.467 |
| Standard Deviation: | - | 0.026 | 0.319 | 0.026 | 0.319 | 0.026 | 0.319 |
| Minimum: | - | 0.550 | 0.000 | 0.550 | 0.000 | 0.550 | 0.000 |
| Maximum: | - | 0.750 | 1.000 | 0.750 | 1.000 | 0.750 | 1.000 |

Table 4: Results for the Safety First Portfolio Selection ( $\mathrm{d}=1 \%$ ).

## Appendix: Technical proofs

Proof of Theorem 1 Fix an $x$ and first note that $E K_{b_{n}}\left(x \mid X^{n}, \tau\right)=$ $K_{b_{n}}(x, P)$ because of stationarity of the $X$-sequence. Now by an argument similar to the proof of Theorem 3.1 in Politis and Romano (1994) we have that

$$
\operatorname{Var}\left(K_{b_{n}}\left(x \mid X^{n}, \tau_{.}\right)\right)=O\left(b_{n} / n\right)+O\left(n^{-1} \sum_{i=1}^{n} \alpha_{X}(i)\right)
$$

note that both terms on the right-hand side above tend to zero under the assumed conditions - the second term being the Cesaro sum of a null sequence, and the first term due to $b_{n} / n \rightarrow 0$.

Therefore, $K_{b_{n}}\left(x \mid X^{n}, \tau.\right)=K_{b_{n}}(x, P)+o_{P}(1)$. Letting $b_{n}=b$ (a constant), part (a) is proved. Now letting $b_{n} \rightarrow \infty$, note that $K_{b_{n}}(x, P) \rightarrow K(x, P)$ by assumption (1). Since $K(x, P)$ is assumed continuous, Polya's theorem yields the uniform consistency in (3).

Finally, if $d_{n} \rightarrow \infty$, then note that $\sup _{x}\left|K_{d_{n}}(x, P)-K(x, P)\right| \rightarrow 0$ as well by (1). Thus (4) is a direct consequence of (3).

Proof of Lemma 1 First note that using arguments similar to those in the proof of Lemma 1 in Bertail et al. (1999) it can be shown here too that

$$
\begin{equation*}
K_{b_{n}}^{-1}\left(t \mid X^{n}, \tau_{.}\right)=K^{-1}(t, P)+o_{P}(1) \tag{26}
\end{equation*}
$$

for any $t \in(0,1)$. Also note that by definition of the quantile transform we have

$$
\begin{equation*}
\bar{K}_{b_{n}}\left(x \tau_{b_{n}}^{-1} \mid X^{n}\right)=K_{b_{n}}\left(x \mid X^{n}, \tau_{.}\right) \tag{27}
\end{equation*}
$$

and thus, for any $t \in(0,1)$,

$$
\begin{equation*}
K_{b_{n}}^{-1}\left(t \mid X^{n}, \tau_{.}\right)=\tau_{b_{n}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right) \tag{28}
\end{equation*}
$$

as well. As a consequence of (26) and (28) we immediately have that (6) is true.

Proof of Lemma 2 The proof follows from the more general Theorem 2 below by letting $J=1$.

Proof of Lemma 3 Note that since $\hat{\gamma}=\gamma+o_{P}\left((\log n)^{-1}\right)$ it follows that $b_{n}^{-\hat{\gamma}} / b_{n}^{-\gamma} \rightarrow 1$ in probability, and $b_{n}^{-\hat{\gamma}}=b_{n}^{-\gamma}\left(1+o_{P}(1)\right)$. So the LHS of (6)
equals $b_{n}^{-\gamma} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=b_{n}^{-\hat{\gamma}}\left(1+o_{P}(1)\right) \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=b_{n}^{-\hat{\gamma}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)+$ $o_{P}\left(b_{n}^{-\hat{\gamma}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)\right)=b_{n}^{-\hat{\gamma}} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)+o_{P}(1)$, since $b_{n}^{\gamma} \bar{K}_{b_{n}}^{-1}\left(t \mid X^{n}\right)=O_{P}(1)$ by (6).

Proof of Theorem 2 Note that since $b_{i, n}=n^{\beta_{i}}, 1>\beta_{1}>\cdots>\beta_{I}>0$, it follows that $b_{i, n} \rightarrow \infty$ and $b_{i, n} / n \rightarrow 0$ for all $i$. Therefore (3) is true, and Lemma 1 validates the ANOVA equation (9), and in particular the fact that $u_{i, j}=o_{P}(1)$. Thus (10) yields

$$
\gamma_{I, J}=\gamma+\frac{\sum_{i=1}^{I}\left(u_{i, .}-\bar{u}\right)\left(\log \left(b_{i, n}\right)-\overline{\log }\right)}{\sum_{i=1}^{I}\left(\log \left(b_{i, n}\right)-\overline{\log }\right)^{2}}
$$

where $\bar{u}=(I J)^{-1} \sum_{i=1}^{I} \sum_{j=1}^{J} u_{i, j}$, and $u_{i, .}=J^{-1} \sum_{j=1}^{J} u_{i, j}$. As in Bertail et al. (1999), note that $\overline{\log }=A \log n$, and $\sum_{i=1}^{I}\left(\log \left(b_{i, n}\right)-\overline{\log }\right)^{2}=B(\log n)^{2}$, for some constants $A, B$; thus, for fixed $I$ and $J$, we have $\gamma_{I, J}=\gamma+o_{P}\left((\log n)^{-1}\right)$.

Proof of Theorem 3 Fix an $x$ and note that

$$
K_{b_{n}}\left(x \mid X^{n}, \widehat{\tau} .\right) \equiv q^{-1} \sum_{i=1}^{q} 1\left\{b_{n}^{-\hat{\gamma}} T_{b_{n}, i} \leq x\right\}
$$

As before, since $\hat{\gamma}=\gamma+o_{P}\left((\log n)^{-1}\right)$ it follows that $b_{n}^{-\hat{\gamma}} / b_{n}^{-\gamma} \rightarrow 1$ in probability, and $b_{n}^{-\hat{\gamma}}=b_{n}^{-\gamma}\left(1+o_{P}(1)\right)$. Therefore, $b_{n}^{-\hat{\gamma}} T_{b_{n}, i}=b_{n}^{-\gamma} T_{b_{n}, i}+o_{P}\left(b_{n}^{-\gamma} T_{b_{n}, i}\right)=$ $b_{n}^{-\gamma} T_{b_{n}, i}+o_{P}(1)$, since $b_{n}^{-\gamma} T_{b_{n}, i}=O_{P}(1)$ by (1).

So from (3) of Theorem 1 equation (11) follows. Equation (12) follows by similar arguments as in the proof of Theorem 5 in Bertail et al. (1999).

Proof of Theorem 5 Equation (15) directly yields that $\widehat{\tau}_{b_{n}}=\tau_{b_{n}}(1+$ $\left.o_{P}(1)\right) . ¿$ From (14) we get that at any $\mathrm{t}_{i}$

$$
\begin{equation*}
\bar{K}_{U, b_{n}}^{-1}\left(t_{i} \mid X^{n}\right)=\mu_{b_{n}+} \widehat{\tau}_{b_{n}}^{-1} K^{-1}\left(t_{i}, P\right)\left(1+o_{P}(1)\right)+o_{P}\left(\tau_{b_{n}}^{-1}\right) \tag{29}
\end{equation*}
$$

It follows that
$\widehat{\mu}_{b_{n}, I}=I^{-1} \sum_{i=1}^{I}\left(K_{U, b_{n}}^{-1}\left(t_{i} \mid X^{n}\right)\right)-\widehat{\tau}_{b_{n}}^{-1} K^{-1}\left(t_{i} \mid X^{n}\right)=\mu_{b_{n}}+I^{-1} \sum_{i=1}^{I} K^{-1}\left(t_{i}, P\right) o_{P}\left(\tau_{b_{n}}^{-1}\right)$

Since $K\left(t_{i}, P\right)$ is continuous at each $\mathrm{t}_{i}$ and I fixed , $I^{-1} \sum_{i=1}^{I} K^{-1}\left(t_{i}, P\right)$ is bounded and the result follows.

Proof of Theorem 7 The proof for $\widehat{\hat{\gamma}}_{I, J}=\gamma+o_{P}\left(\log n^{-1}\right)$ and $\widehat{C}_{0}=$ $C_{0}+o(1)$ is similar to the proof of theorem 2. It essentially remains to prove the statement on $\widehat{L}$. By construction of the $b_{n, i}^{*}$, we have $n^{1 / 2} \leq b_{n, j}^{*}<e^{j} n^{1 / 2}$, $j=1, \ldots . M_{n}$. Since $M_{n}=o(\log n), \mathrm{b}_{n, M_{n}}^{*} / n \rightarrow 0$ and all the $b_{n, j}^{*}$ satisfy the hypothesis of Theorem 1b and are thus admissible for the construction of the subsampling distributions. Since $\widehat{\hat{\gamma}}_{I, J}=\gamma+o_{P}\left(\log n^{-1}\right)$ holds and since we have fixed $y\left(t_{1,0}, t_{2,0}, \infty\right)=0$, we have directly from (17)

$$
\widehat{h}_{1}\left(\log \left(b_{n}^{*}\right)\right)=h_{1}\left(\log b_{n}^{*}\right)+o_{P}(1)
$$

It follows that for any $x \in\left[b_{n, i}^{*}, b_{n, i+1}^{*}\right]$,

$$
\begin{gathered}
\widehat{L}(x)=\left\{L\left(b_{n, i}^{*}\right)\left(1+o_{p}(1)\right)+\left(L\left(b_{n, i+1}^{*}\right)\left(1+o_{p}(1)\right)\right.\right. \\
\left.\left.-L\left(b_{n, i}^{*}\right)\left(1+o_{p}(1)\right)\right) \int_{0}^{\left(x-b_{n, i}^{*}\right) /\left(b_{n, i+1}^{*}-b_{n, i}^{*}\right)} s(t) d t+o_{P}(1)\right\} .
\end{gathered}
$$

But according to Bingham et al. (1987, p. 15 -see the proof of their theorems 1.3.3 and 1.3.4), uniformly in $x \in\left[b_{n, i}^{*}, b_{n, i+1}^{*}\right]$,

$$
L_{1}(x)=L\left(b_{n, i}^{*}\right)+\left(L\left(b_{n, i+1}^{*}\right)-L\left(b_{n, i}^{*}\right)\right) \int_{0}^{\left(x-b_{n, i}^{*}\right) /\left(b_{n, i+1}^{*}-b_{n, i}^{*}\right)} s(t) d t
$$

is asymptotically equivalent to $L(x)$. Since $s(t)$ is uniformly bounded on each interval, it follows that

$$
\widehat{L}(x)=\left(L(x)+o_{P}(1)+r_{n}\right)
$$

with

$$
r_{n}=\left(L\left(b_{n, i}^{*}\right) o_{p}(1)-L\left(b_{n, i+1}^{*}\right) o_{p}(1)\right)
$$

But by the uniform convergence theorem (see theorem 1.2.1 of Bingham et al. (1987)), $L\left(b_{n, i}^{*}\right) / L(x) \rightarrow 1, L\left(b_{n, i+1}^{*}\right) / L(x) \rightarrow 1$ and $L\left(b_{n, i+1}^{*}\right) / L\left(b_{n, i}^{*}\right) \rightarrow 1$ uniformly in $x \in\left[b_{n, i}^{*}, b_{n, i+1}^{*}\right]$. We deduce from this that $r_{n}=L(x) o_{p}(1)$ and thus that

$$
\widehat{L}(x)=L(x)\left(1+o_{P}(1)\right)
$$

yielding the result of the theorem.
The last result of the theorem follows from the fact that
$\widehat{L}\left(n^{1 / 2}\right)=L\left(n^{1 / 2}\right)(1+o(1))=L(n) \exp \left(C_{0}\right)(1+o(1))=L(n) \exp \left(C_{0}\right)(1+o(1)$.
This combined with $\widehat{\hat{\gamma}}_{I, J}=\gamma+o_{P}\left(\log n^{-1}\right)$ and $\widehat{C}_{0}=C_{0}+o(1)$ yields the convergence of the estimated rate.


Figure 1: Optimal Block Lengths for 1\% Value-at-Risk


Figure 2: Optimal Block Lengths for 5\% Value-at-Risk


Figure 3: In Sample Exceedances for 1\% Value-at-Risk


Figure 4: In Sample Exceedances for 5\% Value-at-Risk


Figure 5: Estimated Divergence Rates for 1\% Value-at-Risk


Figure 6: Estimated Divergence Rates for 5\% Value-at-Risk


Figure 7: Estimates for 1 Day Value-at-Risk


Figure 8: Estimates for 10 Day Value-at-Risk


Figure 9: Estimates for 23 Day Value-at-Risk


Figure 10: Portfolio returns for $s=0.1 \%$.


Figure 11: Portfolio returns for $s=0.5 \%$.


Figure 12: Portfolio returns for $s=1 \%$.

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[^0]:    *INRA-CORELA, 65 Bd. de Brandebourg, 94205 Ivry/Seine, France.
    ${ }^{\dagger}$ Department of Economics, University of California at San Diego, La Jolla, CA 92093 USA.
    $\ddagger$ Department of Mathematics, University of California at San Diego, La Jolla, CA 92093 USA.

[^1]:    ${ }^{1}$ The case where $T_{n} \xrightarrow{P}-\infty$ can be similarly treated by noting that in that case $-T_{n} \xrightarrow{P}$ $\infty$. Note that $\xrightarrow{P}$ indicates convergence in probability; the order notations 'in probability',

[^2]:    ${ }^{2}$ Note, however, that the effect of the centering $\mu_{n}$ will be negligible if $\tau_{n} \mu_{n} \rightarrow 0$, e.g., when the standard deviation of the general statistic $U_{n}$ increases faster than its mean $E U_{n}$; this is actually what happens in the case of the maximum of heavy-tailed observations of our Example 3 - see also section 3.1 in what follows.

[^3]:    ${ }^{3}$ Note that there is no real discrepancy between eq. (22) and the type II extreme distribution of our Section 3.1; the only difference is that now we are directly interested in the scaling constant $\Gamma$ as well, and thus we cannot allow the presence of an arbitrary 'floating' proportionality constant in the rate $\tau_{n}$. In other words, whereas in eq. (22) (and throughout Section 4) we explicitly take $\tau_{n}=n^{-1 / \alpha}$, the rate $\tau_{n}$ that figures in the type II extreme distribution of our Section 3.1 is proportional (but not equal) to $n^{-1 / \alpha}$; as a matter of fact, under eq. (21), the rate $\tau_{n}$ that figures in the type II extreme distribution of Section 3.1 satisfies $\tau_{n}=(\Gamma n)^{-1 / \alpha}$ instead.

[^4]:    ${ }^{4}$ We compute $Q L L=-n^{-1} \sum_{i=1}^{n}\left|y_{i}-\widehat{V a R}_{i}\right| \cdot\left(p 1_{\left[y_{i} \geq \widehat{V a R}_{i}\right]}+(1-p) 1_{\left[y_{i}<\widehat{V a R}_{i}\right]}\right)$. This measure weights the observed deviation from the $V a R$ with the probability with which it is supposed to occur.

