

# Nonparametric estimation of conditional Value-at-Risk and Expected Shortfall based on Extreme Value Theory

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# Preliminaries

- ▶ Understanding and modeling price volatility
- ▶ Pushing the research frontier
- ▶ Identifying periods of excessive price volatility

In empirical finance there is often an interest in stochastic models for log returns

$$Y_t = \log \frac{P_t}{P_{t-1}} \text{ where } t \in \{0, \pm 1, \dots\}.$$

# Motivation

- ▶  $\{Y_t\}_{t \in \mathbb{Z}}$  be a stochastic process
- ▶  $F_{Y_t | \mathbf{x}_t = \mathbf{x}}, \mathbf{x}_t \in \mathbb{R}^d$ .
- ▶ Normally,  $\mathbf{x}'_t = ( Y_{t-1} \ \cdots \ Y_{t-m} \ W'_t )$  for  $m \in \mathbb{N}$ .

For  $a \in (0, 1)$ ,

- ▶  $a\text{-CVaR}(\mathbf{x})$  is the  $a$ -quantile associated with  $F_{Y_t | \mathbf{x}_t = \mathbf{x}}$ ,
- ▶  $a\text{-CES}(\mathbf{x})$  is the  $E(Y_t | Y_t > a\text{-CVaR}(\mathbf{x}))$ .

These are frequently used as synthetic measures of risk by regulators, portfolio managers, etc.

# How does $Y_t$ evolve through time?

We consider the following conditional location-scale model

$$Y_t = m(\mathbf{X}_t) + h^{1/2}(\mathbf{X}_t)\varepsilon_t, \text{ where } t = 1, \dots, n.$$

- ▶  $m, h : \mathbb{R}^d \rightarrow \mathbb{R}$  are suitably restricted real valued functions
- ▶  $E(\varepsilon_t | \mathbf{X}_t = \mathbf{x}) = 0$  and  $V(\varepsilon_t | \mathbf{X}_t = \mathbf{x}) = 1$
- ▶  $\varepsilon_t$  has a strictly increasing absolutely continuous distribution  $F$  which belongs to the domain of attraction of an extremal distribution [Leadbetter (1983), Resnick (1987)].

## A result of Gnedenko (1943)

- ▶ Let  $\{X_t\}_{t \geq 1}$  be a sequence of iid random variables with distribution  $F$  and let  $m_n = \max\{X_1, \dots, X_n\}$ . Then,

$$P(m_n \leq x) = P(X_t \leq x, \forall t) = F(x)^n$$

Suppose there exists  $a_n > 0$ ,  $b_n \in \mathbb{R}$  such that as  $n \rightarrow \infty$

$$P\left(\frac{m_n - b_n}{a_n} \leq x\right) = F(a_n x + b_n)^n \rightarrow E(x)$$

then  $E(x)$  is either

1.  $\Phi_\alpha(x) = e^{-x^{-\alpha}}$  for  $x \geq 0$  (Fréchet)
  2.  $\Psi_\alpha(x) = e^{-(-x)^\alpha}$  for  $x < 0$  (reverse Weibull) and 1 for  $x \geq 0$ ,
  3.  $\Lambda(x) = e^{-e^{-x}}$  for  $x \in \mathbb{R}$  (Gumbel).
- ▶ There are  $F$ 's that are not in the domain of attraction of  $E$  but they constitute rather pathological cases (Leadbetter et al., 1983).

## Some restrictions of the location scale model

1. AR(m), ARCH(p):  $\mathbf{X}_t = ( 1 \quad Y_{t-1} \quad \cdots \quad Y_{t-m} )$

$$m(\mathbf{X}_t) = \mathbf{X}_t' \mathbf{b}$$

$$h(\mathbf{X}_t) = ( 1 \quad Y_{t-1}^2 \quad \cdots \quad Y_{t-p}^2 ) \mathbf{a}$$

2. CHARN Model of Diebolt and Guègan (1993), Härdle and Tsybakov (1997), Hafner (1998).

$$\mathbf{X}_t = ( Y_{t-1} )$$

$$m(\mathbf{X}_t) = m(Y_{t-1})$$

$$h(\mathbf{X}_t) = h(Y_{t-1})$$

3. Nonparametric autoregression of Fan and Yao (1998, Biometrika)

# CVaR and CES

For  $a \in (0, 1)$ ,

$$a - \text{CVaR}(\mathbf{x}) = q_{Y_t|\mathbf{x}_t=\mathbf{x}}(a) = m(\mathbf{x}) + h^{1/2}(\mathbf{x})q(a)$$

and

$$a\text{-CES}(\mathbf{x}) = E(Y_t | Y_t > q_{Y_t|\mathbf{x}_t=\mathbf{x}}(a)) = m(\mathbf{x}) + h^{1/2}(\mathbf{x})E(\varepsilon_t | \varepsilon_t > q(a))$$

where  $q(a)$  is the  $a$ -quantile associated with  $F$ .

The sequence  $\{\varepsilon_t\}$  is not observed.

Motivation: McNeill and Frey (2000), Martins-Filho and Yao (2006).

## Estimation

Given a sample  $\{(Y_t, \mathbf{X}_t^T)\}_{t=1}^n$  and estimators  $\hat{m}(\mathbf{x})$  and  $\hat{h}(\mathbf{x})$  it is possible to obtain a sequence of standardized nonparametric residuals

$$\hat{\varepsilon}_t = \frac{Y_t - \hat{m}(\mathbf{X}_t)}{\hat{h}^{1/2}(\mathbf{X}_t)} \chi_{\{\hat{h}(\mathbf{X}_t) > 0\}} \text{ for } t = 1, \dots, n,$$

These can be used to construct

$$\hat{q}_{Y|X_i=x}(a) = \hat{m}(\mathbf{x}) + \hat{h}^{1/2}(\mathbf{x})\hat{q}(a)$$

$$\hat{E}(Y_t | Y_t > q_{Y_t|\mathbf{X}_t=\mathbf{x}}(a)) = \hat{m}(\mathbf{x}) + \hat{h}^{1/2}(\mathbf{x})\hat{E}(\varepsilon_t | \varepsilon_t > q(a))$$



## Pickands' result

- ▶ We are interested in the case where  $a$  is in the vicinity of 1.
- ▶ The restriction that  $a$  is in a neighborhood of 1 is useful in estimation. The result is due to Pickands (1975).

$F(x) \in D(E)$  if, and only if,

$$\lim_{\xi \rightarrow u_{\infty}} \sup_{0 < u < u_{\infty} - \xi} |F_{\xi}(u) - G(u; 0, \sigma(\xi), k)| = 0,$$

where

- ▶  $F_{\xi}(u) = \frac{F(u+\xi) - F(\xi)}{1 - F(\xi)}$ ,
- ▶  $G$  is a generalized Pareto distribution (GPD), i.e.,

$$G(y; \mu, \sigma, k) = \begin{cases} 1 - (1 - k(y - \mu)/\sigma)^{1/k} & \text{if } k \neq 0, \sigma > 0 \\ 1 - \exp(-(y - \mu)/\sigma) & \text{if } k = 0, \sigma > 0 \end{cases}$$

with  $\mu \leq y < \infty$  if  $k \leq 0$ ,  $\mu \leq y \leq \mu + \sigma/k$  if  $k > 0$

## A restriction on $F$

- ▶ Index of regular variation: If for  $x > 0$ ,  $\lim_{t \rightarrow \infty} \frac{1-F(tx)}{1-F(t)} = x^\alpha$  we say that  $1 - F$  is regularly varying at  $\infty$  with index  $\alpha$ . If  $\alpha = 0$  we say that  $1 - F$  is slowly varying at  $\infty$ .
- ▶  $F \in \Phi_\alpha(x) \Leftrightarrow \lim_{t \rightarrow \infty} \frac{1-F(tx)}{1-F(t)} = x^{-\alpha} \Leftrightarrow x^\alpha(1 - F(x))$  is slowly varying at  $\infty$ .
- ▶ If  $F \in D(\Psi_\alpha)$  its endpoint  $u_\infty$  is finite.
- ▶ If  $F \in D(\Lambda)$  and its endpoint  $u_\infty$  is not finite,  $1 - F$  is rapidly varying, a situation we will (must?) avoid.
- ▶ If  $F$  belongs to the domain of attraction of a Fréchet distribution ( $\Phi_\alpha$ ) with parameter  $\alpha$ , then  $k = -\frac{1}{\alpha}$  and  $\sigma(\xi) = \xi/\alpha$ .
- ▶ An estimator for  $q(a)$  can be obtained from the estimation of the parameters  $k$  and  $\sigma(\xi)$ .

## Estimation procedure

First stage: a) We consider the local linear (LL) estimator

$\hat{m}(\mathbf{x}) \equiv \hat{\beta}_0$  where

$$(\hat{\beta}_0, \hat{\beta}) \equiv \underset{\beta_0, \beta}{\operatorname{argmin}} \sum_{t=1}^n \left( Y_t - \beta_0 - (\mathbf{X}_t^T - \mathbf{x}^T) \beta \right)^2 K_1 \left( \frac{\mathbf{X}_t - \mathbf{x}}{h_{1n}} \right),$$

$K_1(\cdot)$  is a multivariate kernel function and  $h_{1n} > 0$  is a bandwidth.

b) We obtain  $\{\hat{U}_t \equiv Y_t - \hat{m}(\mathbf{X}_t)\}_{t=1}^n$  and define  $\hat{h}(\mathbf{x}) \equiv \hat{\eta}$  where

$$(\hat{\eta}, \hat{\eta}_1) \equiv \underset{\eta, \eta_1}{\operatorname{argmin}} \sum_{t=1}^n \left( \hat{U}_t^2 - \eta - (\mathbf{X}_t^T - \mathbf{x}^T) \eta_1 \right)^2 K_2 \left( \frac{\mathbf{X}_t - \mathbf{x}}{h_{2n}} \right),$$

$K_2(\cdot)$  is a multivariate kernel function and  $h_{2n} > 0$  is a bandwidth.

# Estimation procedure

Second stage:

- ▶ We use  $\{\hat{\varepsilon}_t\}_{t=1}^n$  to estimate  $F$  as

$$\tilde{F}(u) = \frac{1}{nh_{3n}} \sum_{t=1}^n \int_{-\infty}^u K_3 \left( \frac{\hat{\varepsilon}_t - y}{h_{3n}} \right) dy \quad (1)$$

where  $K_3(\cdot)$  is a univariate kernel and  $h_{3n} > 0$  is a bandwidth.

- ▶ Let  $\tilde{q}(a)$  be the solution for  $\tilde{F}(\tilde{q}(a)) = a$ . Letting  $0 < a_n < a < 1$  be such that  $a_n \rightarrow 1$  as  $n \rightarrow \infty$  we use  $N_s$  residuals that exceed  $\tilde{q}(a_n)$  to form

$$\{\tilde{Z}_i\}_{i=1}^{N_s} = \{\hat{\varepsilon}_{(n-N_s+i)} - \tilde{q}(a_n)\}_{i=1}^{N_s}$$

where  $\{\hat{\varepsilon}_{(t)}\}_{t=1}^n$  denotes the order statistics associated with  $\{\hat{\varepsilon}_t\}_{t=1}^n$ .

## Estimation procedure

- ▶  $\{\tilde{Z}_i\}_{i=1}^{N_s}$  is used to obtain maximum likelihood estimators for  $\sigma$  and  $k$  based on  $g(z; \sigma, k) = \frac{1}{\sigma} \left(1 - \frac{kz}{\sigma}\right)^{1/k-1}$ . That is,

$$\frac{\partial}{\partial \sigma} \frac{1}{N_s} \sum_{i=1}^{N_s} \log g(\tilde{Z}_i; \tilde{\sigma}_{\tilde{q}(a_n)}, \tilde{k}) = 0 \quad (2)$$

$$\frac{\partial}{\partial k} \frac{1}{N_s} \sum_{i=1}^{N_s} \log g(\tilde{Z}_i; \tilde{\sigma}_{\tilde{q}(a_n)}, \tilde{k}) = 0. \quad (3)$$

- ▶ Based on Pickands approximation

$$F_{\tilde{q}(a_n)}(y) = \frac{F(y + \tilde{q}(a_n)) - F(\tilde{q}(a_n))}{1 - F(\tilde{q}(a_n))} \approx 1 - \left(1 - \frac{ky}{\sigma_{\tilde{q}(a_n)}}\right)^{1/k}$$

## Estimation procedure

- ▶ For  $a \in (a_n, 1)$ ,  $q(a) = \tilde{q}(a_n) + y_{\tilde{q}(a_n),a}$  where by construction  $F(\tilde{q}(a_n) + y_{\tilde{q}(a_n),a}) = a$ .

Then,

$$\frac{1-a}{1-F(\tilde{q}(a_n))} \approx \left(1 - \frac{k y_{\tilde{q}(a_n),a}}{\sigma_{\tilde{q}(a_n)}}\right)^{1/k}.$$

- ▶ If  $F$  is estimated by  $\tilde{F}$ , and noting that  $1 - \tilde{F}(\tilde{q}(a_n)) = 1 - a_n$ , we have

$$y_{\tilde{q}(a_n),a} \approx \frac{\sigma_{\tilde{q}(a_n)}}{k} \left(1 - \left(\frac{1-a}{1-a_n}\right)^k\right).$$

- ▶ Thus,

$$\hat{q}(a) = \tilde{q}(a_n) + \frac{\tilde{\sigma}_{\tilde{q}(a_n)}}{\tilde{k}} \left(1 - \left(\frac{1-a}{1-a_n}\right)^{\tilde{k}}\right).$$

## Estimation procedure

- ▶ If  $\varepsilon_t - q(a)$  were distributed *exactly* as  $g(z; \sigma, k)$ , then integration by parts gives

$$E(\varepsilon_t | \varepsilon_t > q(a)) = q(a) + \frac{\sigma}{1+k}$$

- ▶ If  $F \in D(\Phi_\alpha)$  and restrictions are placed on how fast  $x^\alpha(1 - F(x))$  is slowly varying, it is easy to show that

$$E(\varepsilon_t | \varepsilon_t > q(a)) = \frac{q(a)}{1+k} + o(1)$$

This motivates the estimator

$$\hat{E}(\varepsilon_t | \varepsilon_t > q(a)) = \frac{\hat{q}(a)}{1 + \tilde{k}}.$$

# Assumptions

There are a number of regularity conditions.

- ▶ Assumption 1 deals with the properties of the three kernels used in estimation. The order  $s$  for  $K_1$  and  $K_2$  are needed to establish that the biases for  $\hat{m}$  and  $\hat{h}$  are, respectively, of order  $O(h_{in}^s)$  for  $i = 1, 2$  in Lemmas 2 and 3. The order  $m_1$  for  $K_3$  is necessary in the proof of Lemma 4.
- ▶ Assumption 3 restricts  $m$  and  $h$  to be sufficiently smooth.
- ▶ Assumption A5 is necessary in Lemma 4 and is directly related to the verification of existence of bounds required to use Lemma A.2 in Gao (2007).



# Assumptions

- ▶ Assumption 2 is important:

1)  $\{(\mathbf{X}_t \ \varepsilon_t)^T\}_{t=1,2,\dots}$  is a strictly stationary  $\alpha$ -mixing process with  $\alpha(l) \leq C l^{-B}$  for some  $B > 2$ ;

2) The joint density of  $\mathbf{X}_t$  and  $\varepsilon_t$  is given by

$$f_{\mathbf{X}\varepsilon}(\mathbf{x}, \varepsilon) = f_{\mathbf{X}}(\mathbf{x})f(\varepsilon);$$

3)  $f_{\mathbf{X}}(\mathbf{x})$  and all of its partial derivatives of order  $< s$  are differentiable and uniformly bounded on  $\mathbb{R}^d$ ;

4)  $0 < \inf_{\mathbf{x} \in \mathcal{G}} f_{\mathbf{X}}(\mathbf{x})$  and  $\sup_{\mathbf{x} \in \mathcal{G}} f_{\mathbf{X}}(\mathbf{x}) \leq C$ .

A2 1) implies that for some  $\delta > 2$  and  $a > 1 - \frac{2}{\delta}$ ,

$\sum_{j=1}^{\infty} j^a \alpha(j)^{1-\frac{2}{\delta}} < \infty$ , a fact that is needed in our proofs. We

note that  $\alpha$ -mixing is the weakest of the mixing concepts

Doukhan (1994) and its use here is only possible due to

Lemma A.2 in Gao (2007), which plays a critical role in the proof of Lemma 4.

## Lemmata

Lemma 2: Assume that the kernel  $K_1$  used to define  $\hat{m}$  satisfies assumption A1 and assumptions A2 and A3 are holding. Assume also that the bandwidth  $h_{1n}$  used to define  $\hat{m}$  satisfies equations

$$n^{1-\frac{2}{a}-2\theta} h_n^d \rightarrow \infty$$

and

$$n^{(B+1.5)(\frac{1}{a}+\theta)-\frac{B}{2}+0.75+\frac{d}{2}} h_n^{-1.75d-\frac{d}{2}(d+B)} (\log n)^{0.25+0.5(B-d)} \rightarrow 0.$$

Then, if  $E(|\varepsilon_t|^a) < \infty$ ,  $E(h^{1/2}(\mathbf{X}_t)^a) < \infty$  for some  $a > 2$  and condition c) in Lemma 1 is holding

$$\sup_{\mathbf{x} \in \mathcal{G}} |\hat{m}(\mathbf{x}) - m(\mathbf{x})| = O_p(L_{1n}), \quad (4)$$

where  $L_{1n} = \left(\frac{\log n}{nh_{1n}^d}\right)^{1/2} + h_{1n}^s$ .

## Lemmata

Lemma 3: Assume that the kernel  $K_2$  used to define  $\hat{h}$  satisfies assumption A1 and assumptions A2 and A3 are holding. Assume also that the bandwidth  $h_{2n}$  used to define  $\hat{h}$  satisfies equations

$$n^{1-\frac{2}{a}-2\theta} h_n^d \rightarrow \infty$$

and

$$n^{(B+1.5)(\frac{1}{a}+\theta)-\frac{B}{2}+0.75+\frac{d}{2}} h_n^{-1.75d-\frac{d}{2}(d+B)} (\log n)^{0.25+0.5(B-d)} \rightarrow 0.$$

Then, under the assumptions in Lemma 2, if  $E(|\varepsilon_t^2 - 1|^a) < \infty$  and  $E(h(\mathbf{X}_t)^a) < \infty$  for some  $a > 2$ ,

$$\sup_{\mathbf{x} \in \mathcal{G}} |\hat{h}(\mathbf{x}) - h(\mathbf{x})| = O_p(L_{1n} + L_{2n}), \quad (5)$$

where  $L_{1n} = \left(\frac{\log n}{nh_{1n}^d}\right)^{1/2} + h_{1n}^s$  and  $L_{2n} = \left(\frac{\log n}{nh_{2n}^d}\right)^{1/2} + h_{2n}^s$ .

# Lemmata

Corollary: Under the assumptions of Lemma 3,

$$\sup_{\mathbf{x} \in \mathcal{G}} |\hat{h}^{1/2}(\mathbf{x}) - h^{1/2}(\mathbf{x})| = O_p(L_{1n} + L_{2n})$$

and

$$\sup_{\mathbf{x} \in \mathcal{G}} |\chi_{\{\hat{h}(\mathbf{x}) > 0\}} - 1| = O_p(L_{1n} + L_{2n}),$$

where  $L_{1n} = \left(\frac{\log n}{nh_{1n}^d}\right)^{1/2} + h_{1n}^s$  and  $L_{2n} = \left(\frac{\log n}{nh_{2n}^d}\right)^{1/2} + h_{2n}^s$ .

► FR1': Assume that for some  $\alpha > 0$  we have

$$\lim_{x \rightarrow \infty} \frac{xf(x)}{1-F(x)} = \alpha.$$

# Lemmata

Lemma 4: Under assumptions A1-A5 and conditions FR1' and FR2, if  $\alpha \geq 1$  we have

$$N^{1/2} \left( \frac{\tilde{q}(a_n) - q_n(a_n)}{q(a_n)} \right) = O_p(1), \text{ where } a_n = 1 - \frac{N}{n}.$$

provided that a)  $h_{1n} \propto n^{-\frac{1}{2s+d}}$ ,  $h_{2n} \propto n^{-\frac{1}{2s+d}}$ ,  $h_{3n} \propto n^{-\frac{s}{2(2s+d)} + \delta}$ ,  $N \propto n^{\frac{2s}{2s+d} - \delta}$  for some  $\delta > 0$  with  $s \geq 2d$  and b)  $E(|\varepsilon_t^2 - 1|^a) < \infty$  and  $E(h(\mathbf{x})^a) < \infty$  for some  $a > 2$ .

# Theorem 1

Assume that FR1' with  $\alpha > 1$ , FR2 and assumptions A1-A5 are holding. In addition, assume that a)  $h_{1n} \propto n^{-\frac{1}{2s+d}}$ ,  $h_{2n} \propto n^{-\frac{1}{2s+d}}$ ,  $h_{3n} \propto n^{-\frac{s}{2(2s+d)} + \delta}$ ,  $N \propto n^{\frac{2s}{2s+d} - \delta}$  for some  $\delta > 0$  and  $s \geq 2d$ , b)  $E(|\varepsilon_t^2 - 1|^a) < \infty$  and  $E(h(\mathbf{x})^a) < \infty$  for some  $a > 2$ . Let  $\tau_1, \tau_2 \in \mathbb{R}$ ,  $0 < \delta_N \rightarrow 0$ ,  $\delta_N N^{1/2} \rightarrow \infty$  as  $N \rightarrow \infty$  and denote arbitrary  $\sigma$  and  $k$  by  $\sigma = \sigma_N(1 + \tau_1 \delta_N)$  and  $k = k_0 + \tau_2 \delta_N$ . We define the log-likelihood function

$$\tilde{L}_{TN}(\tau_1, \tau_2) = \frac{1}{N} \sum_{i=1}^{N_s} \log g(\tilde{Z}_i; \sigma_N(1 + \tau_1 \delta_N), k_0 + \tau_2 \delta_N),$$

where  $\tilde{Z}_i = \hat{\varepsilon}_{(n-N_s+i)} - \tilde{q}(a_n)$ ,  $a_n = 1 - \frac{N}{n}$ . Then, as  $n \rightarrow \infty$ ,  $\frac{1}{\delta_N^2} \tilde{L}_{TN}(\tau_1, \tau_2)$  has, with probability approaching 1, a local maximum  $(\tau_1^*, \tau_2^*)$  on  $S_T = \{(\tau_1, \tau_2) : \tau_1^2 + \tau_2^2 < 1\}$  at which  $\frac{1}{\delta_N} \frac{\partial}{\partial \tau_1} \tilde{L}_{TN}(\tau_1^*, \tau_2^*) = 0$  and  $\frac{1}{\delta_N} \frac{\partial}{\partial \tau_2} \tilde{L}_{TN}(\tau_1^*, \tau_2^*) = 0$ .

# Comments on Theorem 1

- ▶ The vector  $(\tau_1^*, \tau_2^*)$  implies a value  $\tilde{\sigma}_{\tilde{q}(a_n)}$  and  $\tilde{k}$  which are solutions for the likelihood equations

$$\frac{\partial}{\partial \sigma} \frac{1}{N} \sum_{j=1}^{N_s} \log g(\tilde{Z}_j; \tilde{\sigma}_{\tilde{q}(a_n)}, \tilde{k}) = 0, \quad \frac{\partial}{\partial k} \frac{1}{N} \sum_{j=1}^{N_s} \log g(\tilde{Z}_j; \tilde{\sigma}_{\tilde{q}(a_n)}, \tilde{k}) = 0.$$

- ▶ There exists, with probability approaching 1, a local maximum  $(\tilde{\sigma}_N = \sigma_N(1 + t^* \delta_N), \tilde{k} = k_0 + \tau^* \delta_N)$  on  $S_R = \{(\sigma, k) : \|(\frac{\sigma}{\sigma_N} - 1, k - k_0)\| < \delta_N\}$  that satisfy the first order conditions
- ▶ Theorem 1 states that the solutions for the first order conditions correspond to a local maximum of the likelihood associated with the GPD in a shrinking neighborhood of the arbitrary point  $(\sigma_N, k_0)$ .

# Comments on Theorem 1

The proof of Theorem 1 depends critically on:



$$\sup_{\mathbf{x} \in \mathcal{G}} |\hat{m}(\mathbf{x}) - m(\mathbf{x})| = O_p(L_{1n}) \quad \text{and} \quad \sup_{\mathbf{x} \in \mathcal{G}} |\hat{h}(\mathbf{x}) - h(\mathbf{x})| = O_p(L_{2n}),$$

$$\text{where } L_{1n} = \left(\frac{\log n}{nh_{1n}^d}\right)^{1/2} + h_{1n}^s \quad \text{and} \quad L_{2n} = \left(\frac{\log n}{nh_{2n}^d}\right)^{1/2} + h_{2n}^s.$$

▶ These orders are sufficient for

$$|\hat{\varepsilon}_t - \varepsilon_t| = O_p(L_{1n}) + (O_p(L_{1n}) + O_p(L_{2n}))|\varepsilon_t|$$

uniformly in  $\mathcal{G}$ .

▶ Lemma 4 shows that  $\tilde{q}(a_n)$  is asymptotically close to  $q_n(a_n)$  by satisfying  $\frac{\tilde{q}(a_n) - q_n(a_n)}{q_n(a_n)} = O_p(N^{-1/2})$

It is here that the stochasticity of the threshold ( $\tilde{q}$ ) is handled and FR1', FR2 and  $\alpha > 1$  is used.



## Asymptotic normality of $\tilde{\gamma}' = (\tilde{\sigma}_N, \tilde{k})$ - Theorem 2

Suppose FR1' with  $\alpha > 1$ , FR2, A1-A5 hold and that  $\frac{C}{\alpha-\rho} N^{1/2} \phi(q(a_n)) \rightarrow \mu \in \mathbb{R}$ . In addition, assume that conditions a) and b) in Theorem 1 are holding. Then, the local maximum  $(\tilde{\sigma}_{\tilde{q}(a_n)}, \tilde{k})$  of the GPD likelihood function, is such that for  $k_0 = -\frac{1}{\alpha}$  and  $\sigma_N = \frac{q(a_n)}{\alpha}$

$$\sqrt{N} \begin{pmatrix} \frac{\tilde{\sigma}_{\tilde{q}(a_n)}}{\sigma_N} - 1 \\ k - k_0 \end{pmatrix} \xrightarrow{d} N \left( \begin{pmatrix} \frac{\mu(1-k_0)(1+2k_0\rho)}{1-k_0+k_0\rho} \\ \frac{\mu(1-k_0)k_0(1+\rho)}{1-k_0+k_0\rho} \end{pmatrix}, H^{-1}V_2H^{-1} \right)$$

where  $V_2 = \begin{pmatrix} \frac{k_0^2-4k_0+2}{(2k_0-1)^2} & \frac{-1}{k_0(k_0-1)} \\ \frac{-1}{k_0(k_0-1)} & \frac{2k_0^3-2k_0^2+2k_0-1}{k_0^2(k_0-1)^2(2k_0-1)} \end{pmatrix}$ .

# Comments on Theorem 1

- ▶ It is easy to show that  $H^{-1}V_2V^{-1} - H^{-1}$  is positive definite.
- ▶ Any additional bias resulting from the use of  $\tilde{Z}_i$  is of second order.
- ▶ The fact that  $\tilde{Z}_i$  is *not* iid as  $Z_i$  does not require the use of a CLT for dependent processes as justified in Lemma 5.

## Asymptotic normality of $\hat{q}(a)$ - Theorem 3

Suppose FR1' with  $\alpha > 1$ , FR2, A1-A5 and  $\frac{C}{\alpha - \rho} N^{1/2} \phi(q(a_n)) \rightarrow \mu$  with  $k_0 = -\frac{1}{\alpha}$  and  $\sigma_N = q(a_n)/\alpha$ . In addition, assume that conditions a) and b) in Theorem 1 are holding. Then, if  $n(1 - a) \propto N$ , for some  $z_a > 0$

$$\sqrt{n(1 - a)} \left( \frac{\hat{q}(a)}{q(a)} - 1 \right) \xrightarrow{d}$$

$$N \left( (-k_0) \left( -\frac{(z_a^\rho - 1)\mu(\alpha - \rho)}{\rho} - c_b^T H^{-1} \lim_{n \rightarrow \infty} \sqrt{N} \begin{pmatrix} b_\sigma \\ b_k \end{pmatrix} \right), \right.$$

$$\left. k_0^2 \left( c_b^T H^{-1} V_2 H^{-1} c_b + 2c_b^T \begin{pmatrix} 2 - k_0 \\ 1 - k_0 \end{pmatrix} + 1 \right) \right).$$

where  $c_b^T = ( -k_0^{-1}(z_a^{-1} - 1) \quad k_0^{-2} \log(z_a) + k_0^{-2}(z_a^{-1} - 1) )$ ,  
 $b_\sigma = E \left( \frac{\partial}{\partial \sigma} \log g(Z_i; \sigma_N, k_0) \sigma_N \right)$  and  $b_k = E \left( \frac{\partial}{\partial k} \log g(Z_i; \sigma_N, k_0) \right)$ .

## Asymptotic normality of $\hat{E}(\varepsilon_t | \varepsilon_t > q(a))$ - Theorem 4

Suppose FR1' with  $\alpha > 1$ , FR2, A1-A5 and  $\frac{C}{\alpha - \rho} N^{1/2} \phi(q(a_n)) \rightarrow \mu$  with  $k_0 = -\frac{1}{\alpha}$  and  $\sigma_N = q(a_n)/\alpha$ . In addition, assume that conditions a) and b) in Theorem 1 are holding. Then, if  $n(1-a) \propto N$ , for some  $z_a > 0$

$$\sqrt{n(1-a)} \left( \frac{\hat{E}(\varepsilon_t | \varepsilon_t > q(a))}{\frac{q(a)}{1+k_0}} - 1 \right) \xrightarrow{d} N \left( k_0 \frac{(z_a^\rho - 1)\mu(\alpha - \rho)}{\rho} + k_0 c_b^T H^{-1} \lim_{n \rightarrow \infty} \sqrt{N} \begin{pmatrix} b_\sigma \\ b_k \end{pmatrix} - \frac{1}{1+k_0} \lim_{n \rightarrow \infty} \sqrt{N} \begin{pmatrix} 0 & 1 \end{pmatrix} H^{-1} \begin{pmatrix} b_\sigma \\ b_k \end{pmatrix}, \Sigma \right),$$

where  $c_b$ ,  $b_\sigma$ ,  $b_k$  are as defined in Theorem 3,

$$\Sigma = k_0^2 \left( c_b^T H^{-1} V_2 H^{-1} c_b + 2c_b^T \begin{pmatrix} 2 - k_0 \\ 1 - k_0 \end{pmatrix} + 1 \right) + 2 \frac{k_0}{1+k_0} \eta^T V_3 \theta + \frac{1}{(1+k_0)^2} \theta^T V_1 \theta,$$

with

$$\eta^T = \left( -c_b^T H^{-1} \quad -c_b^T H^{-1} \begin{pmatrix} b_\sigma \\ b_k \end{pmatrix} \quad 1 \right),$$

$$\theta^T = \left( (0 \ 1) H^{-1} \quad (0 \ 1) H^{-1} \begin{pmatrix} b_1 \\ b_2 \end{pmatrix} \quad 0 \right),$$

$$V_3 = \begin{pmatrix} \frac{1}{1-2k_0} & -\frac{1}{(k_0-1)(2k_0-1)} & 0 & 0 \\ -\frac{1}{(k_0-1)(2k_0-1)} & \frac{1}{(k_0-1)(2k_0-1)} & 0 & 0 \\ 0 & 0 & k_0^2 & -k_0 \\ 0 & 0 & -k_0 & 1 \end{pmatrix},$$

$$b_1 = \frac{1-k_0}{k_0(2k_0-1)} \text{ and } b_2 = \frac{1}{k_0^2} \left( \frac{k_0-1}{2k_0-1} - \frac{1}{k_0-1} \right).$$

# Consistency

Given Lemmas 2, 3, Theorems 3 and 4 we have that for all  $a \in (0, 1)$ ,

$$\hat{q}_{Y_t|\mathbf{x}_t=\mathbf{x}}(a) = \hat{m}(\mathbf{x}) + \hat{h}^{1/2}(\mathbf{x})\hat{q}(a) \xrightarrow{P} m(\mathbf{x}) + h^{1/2}(\mathbf{x})q(a) = a\text{-CVaR}(\mathbf{x})$$

and

$$\hat{E}(Y_t | Y_t > q_{Y_t|\mathbf{x}_t=\mathbf{x}}(a)) = \hat{m}(\mathbf{x}) + \hat{h}^{1/2}(\mathbf{x})\hat{E}(\varepsilon_t | \varepsilon_t > q(a)) \xrightarrow{P} m(\mathbf{x}) + h^{1/2}(\mathbf{x})E(\varepsilon_t | \varepsilon_t > q(a)) = a\text{-CES}(\mathbf{x}).$$

## Comments

- ▶ Estimation of  $q_{Y|X=x}(a)$  when  $a$  is in the vicinity of 0 has been considered by Chernozhukov (2005) when  $q_{Y|X=x}(a) = x\beta(a)$ ,  $\beta(a) \in \mathbb{R}^d$ .
- ▶ He provides a complete asymptotic characterization of the quantile regression estimator of  $\beta(a)$ .
- ▶ Here  $q_{Y|X=x}(a)$  is nonparametric, it is in this sense more general than the one considered by Chernozhukov.  $a$  approaches 1 at a speed that is slower than the sample size ( $n(1-a) \propto N \rightarrow \infty$  in Theorem 2).
- ▶ Furthermore, similar to Smith (1987) and Hall (1982), our proofs require the specification of the speed at which the tail  $1 - F(x)$  behaves asymptotically as a power function. Theorem 1 specifies this speed to be proportional to  $\sqrt{N}$ .

TABLE 7 BACKTEST RESULTS FOR  $a$ -CONDITIONAL VALUE-AT-RISK ( $q$ ) AND EXPECTED SHORTFALL( $E$ ) ON  $m - n = 500$  OBSERVATIONS, EXPECTED VIOLATIONS =  $(m - n)(1 - a)$ .  
 $q$ : NUMBER OF VIOLATIONS AND P-VALUE (IN BRACKETS).  
 $E$ : P-VALUE FOR EXCEEDANCE RESIDUALS TO HAVE ZERO MEAN.

	$q$			$E$		
	$a = 0.95$	$a = 0.99$	$a = 0.995$	$a = 0.95$	$a = 0.99$	$a = 0.995$
	EXPECTED VIOLATIONS					
	25	5	2.5			
Maize	18 (.151)	5(1)	2(.751)	0	.161	.735
Rice	29(.412)	4(.653)	2(.751)	0	.081	.248
Soybean	21(.412)	3(.369)	2(.751)	0	.302	.244
Wheatcbot	30(.305)	6(.653)	2(.751)	.001	.339	.273
Wheatkcbt	25(1)	5(1)	2(.751)	0	.082	.239

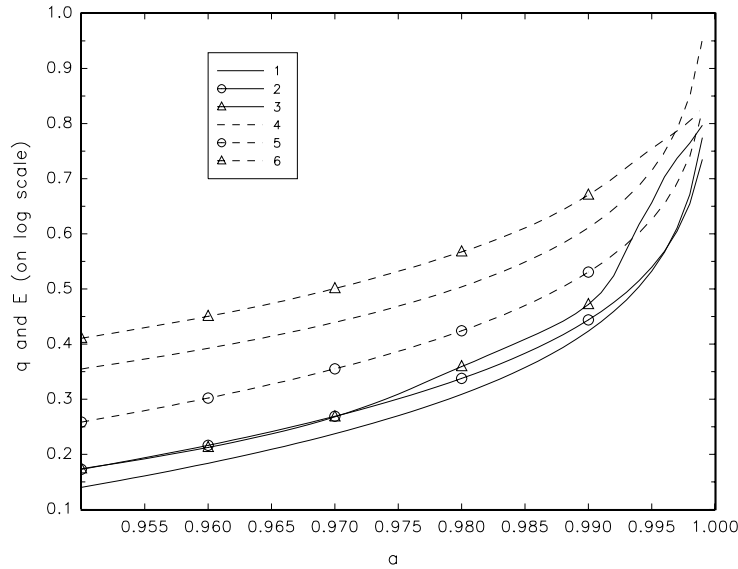


Figure 1: Plot of conditional value-at-risk ( $q$ ) and expected shortfall ( $E$ ) estimates evaluated at sample mean across different  $a$ , with  $n = 1000$ ,  $h_1(Y_{t-1}) = 1 + 0.01Y_{t-1}^2 + 0.5\sin(Y_{t-1})$ ,  $\theta = 0$  and student-t distributed  $\varepsilon_t$  with  $v = 3$ . 1 : true  $q$ , 2 :  $\hat{q}$ , 3 :  $\dot{q}$ , 4 : true  $E$ , 5 :  $\hat{E}$ , and 6 :  $\dot{E}$ .