Asymptotic Properties of the Empirical Spatial Extremogram

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ABSTRACT. The extremogram is a useful tool for measuring extremal dependence and checking model adequacy in a time series. We define the extremogram in the spatial domain when the data is observed on a lattice or at locations distributed as a Poisson point process in d-dimensional space. We establish a central limit theorem for the empirical spatial extremogram. We show these conditions are applicable for max-moving average processes and Brown-Resnick processes and illustrate the empirical extremogram's performance via simulation. We also demonstrate its practical use with a data set related to rainfall in a region in Florida.

Keywords: extremal dependence; extremogram; max moving average; max stable process; spatial dependence

1 Introduction

Extreme events can affect our lives in many dimensions. Events like large swings in financial markets or extreme weather conditions such as floods and hurricanes can cause large financial/property losses and numerous casualties. Extreme events often appear to cluster and that has resulted in a growing interest in measuring *extremal dependence* in many areas including finance, insurance, and atmospheric science.

Extremal dependence between two random vectors X and Y can be viewed as the probability that X is extreme given Y belongs to an extreme set. The *extremogram*, proposed by Davis and Mikosch (2009), is a versatile tool for assessing extremal dependence in a stationary time series. The extremogram has two main features:

- It can be viewed as the extreme-value analog of the autocorrelation function of a stationary time series, i.e., extremal dependence is expressed as a function of lag.
- It allows for measuring dependence between random variables belonging in a large variety of extremal sets. Depending on choices of sets, many of the commonly used extremal dependence measures right tail dependence, left tail dependence, or dependence among large absolute values can be treated as a special case of the extremogram. The flexibility coming from arbitrary choices of extreme sets have made it especially well suited for time series applications such as high-frequency FX rates (Davis and Mikosch (2009)), cross-sectional stock indices (Davis et al. (2012)), and CDS spreads (Cont and Kan (2011)).

In this paper, we will define the notion of the extremogram for random fields defined on \mathbb{R}^d for some d > 1 and investigate the asymptotic properties of its corresponding empirical estimate. Let $\{X_s, s \in \mathbb{R}^d\}$ be a stationary \mathbb{R}^k -valued random field. For measurable sets $A, B \subset \mathbb{R}^k$ bounded away from $\mathbf{0}$, we define the spatial extremogram as

$$\rho_{AB}(h) = \lim_{x \to \infty} P(X_h \in xB | X_0 \in xA), \quad h \in \mathbb{R}^d,$$
(1.1)

provided the limit exists. We call (1.1) the *spatial extremogram* to emphasize that it is for a random field in \mathbb{R}^d . If one takes $A = B = (1, \infty)$ in the k = 1 case, then we recover the tail dependence coefficient between X_h and X_0 . For light tailed time series, such as stationary Gaussian processes, $\rho_{AB}(h) = 0$ for $h \neq \mathbf{0}$ in which case there is no extremal dependence. However, for heavy tailed processes in either time or space, $\rho_{AB}(h)$ is often non-zero for many lags $h \neq \mathbf{0}$ and for most choices of sets A and B bounded away from the origin.

We will consider estimates of $\rho_{AB}(h)$ under two different sampling scenarios. In the first, observations are taken on the lattice \mathbb{Z}^d . Analogous to Davis and Mikosch (2009), we define the *empirical spatial extremogram* (ESE) as

$$\hat{\rho}_{AB,m}(h) = \frac{\sum_{s,t \in \Lambda_n, s-t=h} I_{\{a_m^{-1}X_s \in A, a_m^{-1}X_t \in B\}} / n(h)}{\sum_{s \in \Lambda_n} I_{\{a_m^{-1}X_s \in A\}} / \# \Lambda_n},$$
(1.2)

where

- $\Lambda_n = \{1, 2, \dots, n\}^d$ is the d-dimensional cube with side length n,
- $h \in \mathbb{Z}^d$ are observed lags in Λ_n .
- $m = m_n$ is an increasing sequence satisfying $m \to \infty$ and $m/n \to 0$ as $n \to \infty$,
- a_m is a sequence such that $P(|X| > a_m) \sim m^{-1}$,
- n(h) is the number of pairs in Λ_n with lag h, and
- $\#\Lambda_n$ is the cardinality of Λ_n .

In the second case, the data are assumed to come from a stationary random field X_s , where the locations $\{s_1, ..., s_N\}$ are assumed to be points of a homogeneous Poisson point process on $S_n \subset \mathbb{R}^d$. We define the empirical spatial extremogram as a kernel estimator of $\rho_{AB}(h)$, in the spirit of the estimate of autocorrelation in space (see Li et al. (2008)). Under suitable growth conditions on S_n and restrictions on the kernel function, we show that the weighted estimator of $\rho_{AB}(h)$ is consistent and asymptotically normal.

The organization of the paper is as follows: In Section 2, we present the asymptotic properties of the ESE for both cases described above. Section 3 provides examples illustrating the results of Section 2 together with a simulation study demonstrating the performance of the ESE. In Section 4, the spatial extremogram is applied to a spatial rainfall data set in Florida. The proofs of all the results are in the Appendix.

2 Asymptotics of the ESE

2.1 Definitions and notation

Let $\{X_s, s \in I\}$ be a k-dimensional strictly stationary random process where I is either \mathbb{R}^d or \mathbb{Z}^d . For $H = \{h_1, \ldots, h_t\} \subset I$, we use X_H to denote $(X_{h_1}, \ldots, X_{h_t})$. The random field is said to be regularly varying

with index $\alpha > 0$ if for any H, the radial part $||X_H||$ satisfies for all y > 0

(C1)
$$\frac{P(||X_H|| > yx)}{P(||X_H|| > x)} \to y^{-\alpha} \text{ as } x \to \infty,$$

and the angular part $\frac{X_H}{\|X_H\|}$ is asymptotically independent of the radial part $\|X_H\|$ for large values of $\|X_H\|$, i.e., there exists a random vector $\Theta_H \in \mathbb{S}^{tk-1}$, the unit sphere in \mathbb{R}^{tk} with respect to $\|\cdot\|$, such that

(C2)
$$P\left(\frac{X_H}{\|X_H\|} \in \cdot \middle| \|X_H\| > x\right) \xrightarrow{w} P\left(\Theta_H \in \cdot\right) \text{ as } x \to \infty,$$

where \xrightarrow{w} denotes weak convergence. The distribution of $P(\Theta_H \in \cdot)$ is called the *spectral measure* of X_H .

An equivalent definition of regular variation is given as follows. There exists a sequence $a_n \to \infty, \alpha > 0$ and a family of non-null Radon measures (μ_H) on the Borel σ -field of $\mathbb{R}^{tk} \setminus \{\mathbf{0}\}$ such that $nP(a_n^{-1}X_H \in \cdot) \xrightarrow{v} \mu_H(\cdot)$ for $t \geq 1$, where the limiting measure satisfies $\mu_H(y \cdot) = y^{-\alpha}\mu_H(\cdot)$ for y > 0. Here, \xrightarrow{v} denotes vague convergence. Under the regularly varying assumption, one can show that (1.1) is well defined. See Section 6.1 of Resnick (2006) for more details.

2.2 Random fields on a lattice

Let $\{X_s, s \in \mathbb{Z}^d\}$ be a strictly stationary random field and suppose we have observations $\{X_s, s \in \Lambda_n = \{1, ..., n\}^d\}$. Let $d(\cdot, \cdot)$ be a metric on \mathbb{Z}^d . We denote the α -mixing coefficient by

$$\alpha_{j,k}(r) = \sup \left\{ \alpha \left(\sigma(X_s, s \in S), \sigma(X_s, s \in T) \right) : S, T \subset \mathbb{Z}^d, \#S \le j, \#T \le k, d(S, T) \ge r \right\},\$$

where for any two σ -fields \mathcal{A} and \mathcal{B} , $\alpha(\mathcal{A}, \mathcal{B}) = \sup\{|P(A \cap B) - P(A)P(B)| : A \in \mathcal{A}, B \in \mathcal{B}\}$ and for any $S, T \subset \mathbb{Z}^d$, $d(S, T) = \inf\{d(s, t) : s \in S, t \in T\}$.

In order to study asymptotic properties of (1.2), we impose regularly varying and certain mixing conditions on the random field. In particular, we use the big/small block argument: the side length of big blocks, m_n , and the distance between big blocks, r_n , have to be coordinated in the right fashion. To be precise, we assume the following conditions.

(M1) Let B_{γ} be the ball of radius γ centered at 0, i.e., $B_{\gamma} = \{s \in \mathbb{Z}^d : d(\mathbf{0}, s) \leq \gamma\}$, and set $c = \#B_{\gamma}$. For a fixed γ , assume that there exist $m_n, r_n \to \infty$ with $m_n^{2+2d}/n^d \to 0$, $r_n^d/m_n \to 0$ such that

$$\lim_{k \to \infty} \limsup_{n \to \infty} m_n \sum_{l \in \mathbb{Z}^d, k < d(\mathbf{0}, l) \le r_n} P\left(\max_{s \in B_\gamma} |X_s| > \epsilon a_m, \max_{s' \in B_\gamma + l} |X_{s'}| > \epsilon a_m \right) = 0 \quad \text{for} \quad \forall \epsilon > 0, \ (2.1)$$

$$\lim_{n \to \infty} m_n \sum_{l \in \mathbb{Z}^d, r_n < d(\mathbf{0}, l)} \alpha_{c,c}(d(\mathbf{0}, l)) = 0, \tag{2.2}$$

$$\sum_{l \in \mathbb{Z}^d} \alpha_{j_1, j_2}(d(\mathbf{0}, l)) < \infty \quad \text{for} \quad 2c \le j_1 + j_2 \le 4c, \tag{2.3}$$

$$\lim_{n \to \infty} n^{d/2} m_n^{1/2} \alpha_{c,cn^d}(m_n) = 0, \tag{2.4}$$

where a_m satisfies $P(|X| > a_m) \sim \frac{1}{m}$.

Condition (2.1) restricts the joint distributions for exceedance as two sets of points become far apart. Conditions (2.2) - (2.4) impose restrictions on the decaying rate of the mixing functions together with the level of the threshold specified by m_n . These conditions are similar to those in Bolthausen (1982) and Davis and Mikosch (2009).

As in Davis and Mikosch (2009), the ESE $\hat{\rho}_{AB,m}(h)$ is centered by the *Pre-Asymptotic* (PA) extremogram

$$\rho_{AB,m}(h) = \frac{\tau_{AB,m}(h)}{p_m(A)},\tag{2.5}$$

where $\tau_{AB,m}(h) = m_n P(X_0 \in a_m A, X_h \in a_m B)$ and $p_m(A) = m_n P(X_0 \in a_m A)$. Notice that (2.5) is the ratio of the expected value of the numerator and denominator in (1.2).

Theorem 2.1. Suppose a strictly stationary regularly varying random field $\{X_s, s \in \mathbb{Z}^d\}$ with index $\alpha > 0$ is observed on $\Lambda_n = \{1, ..., n\}^d$. For any finite set of non-zero lags H in \mathbb{Z}^d , assume $(\mathbf{M1})$, where $B_{\gamma} \supseteq H$ for some γ . Then

$$\sqrt{\frac{n^d}{m_n}} \left[\hat{\rho}_{AB,m}(h) - \rho_{AB,m}(h) \right]_{h \in H} \xrightarrow{d} N(\boldsymbol{0}, \Sigma),$$

where the matrix Σ in normal distribution is specified in Appendix A.

We present the proof of Theorem 2.1 in Appendix A. Examples of heavy-tailed processes satisfying (M1) are presented in Section 3.

Remark 1. In Theorem 2.1, the pre-asymptotic extremogram $\rho_{AB,m}(h)$ is replaced by the extremogram $\rho_{AB}(h)$ if

$$\lim_{n \to \infty} \sqrt{\frac{n^d}{m_n}} |\rho_{AB,m}(h) - \rho_{AB}(h)| = 0, \quad for \quad h \in H.$$
 (2.6)

2.3 Random fields on \mathbb{R}^d

Now consider the case of a random field defined on \mathbb{R}^d and the sampling locations are given by points of a Poisson process. In this case, we adopt the ideas from Karr (1986) and Li et al. (2008) and use a kernel estimate of the extremogram. For convenience, we restrict our attention to \mathbb{R}^2 . The extension to \mathbb{R}^d is straightforward, but notationally more complex.

Let $\{X_s, s \in \mathbb{R}^2\}$ be a stationary regularly varying random field with index $\alpha > 0$. Suppose N is a homogeneous 2-dimensional Poisson process with intensity parameter ν and is independent of X. Define $N^{(2)}(ds_1, ds_2) = N(ds_1)N(ds_2)I(s_1 \neq s_2)$. Now consider a sequence of compact and convex sets $S_n \subset \mathbb{R}^2$ with Lebesgue measure $|S_n| \to \infty$ as $n \to \infty$. Assume that for each $y \in \mathbb{R}^2$

$$\lim_{n \to \infty} \frac{|S_n \cap (S_n - y)|}{|S_n|} = 1,\tag{2.7}$$

where $S_n - y = \{x - y : x \in S_n\},\$

$$|S_n| = O(n^2), \qquad |\partial S_n| = O(n), \tag{2.8}$$

and ∂S_n denotes the boundary of S_n .

The spatial extremogram in (1.1) is estimated by $\hat{\rho}_{AB,m}(h) = \hat{\tau}_{AB,m}(h)/\hat{p}_m(A)$, where

$$\hat{p}_m(A) = \frac{m_n}{\nu |S_n|} \int_{S_n} I\left(\frac{X_{s_1}}{a_m} \in A\right) N(ds_1), \tag{2.9}$$

$$\hat{\tau}_{AB,m}(h) = \frac{m_n}{\nu^2} \frac{1}{|S_n|} \int_{S_n} \int_{S_n} w_n(h + s_1 - s_2) I\left(\frac{X_{s_1}}{a_m} \in A\right) I\left(\frac{X_{s_2}}{a_m} \in B\right) N^{(2)}(ds_1, ds_2). \tag{2.10}$$

Here $w_n(\cdot) = \frac{1}{\lambda_n^2} w(\frac{\cdot}{\lambda_n})$ is a sequence of weight functions, where $w(\cdot)$ on \mathbb{R}^2 is a positive, bounded, isotropic probability density function and λ_n is the bandwidth satisfying $\lambda_n \to 0$ and $\lambda_n^2 |S_n| \to \infty$. To establish a central limit theorem for $\hat{\rho}_{AB,m}(h)$, we derive asymptotics of the denominator $\hat{p}_m(A)$ and numerator $\hat{\tau}_{AB,m}(h)$. In order to show consistency of $\hat{p}_m(A)$, we assume the following conditions, which are the non-lattice analogs of (2.1) and (2.2).

(M2) There exist an increasing sequence m_n and r_n with $m_n = o(n)$ and $r_n^2 = o(m_n)$ such that

$$\lim_{k \to \infty} \limsup_{n \to \infty} \int_{B[k, r_n]} m_n P(|X_y| > \epsilon a_m, |X_0| > \epsilon a_m) dy = 0 \quad \text{for} \quad \forall \epsilon > 0,$$
(2.11)

$$\lim_{n \to \infty} \int_{\mathbb{R}^2 \setminus B[0, r_n)} m_n \alpha_{1,1}(y) dy = 0, \tag{2.12}$$

$$\int_{\mathbb{R}^2} \tau_{AA}(y) dy < \infty, \tag{2.13}$$

where $B[a, b) = \{s : a \le d(\mathbf{0}, s) < b, s \in \mathbb{R}^2\}$ and $\tau_{AA}(y) = \lim_{n \to \infty} \tau_{AA, m}(y)$.

For a central limit theorem for $\hat{\tau}_{AB,m}(h)$, the following conditions are required.

(M3) Consider a cube $B_n \subset S_n$ with $|B_n| = O(n^{2a})$ and $|\partial B_n| = O(n^a)$ for 0 < a < 1. Assume that there exist an increasing sequence m_n with $m_n = o(n^a)$ and $\lambda_n^2 m_n \to 0$ such that

$$\sup_{n} E\left\{ \sqrt{\frac{|B_n|\lambda_n^2}{m_n}} \left| \hat{\tau}_{AB,m}(h:B_n) - E\hat{\tau}_{AB,m}(h:B_n) \right|^{2+\delta} \right\} \le C_{\delta}, \quad \delta > 0, \ C_{\delta} < \infty, \tag{2.14}$$

where $\hat{\tau}_{AB,m}(h:B_n)$ is the quantity (2.10) with S_n replaced by B_n on the right-hand side. Further assume

$$\int_{\mathbb{R}^2} \tau_{AB}(y) dy < \infty, \quad \int_{\mathbb{R}^2} \alpha_{2,2}(d(\mathbf{0}, y)) dy < \infty, \tag{2.15}$$

and

$$\sup_{l} \frac{\alpha_{l,l}(h)}{l^2} = O(h^{-\epsilon}) \quad \text{for some } \epsilon > 0.$$
 (2.16)

Lastly, the proof requires some smoothness of the random field.

Definition 2.2. A stationary regularly varying random field $\{X_s, s \in \mathbb{R}^d\}$ satisfies a local uniform negligibility condition (LUNC) if for an increasing sequence a_n satisfying $P(|X| > a_n) \sim \frac{1}{n}$ and for all $\epsilon, \delta > 0$, there exists $\delta' > 0$ such that

$$\lim_{n} \sup_{n} nP \left(\sup_{||s|| < \delta'} \frac{|X_s - X_0|}{a_n} > \delta \right) < \epsilon. \tag{2.17}$$

Theorem 2.3. Let $\{X_s, s \in \mathbb{R}^2\}$ be a stationary regularly varying random field with index $\alpha > 0$ satisfying LUNC. Assume N is a homogeneous 2-dimensional Poisson process with intensity parameter ν and is independent of X. Consider a sequence of compact and convex sets $S_n \subset \mathbb{R}^2$ satisfying $|S_n| \to \infty$ as $n \to \infty$.

Assume conditions (M2) and (M3). Then for any finite set of non-zero lags H in \mathbb{R}^2 ,

$$\sqrt{\frac{|S_n|\lambda_n^2}{m_n}} \left[\hat{\rho}_{AB,m}(h) - \rho_{AB,m}(h)\right]_{h\in H} \to N(\boldsymbol{\mathcal{O}}, \Sigma), \tag{2.18}$$

where the matrix Σ is specified in the proof of Theorem 2.1 in Appendix A.

We present the proof of Theorem 2.3 in Appendix B. As in Remark 1, $\rho_{AB,m}(h)$ can be replaced by $\rho_{AB}(h)$ if $\rho_{AB,m}(h)$ converges fast enough.

Remark 2. In (2.18), $\rho_{AB,m}(h)$ can be replaced by $\rho_{AB}(h)$ if

$$\lim_{n \to \infty} \sqrt{\frac{|S_n|\lambda_n^2}{m_n}} |\rho_{AB,m}(h) - \rho_{AB}(h)| = 0 \qquad \text{for } h \in H.$$
 (2.19)

3 Examples

Here we provide two max-stable processes to illustrate the results of Section 2. For background on max-stable processes, see de Haan (1984) and de Haan and Ferreira (2006). In order to check conditions, we need the result from Dombry and Eyi-Minko (2012).

Proposition 3.1 (Dombry and Eyi-Minko (2012)). Suppose $\{X_s, s \in S\}$ is a max-stable random field with unit Fréchet marginals. If S_1 and S_2 are finite or countable disjoint closed subsets of S, and S_1 and S_2 are the respective σ -fields generated by each set, then

$$\beta(S_1, S_2) \le 4 \sum_{s_1 \in S_1} \sum_{s_2 \in S_2} \rho_{(1,\infty)(1,\infty)}(||s_1 - s_2||)$$
(3.1)

where $\beta(\cdot,\cdot)$ is the β -mixing coefficient. We refer to Lemma 2 in Davis et al. (2013).

Notice that (3.1) provides the upper bound for α -mixing coefficient since $2\alpha(S_1, S_2) \leq \beta(S_1, S_2)$. See Bradley (1993).

3.1 Max Moving Average (MMA)

Let $\{Z_s, s \in \mathbb{Z}^2\}$ be an iid sequence of unit Fréchet random variables. The max-moving average (MMA) process is defined by

$$X_t = \max_{s \in \mathbb{Z}^2} w(s) Z_{t-s}, \tag{3.2}$$

where w(s) > 0 and $\sum_{s \in \mathbb{Z}^2} w(s) < \infty$. Note that the summability of $w(\cdot)$ implies the process is well defined. Also, notice that $a_m = O(m)$ since marginal distributions are Fréchet. Consider the Euclidean metric $d(\cdot, \cdot)$ and write $||l|| = d(\mathbf{0}, l)$ for notational convenience. With $w(s) = I(||s|| \le 1)$, the process (3.2) becomes the MMA(1): $X_t = \max_{||s|| \le 1} Z_{t-s}$. Using $A = B = (1, \infty)$, the extremogram for the MMA(1) is then

$$\rho_{AB}(h) = \lim_{n \to \infty} P(X_h > a_{m_n} | X_0 > a_{m_n}) = \begin{cases} 1, & \text{if } ||h|| = 0, \\ 2/5, & \text{if } ||h|| = 1, \sqrt{2}, \\ 1/5, & \text{if } ||h|| = 2, \\ 0, & \text{if } ||h|| > 2. \end{cases}$$
(3.3)

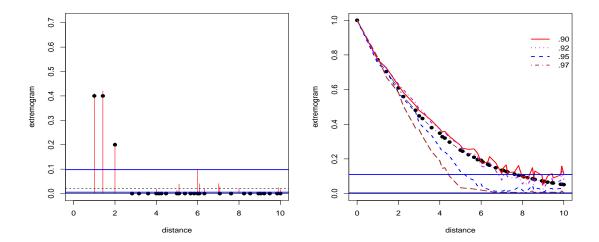


Figure 1: $\rho_{AA}(h)$ and $\hat{\rho}_{AA,m}(h)$, where $A=(1,\infty)$, from a realization of an MMA(1) (left) and the process (3.4) (right). For the ESE, $a_m=.97$ (left) and $a_m=(.90,.92,.95,.97)$ quantile (right) are used. For both cases, the ESE closely tracks the extremogram. Two horizontal lines are 95% random permutation confidence bands.

Since the process is 2-dependent, conditions for Theorem 2.1 are easily checked.

Figure 1 (left) shows $\rho_{AB}(h)$ and $\hat{\rho}_{AB,m}(h)$ from a realization of MMA(1) generated by rmaxstab in the SpatialExtremes package¹ in R. We use 1600 points $(\Lambda_n = \{1, ..., 40\}^2 \in \mathbb{Z}^2)$ and set $A = B = (1, \infty)$ and $a_m = .97$ quantile of the process. In the figure, the dots and the bars correspond to $\rho_{AB}(h)$ and $\hat{\rho}_{AB,m}(h)$ for observed distances in the sample. The dashed line corresponds to 0.03 (= 1 – 0.97) and two horizontal lines are 95% $random\ permutation$ confidence bands to check the existence of extremal dependence (see Davis et al. (2012)). The bands suggest $\rho_{(1,\infty)(1,\infty),m}(h) = 0$ for h > 2, which is consistent with (3.3).

Now consider $w(s) = \phi^{||s||}$ where $0 < \phi < 1$. Then the process (3.2) becomes

$$X_{t} = \max_{s \in \mathbb{Z}^{2}} \phi^{||s||} Z_{t-s} \quad \text{for} \quad \sum_{l \in \mathbb{Z}^{2}} \phi^{||l||} = \sum_{0 \le ||l|| < \infty} \phi^{||l||} p(||l||) < \infty, \tag{3.4}$$

where $p(||l||) = \#\{s \in \mathbb{Z}^2 : d(\mathbf{0}, s) = ||l||\}$. Observe that the process (3.4) is isotropic and that p(||l||) = O(||l||) from Lemma A.1 in Jenish and Prucha (2009), and

$$P(X_{t} \leq x) = \exp \left\{ -\frac{1}{x} \sum_{0 \leq ||l|| < \infty} \phi^{||l||} p(||l||) \right\},$$

$$P(X_{0} \leq x, X_{h} \leq x) = \exp \left\{ -\frac{1}{x} \sum_{s \in \mathbb{Z}^{2}} \max \left(\phi^{||s||}, \phi^{||h+s||} \right) \right\}$$

$$= \exp \left\{ -\frac{1}{x} \sum_{0 \leq ||l|| < \infty} \phi^{||l||} q(||l||) \right\},$$
(3.5)

where $q(||l||) = \#\{s \in \mathbb{Z}^2 : \min(||s||, ||h+s||) = ||l||\}$, the number of observations with minimum distance to $\mathbf{0}$ or h equals ||l||. For a given h, if $||l|| < \frac{||h||}{2}$, there are p(||l||) pairs from both $\mathbf{0}$ and h while $q(||l||)/p(||l||) \to 1$ as $||l|| \to \infty$. In other words,

¹http://cran.r-project.org/web/packages/SpatialExtremes/SpatialExtremes.pdf

$$q(||l||) = 2p(||l||)$$
 for $||l|| < \frac{||h||}{2}$ and $\lim_{||l|| \to \infty} \frac{q(||l||)}{p(||l||)} = 1$.

Using the joint distribution in (3.6) and a Taylor series expansion, the extremogram with $A = B = (1, \infty)$ is

$$\rho_{(1,\infty)(1,\infty)}(h) = \frac{\sum_{\substack{||h||\\2} \le ||l|| < \infty} \phi^{||l||} [2p(||l||) - q(||l||)]}{\sum_{0 \le ||l|| < \infty} \phi^{||l||} p(||l||)}.$$
(3.7)

Example 3.2. For the process (3.4), the conditions (2.1)-(2.4) in Theorem 2.1 are satisfied if $r_n^2 = o(m_n)$, $\log m_n = o(r_n)$ and $\log n = o(m_n)$.

Proof. Observe that (3.4) is isotropic. By Lemma A.1 in Jenish and Prucha (2009), p(||l||) = O(||l||). Thus, (3.1) implies that

$$\alpha_{c,c}(k) \le const \int_{\frac{k}{2}}^{\infty} j\phi^j dj = O(k\phi^k)$$
 for any $k > 0$.

Then (2.2) is satisfied if $\log m_n = o(r_n)$ since

$$m_n \sum_{l \in \mathbb{Z}^2, r_n < ||l||} \alpha_{c,c}(||l||) = m_n \sum_{r_n < ||l||} p(||l||) \alpha_{c,c}(||l||) = O\left(m_n \ r_n^2 \phi^{r_n}\right).$$

Similarly, (2.3) can be shown. If $\log n = o(m_n)$, (2.4) holds since (3.1) implies

$$n^{d/2} m_n^{1/2} \alpha_{c,cn^d}(m_n) \le const \ n^{3d/2} m_n^{1/2} m_n \phi^{m_n}$$
.

Turning to (2.1), notice from (3.5) and (3.6) that

$$\begin{split} P\left(\max_{\mathbf{s}\in B_{\gamma}}|X_{\mathbf{s}}|>\epsilon a_{m},\max_{\mathbf{s}'\in B_{\gamma}+l}|X_{\mathbf{s}'}|>\epsilon a_{m}\right) &\leq \sum_{s\in B_{\gamma}}\sum_{s'\in B_{\gamma}+l}P\left(X_{\mathbf{s}}>\epsilon a_{m},X_{\mathbf{s}'}>\epsilon a_{m}\right) \\ &\leq \sum_{s\in B_{\gamma}}\sum_{s'\in B_{\gamma}+l}\left[\frac{const}{\epsilon a_{m}}\sum_{\frac{d(s,s')}{2}\leq j<\infty}\phi^{j}j+O\left(\frac{1}{a_{m}^{2}}\right)\right] \\ &\leq const\;\frac{\phi^{||l||}||l||}{\epsilon a_{m}}+O\left(\frac{1}{a_{m}^{2}}\right). \end{split}$$

Hence the term in (2.1) is bounded by

$$\limsup_{n \to \infty} m_n \sum_{l \in \mathbb{Z}^2, k < ||l|| \le r_n} \left[const \ \frac{\phi^{||l||}||l||}{\epsilon a_m} + O\left(\frac{1}{a_m^2}\right) \right] = \sum_{k < ||l|| < \infty} const \ \phi^{||l||}||l||^2 + \limsup_{n \to \infty} O\left(\frac{m_n r_n^2}{a_m^2}\right),$$

where the second term is 0 since $a_m = O(m_n)$ and $r_n^2 = o(m_n)$. Now letting $k \to \infty$, we obtain (2.1).

Figure 1 (right) shows $\rho_{AB}(h)$ and $\hat{\rho}_{AB,m}(h)$ from a realization of the process (3.4) with $\phi = 0.5$. Here, $A = B = (1, \infty)$ and $a_m = (.90, .92, .95, .97)$ quantiles. The dots are $\rho_{AB}(h)$ and the dashed lines are $\hat{\rho}_{AB,m}(h)$ with different a_m . The ESE with $a_m = .90$ and .92 are close to the extremogram for all observed distances while the ESE with $a_m = .95$ and .97 quantiles decay faster for the observed distances greater than 3. The two horizontal lines are 95% confidence bands based on random permutations.

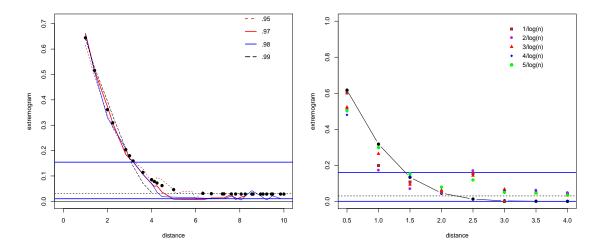


Figure 2: $\hat{\rho}_{(1,\infty)(1,\infty),m}(h)$ from a realization of Brown-Resnick process on lattice (left) and non-lattice (right). For lattice case, the ESE with $a_m = (.95, .97, .98, .99)$ upper quantiles are presented. For non-lattice case, the ESE with different bandwidths, $\frac{c}{\log n}$ with c = 1, 2, 3, 4, and 5, are displayed. Two horizontal lines are 95% random permutation confidence bands.

3.2 Brown-Resnick process

We begin with the definition of the Brown-Resnick process with Fréchet marginals. Details can be found in Kabluchko et al. (2009) or Davis et al. (2013). Consider a stationary Gaussian process $\{Z_s, s \in \mathbb{R}^d\}$ with mean 0 and variance 1 and use $\{Z_s^j, s \in \mathbb{R}^d\}$, $j \in 1, ..., n$, to denote independent replications of $\{Z_s, s \in \mathbb{R}^d\}$. For the correlation function $\rho(h) = E[Z_s Z_{s+h}]$, assume that there exist sequences $d_n \to 0$ such that

$$\log(n)\{1-\rho(d_nh)\}\to \delta(h)>0$$
, as $n\to\infty$.

Then, the random fields defined by

$$X_s(n) = \frac{1}{n} \bigvee_{i=1}^n -\frac{1}{\log\left(\Phi(Z_s^i)\right)}, \qquad s \in \mathbb{R}^d, n \in \mathbb{N},$$
(3.8)

converge weakly in the space of continuous function to the stationary Brown-Resnick process

$$X_{s} = \sup_{j \ge 1} \Gamma_{j}^{-1} Y_{s}^{j} = \sup_{j \ge 1} \Gamma_{j}^{-1} \exp\{W_{s}^{j} - \delta(s)\}, \qquad s \in \mathbb{R}^{d},$$
(3.9)

where $(\Gamma_i)_{i\geq 1}$ is an increasing enumeration of a unit rate Poisson process, $\{Y_s^j, s \in \mathbb{R}^d\}, j \in \mathbb{N}$, are iid sequences of random fields independent of $(\Gamma_i)_{i\geq 1}$, and $\{W_s^j, s \in \mathbb{R}^d\}, j \in \mathbb{N}$, are independent replications of a Gaussian random field with stationary increments, $W_0 = 0$ and $E[W_s] = 0$ and covariance function by $\operatorname{cov}(W_{s_1}, W_{s_2}) = \delta(s_1) + \delta(s_2) - \delta(s_1 - s_2)$. Here, Φ is the cumulative distribution function of N(0, 1).

The extremogram for the Brown-Resnick process $\{X_s, s \in \mathbb{R}^d\}$ with $A = (c_A, \infty)$ and $B = (c_B, \infty)$ is

$$\rho_{AB}(h) = \bar{\Phi}_{c_A,c_B} \left(\delta(h) \right) + \frac{c_A}{c_B} \bar{\Phi}_{c_B,c_A} \left(\delta(h) \right), \tag{3.10}$$

where $\Phi_{y_1,y_2}\left(\delta(h)\right) = \Phi\left(\frac{\log(y_2/y_1)}{2\sqrt{\delta(h)}} + \sqrt{\delta(h)}\right)$. To see (3.10), recall from Hüsler and Reiss (1989) that

$$F(y_1, y_2) := P(X_0 \le y_1, X_h \le y_2) = \exp\left\{-\frac{1}{y_1} \Phi\left(\frac{\log(y_2/y_1)}{2\sqrt{\delta(h)}} + \sqrt{\delta(h)}\right) - \frac{1}{y_2} \Phi\left(\frac{\log(y_1/y_2)}{2\sqrt{\delta(h)}} + \sqrt{\delta(h)}\right)\right\}.$$

As $a_m = O(m_n)$, we assume without loss of generality that $\lim_{n\to\infty} \frac{m_n}{a_m} = 1$. Then we have $p_m(A) = m_n \left(1 - e^{-\frac{1}{c_A a_m}}\right) = \frac{m_n}{c_A a_m} + O(\frac{m_n}{a_m^2}) \to \frac{1}{c_A} = \mu(A)$ and

$$\tau_{AB,m}(h) = m_n \left[1 - e^{-\frac{1}{c_A a_m}} - e^{-\frac{1}{c_B a_m}} + F(a_m c_A, a_m c_B) \right] \to \frac{1}{c_A} \bar{\Phi}_{c_A, c_B}(\delta(h)) + \frac{1}{c_B} \bar{\Phi}_{c_B, c_A}(\delta(h)), \quad (3.11)$$

which proves (3.10).

Similar to Lemma 2 in Davis et al. (2013), \(\alpha\)-mixing coefficient of the process is bounded by

$$\alpha_{m,n}(||h||) \le const \sup_{l \ge ||h||} \frac{1}{\sqrt{\delta(l)}} e^{-\delta(l)/2}. \tag{3.12}$$

In the following examples, the correlation function $\rho(h)$ of a Gaussian process $\{Z_s, s \in \mathbb{R}^d\}$ is assumed to have an expansion around zero as

$$\rho(h) = 1 - \theta ||h||^{\alpha} + o(||h||^{\alpha}), \quad h \in \mathbb{R}^d,$$
(3.13)

where $\alpha \in (0, 2]$ and $\theta > 0$. For this choice of correlation function, we have $\delta(h) = \theta ||h||^{\alpha}$ as mentioned in Davis et al. (2013), Remark 1.

Example 3.3. Consider the Brown-Resnick process $\{X_s, s \in \mathbb{Z}^d\}$ with $\delta(h) = \theta ||h||^{\alpha}$ for $0 < \alpha \le 2$ and $\theta > 0$. The conditions of Theorem 2.1 hold if $\log n = o(m_n^{\alpha})$, $\log m_n = o(r_n^{\alpha})$ and $r_n^d/m_n \to 0$. In this case, (2.6) is not satisfied for d > 0.

Proof. From (3.12), we have $\alpha_{c,c}(||h||) \leq const ||h||^{-\alpha/2} e^{-\theta||h||^{\alpha}/2}$. If $\log m_n = o(r_n^{\alpha})$, (2.2) holds since

$$m_n \sum_{l \in \mathbb{Z}^d, r_n \le ||l||} \alpha_{c,c}(||l||) \le const \ m_n \sum_{r_n \le ||l|| < \infty} ||l||^{d-1} \alpha_{c,c}(||l||) \le const \ m_n \sum_{r_n \le ||l|| < \infty} ||l||^{d-1-\alpha/2} e^{-\theta||l||^{\alpha}/2} \to 0.$$

Similarly, (2.3) can be checked. For (2.4), Proposition 3.1 implies that

$$n^{d/2} m_n^{1/2} \alpha_{c,cn^d}(m_n) \le const \ n^{3d/2} m_n^{(1-\alpha)/2} \exp\{-\theta m_n^{\alpha}/2\}$$

which converges to 0 if $\log n = o(m_n^n)$. Showing (2.1) is similar to Example 3.2. From (3.11),

$$P\left(\max_{\mathbf{s}\in B_{\gamma}}|X_{\mathbf{s}}|>\epsilon a_{m},\max_{\mathbf{s}'\in B_{\gamma}+l}|X_{\mathbf{s}'}|>\epsilon a_{m}\right) \leq const \frac{\bar{\Phi}_{(1,\infty),(1,\infty)}(\sqrt{\delta(||l||)})}{\epsilon a_{m}}+O\left(\frac{1}{a_{m}^{2}}\right).$$

Hence the term in (2.1) is bounded by

$$\limsup_{n \to \infty} \sum_{l \in \mathbb{Z}^d, k < ||l|| \le r_n} \left[const \ m_n \frac{\bar{\Phi}_{(1,\infty),(1,\infty)}(\sqrt{\delta(||l||)})}{\epsilon a_m} + O\left(\frac{1}{a_m^2}\right) \right] \\
\le const \sum_{k < ||l|| < \infty} ||l||^{d-1} e^{-\frac{\theta||l||^{\alpha}}{2}} + \limsup_{n \to \infty} O\left(\frac{r_n^d m_n}{a_m^2}\right),$$

where the second term is 0 since $r_n^d = o(m_n)$. Letting $k \to \infty$, (2.1) is obtained.

For the last statement in Example 3.3, to show (2.6) not hold, note that $|\rho_{AB,m}(h) - \rho_{AB}(h)| = O(1/m_n)$ from (3.11) and a Taylor series expansion and that $a_m = O(m_n)$ and $\log n = o(m_n^{\alpha})$.

In Figure 2 (left), we have $\rho_{AB,m}(h)$ and $\hat{\rho}_{AB,m}(h)$ from a realization of the Brown-Resnick process with $\delta(h) = \frac{2}{9}||h||^2$. We use 1600 points ($\{1,...,40\}^2 \in \mathbb{Z}^2$) to compute the extremogram with $A = B = (1,\infty)$ and $a_m = (.95, .97, .98, .99)$ upper quantiles. The extremogram is marked by dots and the ESE with different line types corresponding to various choices of a_m . From the figure, the ESE is not overly sensitive to different a_m , but $\hat{\rho}_{(1,\infty)(1,\infty),m}(h)$ with $a_m = 0.97$ quantile looks most robust. Also the extremal dependence seems to disappear for h > 4 based on the random permutation bands (two horizontal lines).

Example 3.4. Consider the Brown-Resnick process $\{X_s, s \in \mathbb{R}^2\}$ with $\delta(h) = \theta ||h||^{\alpha}$ for $\alpha \in (0, 2]$ and $\theta > 0$. Assume that $\log m_n = o(r_n^{\alpha})$ and

$$\sup_{n} \frac{\lambda_n^2 n^{2a}}{m_n} < \infty \quad \text{and} \quad \sup_{n} \frac{m_n}{\lambda_n^2 n^{2a}} < \infty \quad \text{for} \quad 0 < a < 1.$$
 (3.14)

Then Theorem 2.3 applies. Furthermore, (2.19) holds if $\frac{|S_n|\lambda_n^2}{m_n^3} \to 0$. See Appendix C for the proof.

Remark 3. Using a similar change of variable technique, as in the proof of Proposition 5.5, one can verify that condition (3.14) implies (2.14) with $\delta = 1$. We omit the details. One of the choices that satisfies condition (3.14) and $\frac{|S_n|\lambda_n^2}{m_n^3} \to 0$ is $a = \frac{7}{12}$, $\lambda_n = n^{-1/3}$ and $m_n = n^{1/2}$.

To simulate the Brown-Resnick process in \mathbb{R}^2 , we use RPbrownresnick in the RandomFields package 2 in R. Here, we consider $\delta(h)=0.5||h||^2$. In each simulation, first we generate 1600 random locations in $\{1,...,40\}^2$, where the process is simulated with the scale of $(1/\log(1600))^{1/a}$ and $\rho(\cdot)=(1+c||\cdot||^a)^{-1}$ with c=1 and a=2. For the ESE computation, we use $A=B=(1,\infty)$, $a_m=.97$ upper quantile. We set $w(\cdot)=I_{\left[-\frac{1}{2},\frac{1}{2}\right]}(\cdot)$, and distances h=(0.5,1,...,4.5,5). In Figure 2 (right), the extremogram and ESE from one realization are displayed. The extremogram $\rho_{AB}(h)$ corresponds to connected solid circles and $\hat{\rho}_{AB,m}(h)$ for different bandwidths λ_n are displayed in different point types. As will be seen in Section 3.3, smaller variances and larger biases are observed for a larger bandwidth. The two horizontal lines are the random permutation bands.

3.3 Simulation study

We use a simulation experiment to examine performances of the ESE. Samples are generated from models with Fréchet marginals for both lattice and non-lattice cases. For lattice cases, we consider MMA(1) and the Brown-Resnick process with $\delta(h) = 0.5||h||^2$. In each simulation, $\hat{\rho}_{AB,m}(h)$ with $A = B = (1, \infty)$ and $a_m = .97$ upper quantile is calculated for observed distances less than 10. This is repeated 1000 times.

Figure 3 (upper left) shows the distributions of $\hat{\rho}_{AB,m}(h)$ (box plots), $\rho_{AB}(h)$ (solid squares) and $\rho_{AB,m}(h)$ (solid circles) for MMA(1). In the figure, we see the distributions are centered at $\rho_{AB,m}(h)$, not $\rho_{AB}(h)$. Notice that $\rho_{AB,m}(h)$ for MMA(1) is computed by

²http://cran.r-project.org/web/packages/RandomFields/RandomFields.pdf

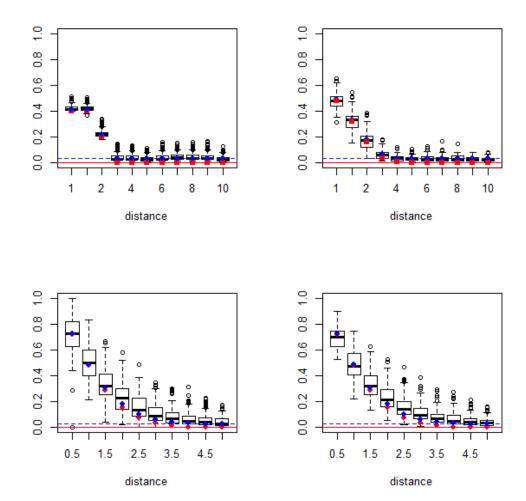


Figure 3: The distribution of the ESE for MMA(1) on lattice (upper left, 1000 simulations); the Brown-Resnick process on lattice (upper right, 1000 simulations); on \mathbb{R}^2 with $\lambda_n = 1/\log n$ (bottom left, 100 simulations); and $\lambda_n = 5/\log n$ (bottom right, 100 simulations). The solid squares are the extremogram. For MMA(1), we see the ESE is centered around PA extremogram (solid circles). For the Brown-Resnick process on \mathbb{R}^2 , we see the impact of bandwidths on the ESE.

$$P(X_h > a_m | X_0 > a_m) = \frac{1 - 2P(X_0 \le a_m) + P(X_h \le a_m, X_0 \le a_m)}{P(X_0 > a_m)}$$

$$= \frac{\frac{\frac{2}{m} - 1 + (1 - \frac{1}{m})^{8/5}}{1/m} \quad \text{for } ||h|| = 1, \sqrt{2},$$

$$= \frac{\frac{2}{m} - 1 + (1 - \frac{1}{m})^{9/5}}{1/m} \quad \text{for } ||h|| = 2,$$

$$\frac{1}{m} \quad \text{for } ||h|| > 2.$$

using $P(X > a_m) = \frac{1}{m}$ and $P(X \le x) = e^{-5/x}$ for x > 0, and $m = 0.03^{-1}$.

The upper right panel of the figure presents the distributions of the ESE with $\rho_{AB}(h)$ (solid squares) and $\rho_{AB,m}(h)$ (solid circles) for the Brown-Resnick process on the lattice. The derivation of $\hat{\rho}_{AB,m}(h)$ is from



Figure 4: The region of Florida rainfall data.

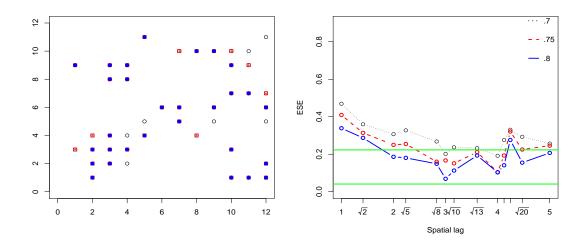


Figure 5: The locations of extremes (left) and the ESE (right) using the 6 year maxima of Florida rainfall data. For example, the ESE with 0.80 upper quantile (solid line, right) is based on the locations of corresponding extremes (solid circles, left). The ESE using the 0.70 upper quantile indicates that no spatial extremal dependence for lags larger than 3.

(3.11). Again, the ESE is centered around PA extremogram.

The bottom panels of Figure 3 are based on the simulation results from the Brown-Resnick process in the non-lattice case. For each simulation, 1600 points are generated from a Poisson process in $\{1, ..., 40\}^2$, from which $\hat{\rho}_{AB,m}(h)$ for h = (0.5, 1, ..., 4.5, 5) is computed using the bandwidths $\lambda_n = 1/\log n$ and $5/\log n$. This is repeated 100 times. Notice that the ESE using $\lambda_n = 1/\log n$ has generally smaller bias but larger variance compared to the ESE using $\lambda_n = 5/\log n$ for $h \leq 2$. For longer lags, the differences is not apparent. This indicates that the ESE with wider bandwidths tends to have smaller variance but larger biases.

4 Application

In this section, we apply the ESE to analyze geographical dependence of heavy rainfall in a region in Florida. The source is Southwest Florida Water Management District. The raw data is total rainfall in 15 minute intervals from 1999 to 2004, measured on a $120 \times 120 \text{ (km)}^2$ region containing 3600 grid locations. The

region of the measurements is shown in Figure 4. For each fixed time, we first calculate the spatial maximum over a non-overlapping block of size $10 \times 10 \text{ (km)}^2$, which provides a $12 \times 12 \text{ grid}$ of spatial maxima. Then, we calculate the annual maxima from 1999 to 2004 and the 6 year maxima from the corresponding time series for each spatial maximum. The 7 spatial data sets on a $12 \times 12 \text{ grid}$ under consideration consist of annual maxima and 6 year maxima of spatial maxima. Since the data are constructed as a maxima over a spatial grid of 25 locations and a temporal resolution of 15 minutes intervals, it is not unreasonable to view these 7 spatial data sets as realizations from a max-stable process.

We first look at the spatial extremal dependence for 6 year maxima rainfall. In Figure 5, the locations of extremes (left) and the ESE (right) are displayed, where the ESE is computed using $A = B = (1, \infty)$ and $a_m = .70$ (dotted line), .75 (dashed line) and .80 (solid line) upper quantiles. Since the number of spatial locations is small (144), we chose modest thresholds in order to ensure enough exceedances for estimation of the ESE. Such thresholds should provide good estimates of the pre-asymptotic extremogram for a max-stable process. The locations of extremes are marked corresponding to choices of a_m by .70 (empty circles), .75 (empty squares) and .80 (solid circles) upper quantiles. For the ESE plot, the horizontal lines are permutation based confidence bands. For example, if extreme events are defined by any rainfall heavier than the .70 upper quantile of the maxima rainfall observed for the entire periods, there is a significant extremal dependence between two clusters at distance 2. On the other hand, using the 0.80 upper quantile, the extremal dependence at the same distance is no longer significant. In the case of 6 year maxima rainfall, the ESE from the 0.70 upper quantile indicates that no spatial extremal dependence for spatial lags larger than 3. A small spike of the ESE at spatial lags around 4 may be the result of two extremal clusters that are 4 units apart, as seen in the left panel of Figure 5.

By looking at the ESE of annual maxima rainfall from 1999 to 2004, we see year-over-year changes in spatial extremal dependence. Figure 6 presents the locations of extremes and the ESE from 1999 to 2004 (left to right, top to bottom). For example, the ESE suggests that the spatial extremal dependence for lags less than 3 in 2000 is stronger than at any other year between 1999 and 2004. Using the .80 upper quantile, there is significant extremal dependence for spatial lag $\sqrt{8}$ in 2000, but not for any other years. In 2002, the spatial extremal dependence is not significant at lag $\sqrt{8}$ using the .80 upper quantile. Similarly, the year-to-year comparisons of the ESE with 0.70 and 0.75 upper quantiles confirm that the spatial extremal dependence for spatial lags up to 3 is stronger in 2000 than in any other years.

5 Appendix: Proofs

The following proposition presented by Li et al. (2008) is used in the proof. The proposition is analogous to Theorem 17.2.1 in Ibragimov and Linnik (1971).

Proposition 5.1 (Lemma A.1. in Li et al. (2008)). Let U and V be two closed and connected sets in \mathbb{R}^d such that $\#U = \#V \leq b$ and $d(U,V) \geq r$ for some constants b and r. For a stationary process X_s , consider ξ and η measurable random variables with respect to $\sigma(X_s : s \in U)$ and $\sigma(X_s : s \in V)$ with $|\xi| \leq C_1$, $|\eta| \leq C_2$. Then $|\cos(\xi, \eta)| \leq 4C_1C_2\alpha_{b,b}(r)$.

5.1 Appendix A: Proof of Theorem 2.1

Theorem 2.1 is derived from Theorem 5.2. For notation, we suppress the dependence of m on n and write m for m_n . Define a vector valued random field by

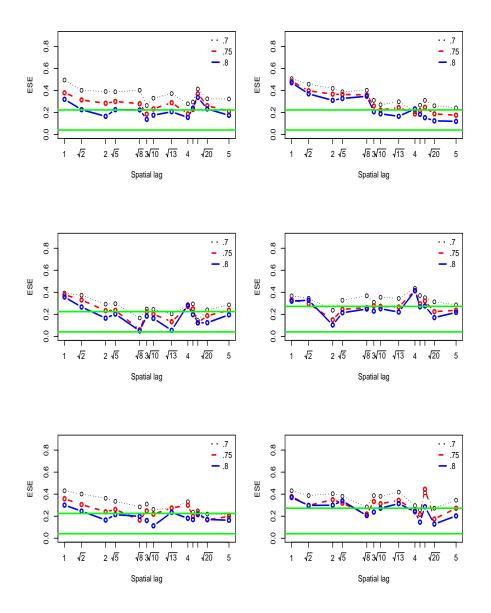


Figure 6: The ESE of the annual maxima of Florida rainfall from 1999 to 2004 (left to right, top to bottom). The ESE indicates that the spatial extremal dependency for spatial lags less than 3 is the strongest in 2000.

$$Y_t = X_{D_t}$$
, where $D_t = t + B_{\gamma} = \{s \in \mathbb{Z}^d : d(t, s) \leq \gamma\}$.

In Theorem 5.2, we will establish a joint central limit theorem for

$$\hat{P}_m(C) = \frac{m_n}{n^d} \sum_{t \in \Lambda_n} I_{\{Y_t/a_m \in C\}} = \frac{m_n}{n^d} \sum_{t \in \Lambda_n^p} I_{\{Y_t/a_m \in C\}} + \frac{m_n}{n^d} \sum_{t \in \Lambda_n \setminus \Lambda_n^p} I_{\{Y_t/a_m \in C\}}, \tag{5.1}$$

where $\Lambda_n^p = \{t \in \Lambda_n : d(t, \partial \Lambda_n) \ge p\}$ and $\partial \cdot$ denotes the boundary. In fact, showing a CLT for the first term in (5.1) is sufficient as the second term is negligible as $n \to \infty$. Recall that

$$p_m(A) = mP(X_0 \in a_m A) \text{ and } \tau_{AB,m}(h) = mP(X_0 \in a_m A, X_h \in a_m B),$$

where A and B are sets bounded away from the origin. Write $\mu(A) = \lim_{n \to \infty} p_m(A)$,

$$\begin{split} \tau_{AB}(h) &= \lim_{n \to \infty} \tau_{AB,m}(h), \\ \mu_A(D_\mathbf{0}) &= \lim_{x \to \infty} P\left(\frac{Y_t}{\|Y_t\|} \in A \middle| \|Y_t\| > x\right), \\ \tau_{A \times B}(D_\mathbf{0} \times D_l) &= \lim_{x \to \infty} P\left(\frac{(Y_\mathbf{0}, Y_l)}{\|\text{vector}\{Y_\mathbf{0}, Y_l\}\|} \in A \times B \middle| \|\text{vector}\{Y_\mathbf{0}, Y_l\}\| > x\right). \end{split}$$

Theorem 5.2. Assume the conditions of Theorem 2.1. Let C be a set bounded away from zero and a continuity set with respect to μ and τ . Then

$$S_n = \left(\frac{m_n}{n^d}\right)^{1/2} \sum_{s \in \Lambda_n} \left[I\left(\frac{Y_s}{a_m} \in C\right) - P\left(\frac{Y_s}{a_m} \in C\right) \right] \xrightarrow{d} N(0, \sigma_Y^2(C))$$

where $\sigma_Y^2(C) = \mu_C(D_0) + \sum_{l \neq 0 \in \mathbb{Z}^d} \tau_{C \times C}(D_0 \times D_l).$

Proof. We use ideas from Bolthausen (1982) and Davis and Mikosch (2009) to show the CLT for quantity in (5.1)

$$\hat{P}_m(C) = m_n \sum_{s \in \Lambda_n} I_s / |\Lambda_n|$$
 where $I_s = I_{\{X_s/a_m \in C\}}$.

The proof for the CLT of X_s replaced by a vector valued random field Y_s in indicator is analogous.

Define $H[a,b]=\{d(s,t):a\leq d(s,t)\leq b\}$ and $||l||=d(\mathbf{0},l)$ for convenience. Assume $m_n^{2+2d}=o(n^d)$, $r_n^d = o(m_n)$, and

$$\lim_{k \to \infty} \limsup_{n \to \infty} m_n \sum_{l \in \mathbb{Z}^d, ||l|| \in H(k, r_n]} P(|X_l| > \varepsilon a_m, |X_0| > \varepsilon a_m) = 0 \quad \text{for} \quad \forall \epsilon > 0,$$
 (5.2)

$$\lim_{n \to \infty} m_n \sum_{l \in \mathbb{Z}^d, ||l|| \in H(r_n, \infty)} \alpha_{1,1}(||l||) = 0,$$

$$\sum_{l \in \mathbb{Z}^d} \alpha_{j_1, j_2}(||l||) < \infty \quad \text{for} \quad 2 \le j_1 + j_2 \le 4,$$
(5.4)

$$\sum_{l \in \mathbb{Z}^d} \alpha_{j_1, j_2}(||l||) < \infty \quad \text{for} \quad 2 \le j_1 + j_2 \le 4, \tag{5.4}$$

$$\lim_{n \to \infty} n^{d/2} m_n^{1/2} \alpha_{1,n^d}(m_n) = 0, \tag{5.5}$$

which are univariate case analog of conditions (2.1) - (2.4).

By the same arguments in Davis and Mikosch (2009),

$$E\hat{P}_m(C) \to \mu(C)$$
 (5.6)

$$\operatorname{var}\left(\hat{P}_{m}(C)\right) \sim \frac{m_{n}}{n^{d}} \left| \mu(C) + \sum_{l \neq \mathbf{0} \in \mathbb{Z}^{d}} \tau_{CC}(l) \right| = \frac{m_{n}}{n^{d}} \sigma_{X}^{2}(C), \tag{5.7}$$

where (5.6) is implied by the regularly varying assumption. To see (5.7), observe that

$$\frac{n^d}{m_n} \operatorname{var}\left(\hat{P}_m(C)\right) = \frac{m_n}{n^d} \sum_{s \in \Lambda_n} \operatorname{var}(I_s) + \frac{m_n}{n^d} \sum_{s, t \in \Lambda_n, s \neq t} \operatorname{cov}(I_s, I_t) = A_1 + A_2.$$
 (5.8)

By the regularly varying assumption, $A_1 = p_m(C) + (p_m(C))^2/m_n \to \mu(C)$. Turning to A_2 , for $k \ge 1$ fixed,

$$\begin{split} A_2 &\sim \frac{m_n}{n^d} \sum_{l=(l_1,\ldots l_d) \neq \mathbf{0}, ||l|| \leq \max \Lambda_n} \Pi_{i=1}^d (n-|l_i|) \mathrm{cov}(I_{\mathbf{0}}, I_l) \\ &= \frac{m_n}{n^d} \left[\sum_{l \in \mathbb{Z}^d, ||l|| \in H(0,k]} \cdot + \sum_{l \in \mathbb{Z}^d, ||l|| \in H(k,r_n]} \cdot + \sum_{l \in \mathbb{Z}^d, ||l|| \in H(r_n, \max \Lambda_n]} \cdot \right] = A_{21} + A_{22} + A_{23} \end{split}$$

where $\max \Lambda_n = \{\max(d(s,t)) : s,t \in \Lambda_n\}$ and $\Pi_{i=1}^d(n-|l_i|)$ counts a number of cubes with lag l in Λ_n . From the regularly varying assumption, $\lim_{k\to\infty} \limsup_{n\to\infty} A_{21} = \sum_{l\neq \mathbf{0}\in\mathbb{Z}^d} \tau_{CC}(l)$ since

$$\limsup_{n \to \infty} A_{21} = \sum_{l \in \mathbb{Z}^d, ||l|| \in H(0,k]} \limsup_{n \to \infty} \left(\tau_{CC,m}(C) - p_m(C) \frac{p_m(C)}{m_n} \right) = \sum_{l \in \mathbb{Z}^d, ||l|| \in H(0,k]} \tau_{CC}(l).$$

Thus, it is sufficient to show

$$\lim_{k \to \infty} \limsup_{n \to \infty} (|A_{22}| + |A_{23}|) = 0$$

to achieve (5.7). Recall that C is bounded away from the origin. Notice that

$$A_{22} \leq const \ m_n \sum_{l \in \mathbb{Z}^d, ||l|| \in H(k, r_n]} \left[P(|X_l| > \varepsilon a_m, |X_{\mathbf{0}}| > \varepsilon a_m) + \left(\frac{p_m(C)}{m_n}\right)^2 \right],$$

$$A_{23} \leq const \ m_n \sum_{l \in \mathbb{Z}^d, ||l|| \in H(r_n, \infty)} \alpha_{1,1}(||l||),$$

so (5.7) holds assuming (5.2), (5.3) and $r_n^d/m_n \to 0$.

Now, we prove

$$\sqrt{\frac{n^d}{m_n}}(\hat{P}_m(C) - p_m(C)) = \sqrt{\frac{m_n}{n^d}} \sum_{s \in \Lambda_n} \bar{I}_s \xrightarrow{d} N(0, \sigma_X^2(C))$$

$$(5.9)$$

where $\bar{I}_t = I\left(\frac{X_t}{a_m} \in C\right) - P\left(\frac{X}{a_m} \in C\right)$. First, infer from (5.8) that

$$\frac{m_n}{n^d} \sum_{s,t \in \Lambda_n} |\operatorname{cov}(\bar{I}_s, \bar{I}_t)| < \infty. \tag{5.10}$$

As the next step, define

$$S_{\alpha,n} = \sum_{\beta \in \Lambda_n, d(\alpha,\beta) < m_n} \sqrt{\frac{m_n}{n^d}} \bar{I}_{\beta}, \ v_n = \sum_{\alpha \in \Lambda_n} E(\sqrt{\frac{m_n}{n^d}} \bar{I}_{\alpha} S_{\alpha,n}), \ \bar{S}_n = v_n^{-1/2} S_n, \ \text{and} \ \bar{S}_{\alpha,n} = v_n^{-1/2} S_{\alpha,n}.$$

From the definition, $v_n \sim \text{var}(S_n) \to \sigma^2(C)$.

Now, use Stein's lemma to show (5.9) as in Bolthausen (1982) by checking $\lim_{n\to\infty} E((i\lambda - \bar{S}_n)e^{i\lambda\bar{S}_n}) = 0$

for all $\lambda \in R$. Write

$$(i\lambda - \bar{S}_n)e^{i\lambda\bar{S}_n} = i\lambda e^{i\lambda\bar{S}_n} (1 - v_n^{-1} \sum_{\alpha \in \Lambda_n} \sqrt{\frac{m_n}{n^d}} \bar{I}_{\alpha} S_{\alpha,n}) - v_n^{-1/2} e^{i\lambda\bar{S}_n} \sum_{\alpha \in \Lambda_n} \sqrt{\frac{m_n}{n^d}} \bar{I}_{\alpha} [1 - e^{-i\lambda\bar{S}_{\alpha,n}} - i\lambda\bar{S}_{\alpha,n}]$$
$$-v_n^{-1/2} \sum_{\alpha \in \Lambda_n} \sqrt{\frac{m_n}{n^d}} \bar{I}_{\alpha} e^{-i\lambda(\bar{S}_n - \bar{S}_{\alpha,n})}$$
$$= B_1 + B_2 + B_3.$$

We will show $E|B_1|^2 \to 0$. From Proposition 5.1, when $d(\alpha, \alpha') \geq 3m_n$,

$$|\operatorname{cov}(\bar{I}_{\alpha}\bar{I}_{\beta},\bar{I}_{\alpha'}\bar{I}_{\beta'})| \leq 4 \alpha_{2,2}(d(\alpha,\alpha')-2m_n).$$

When $d(\alpha, \alpha') < 3m_n$, let $j = \min\{d(\alpha, \alpha'), d(\alpha, \beta'), d(\beta, \alpha'), d(\beta, \beta')\}$. Then

$$|\operatorname{cov}(\bar{I}_{\alpha}\bar{I}_{\beta}, \bar{I}_{\alpha'}\bar{I}_{\beta'})| \leq 4 \alpha_{p,q}(j)$$

for $2 \le p + q \le 4$. Given $m_n^{2+2d} = o(n^d)$, we have $E|B_1|^2 \to 0$ since

$$E|B_{1}|^{2}$$

$$= \lambda^{2} v_{n}^{-2} \sum_{\alpha,\alpha',\beta,\beta',d(\alpha,\beta) \leq m_{n},d(\alpha',\beta') \leq m_{n}} \frac{m_{n}^{2}}{n^{2d}} \operatorname{cov}(\bar{I}_{\alpha}\bar{I}_{\beta},\bar{I}_{\alpha'}\bar{I}_{\beta'})$$

$$\leq \frac{\lambda^{2}}{v_{n}^{2}} \frac{m_{n}^{2}}{n^{2d}} \left[\sum_{\alpha \in \Lambda_{n}} \sum_{\alpha' \in \Lambda_{n} \cap \{d(\alpha,\alpha') > 3m_{n}\}} \sum_{\beta,\beta'} \left| \operatorname{cov}(\bar{I}_{\alpha}\bar{I}_{\beta},\bar{I}_{\alpha'}\bar{I}_{\beta'}) \right| + \sum_{\alpha \in \Lambda_{n}} \sum_{\alpha' \in \Lambda_{n} \cap \{d(\alpha,\alpha') \leq 3m_{n}\}} \sum_{\beta,\beta'} \left| \operatorname{cov}(\bar{I}_{\alpha}\bar{I}_{\beta},\bar{I}_{\alpha'}\bar{I}_{\beta'}) \right| \right]$$

$$\leq \frac{\lambda^{2}}{v_{n}^{2}} \frac{m_{n}^{2}}{n^{2d}} 4 \left[\sum_{\alpha \in \Lambda_{n}} \sum_{\alpha' \in \Lambda_{n} \cap \{d(\alpha,\alpha') > 3m_{n}\}} \sum_{\beta,\beta'} \alpha_{2,2}(d(\alpha,\alpha') - 2m_{n}) + \sum_{\alpha \in \Lambda_{n}} \sum_{\alpha' \in \Lambda_{n} \cap \{d(\alpha,\alpha') \leq 3m_{n}\}} \sum_{\beta,\beta'} \alpha_{p,q}(j) \right]$$

$$\leq \frac{\operatorname{const}\lambda^{2} m_{n}^{2}}{v_{n}^{2} n^{2d}} n^{d} m_{n}^{2d} \left[\sum_{l \in \mathbb{Z}^{d}, ||l|| \in H(3m_{n},\infty)} \alpha_{2,2}(||l|| - 2m_{n}) + \sum_{l \in \mathbb{Z}^{d}, ||l|| \in H[0,3m_{n}]} \alpha_{p,q}(||l||) \right]$$

$$= O(m_{n}^{2+2d}/n^{d}).$$
(5.11)

Notice that in (5.11), $n^d m_n^{2d}$ is from summing over α (giving n^d), β (giving $O(m_n^d)$), and β' (giving $O(m_n^d)$) for the first summation. Similarly, for the second summation, $n^d m_n^{2d}$ is from summing over α , β and α' or β' depending on the location of points. The last equation is from (5.4).

Now we show $E|B_2| \to 0$ provided $m_n^{2+2d} = o(n^d)$. Recall that $|e^{ix} - 1 - ix| \le \frac{1}{2}x^2$. Then

$$\begin{split} E|B_2| & \leq cv_n^{-1/2}n^d\sqrt{\frac{m_n}{n^d}}E\bar{S}_{\alpha,n}^2\\ & = cv_n^{-1/2}\sqrt{\frac{m_n}{n^d}}m_n\sum_{\beta,\beta',d(\mathbf{0},\beta)\leq m_n,d(\mathbf{0},\beta')\leq m_n}E(\bar{I}_{\beta}\bar{I}_{\beta'})\\ & \leq c\sqrt{\frac{m_n}{n^d}}m_n^{d+1}\sum_{l\in\Lambda_n}E(\bar{I}_{\mathbf{0}}\bar{I}_l)\\ & = O\left(\sqrt{\frac{m_n^{1+2d}}{n^d}}\right) \end{split}$$

where $m_n \sum_{l \in \Lambda_n} E(\bar{I}_0 \bar{I}_l) < \infty$ is inferred from (5.10).

Lastly, the condition (5.5) implies $|EB_3| \to 0$ since

$$|EB_3| \le cv_n^{-1/2}n^d\sqrt{\frac{m_n}{n^d}}\alpha_{1,n^d}(m_n) = cn^{d/2}m_n^{1/2}\alpha_{1,n^d}(m_n).$$

Thus, Stein's lemma is satisfied, which completes the proof.

Remark 4. $\hat{P}_m(C)$ is a consistent estimator of $\mu(C)$. If $\mu(C) = 0$, var $(\hat{P}_m(C)) = o(m_n/n^d)$.

Remark 5. The conditions (2.1) - (2.4) are derived from (5.2) - (5.5) by replacing univariate process (X_t) by vectorized process (Y_t) . In order to see (2.1) is derived from (5.2), for example, consider Euclidean norm for (Y_t) process. Then, the vectorized analog of (5.2) is

$$\lim_{k \to \infty} \limsup_{n \to \infty} m_n \sum_{l \in \mathbb{Z}^d, ||l|| \in H(k, r_n)} P(||Y_0|| > \epsilon a_m, ||Y_l|| > \epsilon a_m) = 0,$$

which holds under (2.1) by triangular inequality, i.e.,

$$P(||Y_{\mathbf{0}}|| > \epsilon a_m, ||Y_l|| > \epsilon a_m) \leq P(\sum_{s \in D_{\mathbf{0}}} |X_s| > \epsilon a_m, \sum_{s' \in D_l} |X_{s'}| > \epsilon a_m) \leq P\left(\max_{s \in D_{\mathbf{0}}} |X_s| > \frac{\epsilon a_m}{|D_{\mathbf{0}}|}, \max_{s' \in D_l} |X_{s'}| > \frac{\epsilon a_m}{|D_l|}\right).$$

The rest of the derivations are straightforward.

Proof of Theorem 2.1. Apply the Cramér-Wold device to Theorem 5.2 to achieve the multivariate central limit theorem, then use δ -method to obtain the central limit theorem for the ESE. To specify the limiting variance Σ , redefine

$$\mu(A) = \lim_{x \to \infty} P\left(\frac{X_t}{\|Y_t\|} \in A \middle| \|Y_t\| > x\right).$$

Then, $\Sigma = \mu(A)^{-4} F \Pi F^t$ where

$$\Pi_{i,i} = \mu_{S_i}(D_{\mathbf{0}}) + \sum_{l \neq \mathbf{0} \in \mathbb{Z}^d} \tau_{S_i \times S_i}(D_{\mathbf{0}} \times D_l)$$

$$\Pi_{i,j} = \mu_{S_i \cap S_j}(D_{\mathbf{0}}) + \sum_{l \neq \mathbf{0} \in \mathbb{Z}^d} \tau_{S_i \times S_j}(D_{\mathbf{0}} \times D_l)$$

$$F = \begin{pmatrix} \mu(S_{(\#H)+1}) & 0 & 0 & \dots & 0 & -\mu_{S_1}(D_{\mathbf{0}}) \\ 0 & \mu(S_{(\#H)+1}) & 0 & \dots & 0 & -\mu_{S_2}(D_{\mathbf{0}}) \\ \vdots & \vdots & \ddots & \dots & \vdots & \vdots \\ 0 & 0 & 0 & \dots & \mu(S_{(\#H)+1}) & -\mu_{S_{(\#H)}}(D_{\mathbf{0}}) \end{pmatrix}$$

where the sets S_i are chosen such that $\{Y_t \in S_i\} = \{X_t \in A, X_s \in B : d(t,s) = h_i\}$ for $h_i \in H$ and i = 1, ..., (#H) and $\{Y_t \in S_{(\#H)+1}\} = \{X_t \in A\}$. For more details, see Davis and Mikosch (2009).

5.2 Appendix B: Proof of Theorem 2.3

Theorem 2.3 is derived from Proposition 5.4 - 5.6. Before proceeding to Proposition 5.4, we present the following result regarding LUNC.

Proposition 5.3. Consider a strictly stationary regularly varying random field $\{X_s, s \in \mathbb{R}^d\}$ with index $\alpha > 0$ satisfying LUNC. For a positive integer k and $\lambda_n \to 0$,

$$nP\left(\frac{X_{\theta}}{a_n} \in A_0, \frac{X_{s_1+\lambda_n}}{a_n} \in A_1, \cdots, \frac{X_{s_k+\lambda_n}}{a_n} \in A_k\right) \to \tau_{A_0, A_1, \cdots, A_k}(s_1, \cdots, s_k)$$

provided $A_0 \times A_1 \times \cdots \times A_k$ is a continuity set of the limit measure

$$\tau_{A_0,A_1,\dots,A_k}(s_1,\dots,s_k) = \lim_{n\to\infty} nP(X_0/a_n \in A_0, X_{s_1}/a_n \in A_1,\dots,X_{s_k}/a_n \in A_k).$$

Proof. Let f be a continuous function with compact support on $\mathbb{R}^{k+1} \setminus \{\mathbf{0}\}$. Since f has compact support, it is uniformly continuous and hence for every $\epsilon > 0$ there exists δ such that $|f(x_1, x_2, \dots, x_{k+1}) - f(y_1, y_2, \dots, y_{k+1})| < \epsilon$ whenever $|(x_1, x_2, \dots, x_{k+1}) - (y_1, y_2, \dots, y_{k+1})| < \delta$.

Let
$$\tilde{X}_n = (X_0, X_{s_1+\lambda_n}, \cdots, X_{s_k+\lambda_n})$$
 and $\tilde{X} = (X_0, X_{s_1}, \cdots, X_{s_k})$. Notice that

$$nE\left|f\left(\frac{\tilde{X}_n}{a_n}\right) - f\left(\frac{\tilde{X}}{a_n}\right)\right| \ = \ nE|\cdot|I_{\{\frac{|\tilde{X}_n - \tilde{X}|}{a_n} > \delta\}} + nE|\cdot|I_{\{\frac{|\tilde{X}_n - \tilde{X}|}{a_n} \le \delta\}} = A_1 + A_2.$$

Let $M = \max f\left(\frac{\mathbf{X}}{a_n}\right)$. By (2.17), there exists $\epsilon > 0$ such that

$$\limsup_{n} A_{1} \leq \limsup_{n} 2Mn \left[P\left(|X_{s_{1}+\lambda_{n}} - X_{s_{1}}| > \frac{\delta a_{n}}{k} \right) + \dots + P\left(|X_{s_{k}+\lambda_{n}} - X_{s_{k}}| > \frac{\delta a_{n}}{k} \right) \right] < 2M\epsilon$$

since $|X_{\lambda_n} - X_{\mathbf{0}}| \le \sup_{|s| < \delta'} |X_s - X_{\mathbf{0}}|$ as $n \to \infty$ for $|\lambda_n| < \delta'$. For A_2 , since the support of $f \in \{|\tilde{X}| > C\} \subset \{|X_{\mathbf{0}}| > \frac{C}{k+1}\} \cup \cdots \cup \{|X_{s_k}| > \frac{C}{k+1}\}$

$$\limsup_{n} A_{2} \leq \limsup_{n} \epsilon n \left[P\left(\frac{|\tilde{X}_{n}|}{a_{n}} > C\right) + P\left(\frac{|\tilde{X}|}{a_{n}} > C\right) \right]$$

$$= \limsup_{n} \epsilon n \ 2(k+1) \ P(|X_{\mathbf{0}}| > a_{n}C/(k+1))$$

$$= \epsilon 2(k+1)\tau_{BB}(\mathbf{0}), \quad \text{where } B = \{x : x > C/(k+1)\}.$$

Take ϵ small by choosing appropriate δ and δ' , then for a positive integer k and $\lambda_n \to 0$,

$$nEf\left(\frac{X_0, X_{s_1+\lambda_n}, \cdots, X_{s_k+\lambda_n}}{a_n}\right) \to \int f(u_1, u_2, \cdots, u_k)\mu(du_1, du_2, \cdots, du_k)$$

for any continuous function with compact support f. Using Portmanteau theorem for vague convergence, we complete the proof. See Theorem 3.2 in Resnick (2006).

We discuss asymptotics of the denominator and the numerator of the ESE in turn.

Proposition 5.4. Under the setting of Theorem 2.3 and condition (M2),

$$E(\hat{p}_m(A)) = p_m(A) \to \mu(A)$$
 and $\frac{|S_n|}{m_n} var(\hat{p}_m(A)) \to \frac{\mu(A)}{\nu} + \int_{\mathbb{R}^2} \tau_{AA}(y) dy$.

Hence, $\hat{p}_m(A) \xrightarrow{p} \mu(A)$.

Proof. By the regularly varying property, $E(\hat{p}_m(A)) = p_m(A) \to \mu(A)$.

For $\operatorname{var}(\hat{p}_m(A))$, recall that $N^{(2)}(ds_1, ds_2) = N(ds_1)N(ds_2)I(s_1 \neq s_2)$ and observe that

$$\begin{split} &E(\hat{p}_{m}(A)^{2}) \\ &= \left(\frac{m_{n}}{\nu|S_{n}|}\right)^{2} E\left[\int_{S_{n}} I\left(\frac{X_{s_{1}}}{a_{m}} \in A\right) N(ds_{1}) + \int_{S_{n}} \int_{S_{n}} I\left(\frac{X_{s_{1}}}{a_{m}} \in A, \frac{X_{s_{2}}}{a_{m}} \in A\right) N^{(2)}(ds_{1}, ds_{2})\right] \\ &= \left(\frac{m_{n}}{\nu|S_{n}|}\right)^{2} \left[\int_{S_{n}} \frac{p_{m}(A)}{m_{n}} \nu ds_{1} + \int_{S_{n}} \int_{S_{n}} \left[P\left(\frac{X_{s_{1}}}{a_{m}} \in A, \frac{X_{s_{2}}}{a_{m}} \in A\right) - \frac{p_{m}(A)^{2}}{m_{n}^{2}}\right] \nu^{2} ds_{1} ds_{2}\right] + E(\hat{p}_{m}(A))^{2} \\ &= \left(\frac{m_{n}}{|S_{n}|}\right) \left[\frac{E(\hat{p}_{m}(A))}{\nu} + \int_{S_{n}-S_{n}} m_{n} \left[\frac{\tau_{AA,m}(y)}{m_{n}} - \frac{p_{m}(A)^{2}}{m_{n}^{2}}\right] \frac{|S_{n} \cap (S_{n}-y)|}{|S_{n}|} dy\right] + E(\hat{p}_{m}(A))^{2} \end{split}$$

where the change of variables $s_2 - s_1 = y$ is used in the last line. Using the above, we show

$$\frac{|S_n|}{m_n} \text{var}(\hat{p}_m(A)) = \frac{E(\hat{p}_m(A))}{\nu} + \int_{S_n - S_n} m_n \left[\frac{\tau_{AA,m}(y)}{m_n} - \frac{p_m(A)^2}{m_n^2} \right] \frac{|S_n \cap (S_n - y)|}{|S_n|} dy
\rightarrow \frac{\mu(A)}{\nu} + \int_{\mathbb{R}^2} \tau_{AA}(y) dy.$$
(5.12)

To see (5.12), notice that for a fixed k > 0

$$\int_{S_n - S_n} m_n \left[\frac{\tau_{AA,m}(y)}{m_n} - \frac{p_m(A)^2}{m_n^2} \right] \frac{|S_n \cap (S_n - y)|}{|S_n|} dy = \int_{B[0,k)} [\cdot] dy + \int_{B[k,r_n]} [\cdot] dy + \int_{(S_n - S_n) \setminus B[0,r_n]} [\cdot] dy \\
= A_1 + A_2 + A_3.$$

For each fixed k > 0, $\lim_{n \to \infty} A_1 = \int_{B[0,k)} \tau_{AA}(y) dy$. Now, we show

$$\lim_{k \to \infty} \limsup_{n \to \infty} (|A_2 + A_3|) = 0.$$

Recall that A is bounded away from the origin. Using (2.11) and $r_n^2 = o(m_n)$,

$$|A_2| \le \int_{B[k,r_n]} m_n P(|X_y| > \epsilon a_m, |X_0| > \epsilon a_m) dy + const \ r_n^2 \frac{p_m(A)^2}{m_n} \to 0$$

From (2.12),
$$\lim_{n} |A_3| \leq \lim_{n} \int_{\mathbb{R}^2 \setminus B[0,r_n)} m_n \alpha_{1,1}(y) dy = 0$$
. This completes the proof.

Proposition 5.5. Assume that a stationary regularly varying random field satisfies LUNC. Further, assume the conditions of Proposition 5.4, and (2.15) in (M3). Then

(i) $E\hat{\tau}_{AB,m}(h) \to \tau_{AB}(h)$,

$$(ii) \ \frac{|S_n|\lambda_n^2}{m_n} \cos \left(\hat{\tau}_{AB,m}(h_1), \hat{\tau}_{AB,m}(h_2)\right) \to \frac{\int_{\mathbb{R}^2} w(y)^2 dy}{\nu^2} \left[\tau_{AB}(h_1) \ I_{\{h_1=h_2\}} + \tau_{A\cap BA\cap B}(h_1) \ I_{\{h_1=-h_2\}}\right], \ and \ u_{AB}(h_1) = 0$$

$$(iii) \ \frac{|S_n|\lambda_n^2}{m_n} \ \mathrm{var} \left(\hat{\tau}_{AB,m}(h) \right) \to \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{AB}(h).$$

Proof. (i) From (2.10) and stationarity of $\{X_s, s \in \mathbb{R}^2\}$

$$E\hat{\tau}_{AB,m}(h) = \frac{m_n}{\nu^2} \frac{1}{|S_n|} \int_{S_n} \int_{S_n} w_n(h + s_1 - s_2) P\left(\frac{X_0}{a_m} \in A, \frac{X_{s_2 - s_1}}{a_m} \in B\right) \nu^2 ds_1 ds_2$$

which after making the transformation $\frac{h+s_1-s_2}{\lambda_n}=y$ and $s_2=u$ becomes

$$\frac{1}{|S_n|} \int_{\frac{S_n - S_n + h}{\lambda_n}} \int_{S_n \cap (S_n - \lambda_n y + h)} w(y) \, \tau_{AB,m}(h - y\lambda_n) du dy$$

$$= \int_{\frac{S_n - S_n + h}{\lambda_n}} w(y) \, \tau_{AB,m}(h - y\lambda_n) \frac{|S_n \cap (S_n - \lambda_n y + h)|}{|S_n|} dy$$

$$\to \tau_{AB}(h).$$

The limit in the last line follows from LUNC and the dominated convergence theorem since

$$\tau_{AB,m}(h-y\lambda_n)\frac{|S_n\cap(S_n-\lambda_ny+h)|}{|S_n|}\leq p_m(A) \text{ and } \int_{\mathbb{R}^2}w(y)\ p_m(A)dy<\infty.$$

(ii) For fixed sets A and B let $\tau_m^*(s_1, s_2, s_3, s_4) = m_n P\left(\frac{X_{s_1}}{a_m} \in A, \frac{X_{s_2}}{a_m} \in B, \frac{X_{s_3}}{a_m} \in A, \frac{X_{s_4}}{a_m} \in B\right)$. Then,

$$\frac{|S_n|\lambda_n^2}{m_n} E\left(\hat{\tau}_{AB,m}(h_1)\hat{\tau}_{AB,m}(h_2)\right) \qquad (5.13)$$

$$= \frac{m_n\lambda_n^2}{\nu^4|S_n|} \iiint_{S_4^4} w_n(h_1+s_1-s_2)w_n(h_2+s_3-s_4) \frac{\tau_m^*(s_1,s_2,s_3,s_4)}{m_n} E[N^{(2)}(ds_1,ds_2)N^{(2)}(ds_3,ds_4)]$$

where $N^{(2)}(ds_1, ds_2) = N(ds_1)N(ds_2)I(s_1 \neq s_2)$ and

$$E[N^{(2)}(ds_1, ds_2)N^{(2)}(ds_3, ds_4)] = \nu^4 ds_1 ds_2 ds_3 ds_4 + \nu^3 ds_1 ds_2 \varepsilon_{s_1}(ds_3) ds_4 + \nu^3 ds_1 ds_2 \varepsilon_{s_2}(ds_3) ds_4 + \nu^3 ds_1 ds_2 ds_3 \varepsilon_{s_1}(ds_4) + \nu^3 ds_1 ds_2 ds_3 \varepsilon_{s_2}(ds_4) + \nu^2 ds_1 ds_2 \varepsilon_{s_1}(ds_3) \varepsilon_{s_2}(ds_4) + \nu^2 ds_1 ds_2 \varepsilon_{s_1}(ds_4) \varepsilon_{s_2}(ds_3)$$

$$(5.14)$$

(see Karr (1986)). Now, let I_i for i = 1, ..., 7, be the integral in (5.13) corresponding to these seven scenarios of (5.14). The only cases that contribute to a non-zero limit are I_1, I_6 , and I_7 . For example, if $h_1 = h_2$,

$$I_{6} = \frac{m_{n}\lambda_{n}^{2}}{\nu^{4}|S_{n}|} \iiint_{S_{n}^{4}} w_{n}(h_{1} + s_{1} - s_{2})w_{n}(h_{2} + s_{3} - s_{4}) \frac{\tau_{m}^{*}(s_{1}, s_{2}, s_{3}, s_{4})}{m_{n}} \nu^{2} ds_{1} ds_{2} \varepsilon_{s_{1}}(ds_{3}) \varepsilon_{s_{2}}(ds_{4})$$

$$= \frac{\lambda_{n}^{2}}{\nu^{2}|S_{n}|} \iint_{S_{n}^{2}} w_{n}(h_{1} + s_{1} - s_{2})w_{n}(h_{1} + s_{1} - s_{2})\tau_{AB,m}(s_{2} - s_{1}) ds_{1} ds_{2}$$

$$= \frac{\lambda_{n}^{2}}{\nu^{2}} \int_{\frac{S_{n} - S_{n} + h_{1}}{\lambda_{n}}} \frac{1}{\lambda_{n}^{2}} w(y)^{2} \tau_{AB,m}(h_{1} - \lambda_{n}y) \frac{|S_{n} \cap (S_{n} + h_{1} - \lambda_{n}y)|}{|S_{n}|} dy$$

$$\rightarrow \frac{1}{\nu^{2}} \left(\int_{\mathbb{R}^{2}} w(y)^{2} dy \right) \tau_{AB}(h_{1})$$

$$(5.15)$$

by taking $y = \frac{h_1 + s_1 - s_2}{\lambda_n}$ and $u = s_2$ in the last equation. The convergence is from the dominated convergence theorem. On the other hand, if $h_1 \neq h_2$,

$$I_{6} = \frac{\lambda_{n}^{2}}{\nu^{2}} \int_{\frac{S_{n} - S_{n} + h_{1}}{\lambda_{n}}} \frac{1}{\lambda_{n}^{2}} w(y) w\left(y + \frac{h_{2} - h_{1}}{\lambda_{n}}\right) \tau_{AB,m}(h_{1} - \lambda_{n}y) \frac{|S_{n} \cap (S_{n} + h_{1} - \lambda_{n}y)|}{|S_{n}|} dy \to 0.$$

Similarly,

$$I_7 \to \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{A \cap BA \cap B}(h_1). \tag{5.16}$$

Turning to I_1 , we claim

$$\left| I_1 - \frac{|S_n| \lambda_n^2}{m_n} E\left(\hat{\tau}_{AB,m}(h_1)\right) E\left(\hat{\tau}_{AB,m}(h_2)\right) \right| \to 0.$$
 (5.17)

To see this, observe that the left-hand side in (5.17) is bounded by

$$\frac{m_{n}\lambda_{n}^{2}}{\nu^{4}|S_{n}|}\iiint_{S_{n}^{4}}w_{n}(h_{1}+s_{1}-s_{2})w_{n}(h_{2}+s_{3}-s_{4})$$

$$\left|\frac{\tau_{m}^{*}(\mathbf{0},s_{2}-s_{1},s_{3}-s_{1},s_{4}-s_{1})}{m_{n}}-\frac{\tau_{AB,m}(s_{2}-s_{1})}{m_{n}}\frac{\tau_{AB,m}(s_{4}-s_{3})}{m_{n}}\right|\nu^{4}ds_{1}ds_{2}ds_{3}ds_{4}$$

$$\leq \lambda_{n}^{2}m_{n}\iiint_{(S_{n}-S_{n})^{3}}w_{n}(h_{1}-v_{1})w_{n}(h_{2}-(v_{3}-v_{2}))\left|\frac{\tau_{m}^{*}(\mathbf{0},v_{1},v_{2},v_{3})}{m_{n}}-\frac{\tau_{AB,m}(v_{1})}{m_{n}}\frac{\tau_{AB,m}(v_{3}-v_{2})}{m_{n}}\right|dv_{1}dv_{2}dv_{3}$$

where the change of variables $v_1 = s_2 - s_1$, $v_2 = s_3 - s_1$, $v_3 = s_4 - s_1$ are used. By taking $u = v_2$, $y_1 = \frac{h_1 - v_1}{\lambda_n}$ and $y_2 = \frac{h_2 - (v_3 - v_2)}{\lambda_n}$, the right-hand side of the inequality is equivalent to

$$\lambda_{n}^{2} m_{n} \int_{\frac{(S_{n} - S_{n}) - (S_{n} - S_{n}) + h_{2}}{\lambda_{n}}} \int_{\frac{S_{n} - S_{n} + h_{1}}{\lambda_{n}}} \int_{S_{n} - S_{n}} w(y_{1}) w(y_{2}) \\
\left| \frac{\tau_{m}^{*}(\mathbf{0}, h_{1} - y_{1}\lambda_{n}, u, u + h_{2} - y_{2}\lambda_{n})}{m_{n}} - \frac{\tau_{AB,m}(h_{1} - y_{1}\lambda_{n})}{m_{n}} \frac{\tau_{AB,m}(h_{2} - y_{2}\lambda_{n})}{m_{n}} \right| du dy_{1} dy_{2} \\
= \lambda_{n}^{2} m_{n} O\left(\int_{\mathbb{R}^{2}} \alpha_{2,2}(||y||) dy\right) \tag{5.18}$$

To see (5.18), observe that $\min d(\{\mathbf{0}, h_1 - y_1 \lambda_n\} \{u, u + h_2 - y_2 \lambda_n\}) \le ||u|| + ||u - h_1 + y_1 \lambda_n|| + ||u + h_2 - y_2 \lambda_n|| + ||u + h_2 - y_2 \lambda_n| - h_1 + y_1 \lambda_n||$. Thus, the integral in (5.18) is bounded by

$$\begin{split} &\int_{\mathbb{R}^2} \alpha_{2,2}(||u||) du \left(\int_{\mathbb{R}^2} w(y_1) dy_1 \right)^2 + \int_{\frac{S_n - S_n + h_1}{\lambda_n}} \int_{S_n - S_n} w(y_1) \alpha_{2,2}(||u - h_1 + y_1 \lambda_n||) du dy_1 \int_{\mathbb{R}^2} w(y_2) dy_2 \\ &+ \int_{\frac{S_n - S_n + h_2}{\lambda_n}} \int_{S_n - S_n} w(y_2) \alpha_{2,2}(||u - h_2 + y_2 \lambda_n||) du dy_2 \int_{\mathbb{R}^2} w(y_1) dy_1 \\ &+ \int_{\frac{S_n - S_n + h_2}{\lambda_n}} \int_{\frac{S_n - S_n + h_1}{\lambda_n}} \int_{S_n - S_n} w(y_1) w(y_2) \alpha_{2,2}(||u + h_2 - h_1 - y_2 \lambda_n + y_1 \lambda_n||) du dy_1 dy_2 \\ &= A_1 + A_2 + A_3 + A_4. \end{split}$$

Notice that $A_1 = \int_{\mathbb{R}^2} \alpha_{2,2}(||u||) du$. Take $x = u - h_1 + y_1 \lambda_n$, then

Similarly $A_3 \leq \int_{\mathbb{R}^2} \alpha_{2,2}(||x||) dx$ can be shown. Using the similar change of variable technique,

Hence, (5.18) is verified, and (5.17) is proved.

Lastly, using the same argument in Lemma A.4. in Li et al. (2008), we have

$$I_i \to 0$$
, if $j = 2, 3, 4, 5$.

Combining the result (5.15)-(5.17), (ii) is proved, which completes the proof.

Next, we establish the asymptotic normality for $\hat{\tau}_{AB,m}(h)$.

Proposition 5.6. Assume that the conditions of Proposition 5.5 and (M3) hold. Then

$$\sqrt{\frac{|S_n|\lambda_n^2}{m_n}} \left(\hat{\tau}_{AB,m}(h) - E\hat{\tau}_{AB,m}(h)\right) \to N(0,\sigma^2),$$

where $\sigma^2 = \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{AB}(h)$. Furthermore, if $E\hat{\tau}_{AB,m}(h) - \tau_{AB}(h) = o\left(\sqrt{\frac{|S_n|\lambda_n^2}{m_n}}\right)$,

$$\sqrt{\frac{|S_n|\lambda_n^2}{m_n}} \left(\hat{\tau}_{AB,m}(h) - \tau_{AB}(h)\right) \to N(0,\sigma^2).$$

Proof. We follow Li et al. (2008) with focusing our attention to \mathbb{R}^2 and using a classical blocking technique. Let D_n^i be non-overlapping cubes that divide S_n for $i=1,...,k_n$, where $k_n=|S_n|/|D_n^i|$. Within each D_n^i , B_n^i is an inner cube sharing the same center and $d(\partial D_n^i, B_n^i) \geq n^\eta$. Let $|D_n^i| = n^{2\alpha}$ and $|B_n^i| = (n^\alpha - n^\eta)^2$ where $6/(2+\epsilon) < \eta < \alpha < 1$ for some $\epsilon > \frac{2+4\alpha}{\eta}$. Let k_n' be the additional number of cubes to cover S_n . From Lemma A.3. in Li et al. (2008),

$$k_n = O(n^{2(1-\alpha)})$$
 and $k'_n = O(n^{1-\alpha})$. (5.19)

Now define

$$A_{n} = \sqrt{\frac{m_{n}\lambda_{n}^{2}}{|S_{n}|}} \frac{1}{\nu^{2}} \iint_{S_{n} \times S_{n}} w_{n}(h + s_{1} - s_{2}) I\left(\frac{X_{s_{1}}}{a_{m}} \in A\right) I\left(\frac{X_{s_{2}}}{a_{m}} \in B\right) N^{(2)}(ds_{1}, ds_{2}),$$

$$a_{ni} = \sqrt{\frac{m_{n}\lambda_{n}^{2}}{|S_{n}|}} \frac{1}{\nu^{2}} \iint_{B_{n}^{i} \times B_{n}^{i}} w_{n}(h + s_{1} - s_{2}) I\left(\frac{X_{s_{1}}}{a_{m}} \in A\right) I\left(\frac{X_{s_{2}}}{a_{m}} \in B\right) N^{(2)}(ds_{1}, ds_{2}),$$

$$= \frac{1}{\sqrt{k_{n}}} \sqrt{\frac{m_{n}\lambda_{n}^{2}}{|D_{n}^{i}|}} \frac{1}{\nu^{2}} \iint_{B_{n}^{i} \times B_{n}^{i}} w_{n}(h + s_{1} - s_{2}) I\left(\frac{X_{s_{1}}}{a_{m}} \in A\right) I\left(\frac{X_{s_{2}}}{a_{m}} \in B\right) N^{(2)}(ds_{1}, ds_{2}),$$

$$\tilde{A}_{n} = A_{n} - EA_{n}, \quad \tilde{a}_{ni} = a_{ni} - Ea_{ni}, \quad a_{n} = \sum_{i=1}^{k_{n}} a_{ni}, \quad \tilde{a}_{n} = \sum_{i=1}^{k_{n}} \tilde{a}_{ni}, \quad \tilde{a}'_{n} = \sum_{i=1}^{k_{n}} \tilde{a}'_{ni},$$

where \tilde{a}'_{ni} denotes an independent copy of \tilde{a}_{ni} .

Step 1. Show $var(\tilde{A}_n - \tilde{a}_n) \to 0$.

We will prove Step 1 by showing:

- i) $\operatorname{var}(\tilde{A}_n) \to \frac{1}{n^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{AB}(h),$
- ii) $\operatorname{cov}(\tilde{A}_n, \tilde{a}_n) \to \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{AB}(h)$, and
- iii) $\operatorname{var}(\tilde{a}_n) \to \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{AB}(h)$.

- i) This follows from Proposition 5.5 (iii).
- ii) Recall $\tau_m^*(s_1, s_2, s_3, s_4)$ defined in Proposition 5.5 (ii). Then

$$E(A_{n}a_{n})$$

$$= \frac{\lambda_{n}^{2}}{\nu^{4}|S_{n}|} \sum_{i=1}^{k_{n}} \iiint_{S_{n} \times B_{n}^{i} \times B_{n}^{i}} w_{n}(h + s_{1} - s_{2}) w_{n}(h + s_{3} - s_{4}) \tau_{m}^{*}(s_{1}, s_{2}, s_{3}, s_{4}) E[N^{(2)}(ds_{1}, ds_{2})N^{(2)}(ds_{3}, ds_{4})]$$

$$= \frac{\lambda_{n}^{2}}{\nu^{4}|S_{n}|} \sum_{i=1}^{k_{n}} \left[\iiint_{S_{n} \times B_{n}^{i} \times B_{n}^{i} \times B_{n}^{i}} \cdot + \iiint_{S_{n} \times B_{n}^{i} \times B_{n}^{i$$

where D_i^j be the integral in D_i corresponding to the seven cases of $E[N^{(2)}(ds_1, ds_2)N^{(2)}(ds_3, ds_4)]$ as in (5.14) for i = 1, ..., 4 and j = 1, ..., 7. As shown in the proof of Proposition 5.5 (ii), non-zero contributions only arise when j = 1, 6, and 7. By the similar arguments in (5.17),

$$\left| \sum_{i=1}^{4} D_i^1 - E(A_n) E(a_n) \right| \to 0.$$

Since j=6 and 7 only occur when $s_1, s_2, s_3, s_4 \in B_n^i$, we only consider $D_4^6 + D_4^7$ which equals to

$$\frac{\lambda_n^2}{\nu^4 |S_n|} \sum_{i=1}^{k_n} \iint_{B_n^i \times B_n^i} \left[w_n(h+s_1-s_2)^2 + w_n(h+s_1-s_2)w_n(h+s_2-s_1) \right] \tau_{AB,m}(s_2-s_1)\nu^2 ds_1 ds_2$$

$$= \frac{m_n \lambda_n^2}{\nu^2 |D_n^1|} \iint_{B_n^1 \times B_n^1} \left[w_n(h+s_1-s_2)^2 + w_n(h+s_1-s_2)w_n(h+s_2-s_1) \right] \tau_{AB,m}(s_2-s_1) ds_1 ds_2$$

$$\rightarrow \frac{1}{v^2} \int_{\mathbb{R}^2} w(y)^2 dy \, \tau_{AB}(h).$$

The convergence is derived from arguments in (5.15) and (5.17). Thus, we conclude

$$cov(\tilde{A}_n, \tilde{a}_n) = \left(\sum_{i=1}^4 \sum_{j=1}^7 D_i^j\right) - E(A_n)E(a_n) = D_4^6 + D_4^7 + o(1) \to \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy\right) \tau_{AB}(h).$$

iii) Let $var(\tilde{a}_n) = \sum_{i=1}^{k_n} var(\tilde{a}_{ni}) + \sum_{1 \leq i \neq j \leq k_n} cov(\tilde{a}_{ni}, \tilde{a}_{nj})$. Note from Proposition 5.5 (iii) that

$$\sum_{i=1}^{k_n} \operatorname{var}(\tilde{a}_{ni}) = k_n \operatorname{var}(a_{n1}) \to \frac{1}{\nu^2} \left(\int_{\mathbb{R}^2} w(y)^2 dy \right) \tau_{AB}(h).$$

Also note that since \tilde{a}_{ni} and \tilde{a}_{nj} are integrals over disjoint sets for $i \neq j$ and X_s is independent of N,

 $E[\tilde{a}_{ni}|N]$ and $E[\tilde{a}_{nj}|N]$ are independent. Thus,

$$\begin{split} \sum_{1 \leq i \neq j \leq k_n} |\text{cov}(\tilde{a}_{ni}, \tilde{a}_{nj})| &= \sum_{1 \leq i \neq j \leq k_n} |E\{\text{cov}(\tilde{a}_{ni}, \tilde{a}_{nj}|N)\} + \text{cov}\{E(\tilde{a}_{ni}|N), E(\tilde{a}_{nj}|N)\}| \\ &= \sum_{1 \leq i \neq j \leq k_n} |E\{\text{cov}(\tilde{a}_{ni}, \tilde{a}_{nj}|N)\}|. \end{split}$$

Notice from Proposition 5.1 and $|a_{ni}| \leq \sqrt{\frac{m_n \lambda_n^2}{|S_n|}} |B_n^i|$ that

$$E\{\text{cov}(\tilde{a}_{ni}, \tilde{a}_{nj}|N)\} \leq const \ \frac{m_n \lambda_n^2}{|S_n|} |B_n^i| |B_n^j| \ |E(\alpha_{M,M}(n^{\eta})|N)| \leq const \ \frac{m_n \lambda_n^2}{|S_n|} |B_n^1|^2 \ E(M^2) n^{-\epsilon \eta} |B_n^1| |B_n^2| |B$$

where $M = \max\{N(B_n^i), N(B_n^j)\}$ and the last inequality is from (2.16). Since $k_n = |S_n|/|D_n^1|$ where $|S_n| = n^2, |D_n^1| = n^{2\alpha}, |B_n^1| = O(n^{2\alpha}),$

$$\sum_{1 \le i \ne j \le k_n} |\text{cov}(\tilde{a}_{ni}, \tilde{a}_{nj})| \le const \ k_n^2 \frac{m_n \lambda_n^2}{|S_n|} |B_n^1|^2 \ |B_n^1|^2 n^{-\epsilon \eta} = O\left(m_n \lambda_n^2 n^{2+4\alpha - \epsilon \eta}\right)$$

which converges to 0 as $m_n \lambda_n^2 \to 0$ and $\epsilon > \frac{2+4\alpha}{\eta}$

Step 2. Show $|\phi_n(x) - \phi'_n(x)| \to 0$ where $\phi_n(x)$ and $\phi'_n(x)$ are the characteristic functions of \tilde{a}_n and \tilde{a}'_n . Analogously to the idea presented in (6.2) in Davis and Mikosch (2009)

$$|\phi_n(x) - \phi'_n(x)| = \left| \sum_{l=1}^{k_n} E \prod_{j=1}^{l-1} e^{ix \frac{\bar{a}_{nj}}{\sqrt{k_n}}} \left(e^{ix \frac{\bar{a}_{nl}}{\sqrt{k_n}}} - e^{ix \frac{\bar{a}'_{nl}}{\sqrt{k_n}}} \right) \prod_{j=l+1}^{k_n} e^{ix \frac{\bar{a}'_{nj}}{\sqrt{k_n}}} \right| \le \sum_{l=1}^{k_n} \left| \operatorname{cov} \left(\prod_{j=1}^{l-1} e^{ix \frac{\bar{a}_{nj}}{\sqrt{k_n}}}, e^{ix \frac{\bar{a}_{nl}}{\sqrt{k_n}}} \right) \right|$$

Using the same technique in Step 1 iii),

$$\left|\operatorname{cov}\left(\prod_{j=1}^{l-1}e^{ix\frac{\tilde{a}_{nj}}{\sqrt{k_n}}},e^{ix\frac{\tilde{a}_{nl}}{\sqrt{k_n}}}\right)\right| \leq \operatorname{const} E\left(\alpha_{M,M}(n^{\eta})\right) \leq \operatorname{const} E\left(M^2\right)n^{-\epsilon\eta} \leq \operatorname{const} l^2|n^{2\alpha}|^2n^{-\epsilon\eta}$$

where $M = N(\bigcup_{j=1}^{l} B_n^j)$. The second and the last inequality is from (2.16) and $|B_n^1| = O(n^{2\alpha})$ respectively. Hence, from $k_n = n^{2-2\alpha}$, we have

$$|\phi_n(x) - \phi'_n(x)| \le const \sum_{l=1}^{k_n} l^2 |n^{2\alpha}|^2 n^{-\epsilon\eta} \le O(n^{6-2\alpha-\epsilon\eta})$$

which converges to 0 from $6/(2+\epsilon) < \eta < \alpha < 1$.

Step 3. Show the central limit theorem holds for
$$\tilde{a}'_n$$
.
Let $I_{ni} = \int_{B_n^i} \int_{B_n^i} w_n(h+s_1-s_2) I\left(\frac{X_{s_1}}{a_m} \in A\right) I\left(\frac{X_{s_2}}{a_m} \in B\right) N^{(2)}(ds_1, ds_2)$. By (2.14), we have

$$E|\sqrt{k_n}\tilde{a}'_{ni}|^{2+\delta} = E\left|\sqrt{\frac{m_n\lambda_n^2}{|D_n^i|}} \frac{1}{\nu^2} \left[I_{ni} - E(I_{ni})\right]\right|^{2+\delta}$$

$$= E\left|\sqrt{\frac{|B_n^i|^2\lambda_n^2}{|D_n^i|m_n}} \left[\hat{\tau}_{AB,m}(h:B_n^i) - E(\hat{\tau}_{AB,m}(h:B_n^i))\right]\right|^{2+\delta} < C_{\delta}$$

As (\tilde{a}'_{ni}) is triangular array of independent random variables with $var(\sum_{i=1}^{k_n} \tilde{a}'_{ni}) = \sigma_n^2 \to \sigma^2$, and

$$\frac{\sum_{i=1}^{k_n} E|\tilde{a}'_{ni}|^{2+\delta}}{(\sigma_n)^{2+\delta}} \le \frac{k_n k_n^{-(1+\delta/2)} C_\delta}{(\sigma_n)^{2+\delta}} \to 0,$$

Lyapunov's condition is satisfied and hence the central limit theorem holds.

Proof of Theorem 2.3. Proposition 5.4 implies $\hat{p}_m(A) \xrightarrow{p} \mu(A)$. By Slutsky's theorem and Proposition 5.6,

$$\sqrt{\frac{|S_n|\lambda_n^2}{m_n}}\left(\frac{\hat{\tau}_{AB,m}(h)}{\hat{p}_m(A)}-\frac{\tau_{AB,m}(h)}{\hat{p}_m(A)}\right) = \sqrt{\frac{|S_n|\lambda_n^2}{m_n}}\left(\hat{\rho}_{AB,m}(h)-\frac{\tau_{AB,m}(h)}{\hat{p}_m(A)}\right) \quad \rightarrow \quad N(0,\sigma^2/\mu(A)^2).$$

Recall from Proposition 5.4 that $var(\hat{p}_m(A)) = O(m_n/|S_n|)$. Then

$$\sqrt{\frac{|S_n|\lambda_n^2}{m_n}} \left(\hat{\rho}_{AB,m}(h) - \frac{\tau_{AB,m}(h)}{\hat{p}_m(A)} \right) = \sqrt{\frac{|S_n|\lambda_n^2}{m_n}} \left(\hat{\rho}_{AB,m}(h) - \rho_{AB,m}(h) \right) + o_p(1).$$

Thus, the central limit theorem for $\sqrt{\frac{|S_n|\lambda_n^2}{m_n}}$ $(\hat{\rho}_{AB,m}(h) - \rho_{AB,m}(h))$ is proved. The joint normality (2.18) is established using the Cramér-Wold device.

5.3 Appendix C: Example 3.4

First, we show that X_s satisfies LUNC in (2.17). Notice that the process has continuous sample paths a.s. since the Gaussian process $\{W_s - \delta(s), s \in \mathbb{R}^2\}$ in (3.9) has continuous sample paths. Notice from Lindgren (2012), Section 2.2, that a Gaussian process with a continuous correlation function satisfying (3.13) has continuous sample paths.

From (3.9), let $X_s=U^1_s\vee U^2_s$, where $U^1_s=\Gamma^{-1}_1Y^1_s$ and $U^2_s=\sup_{j\geq 2}\Gamma^{-1}_jY^j_s$. Then

$$nP\left(\sup_{||s||<\delta'} \frac{|X_s - X_{\mathbf{0}}|}{a_n} > \delta\right) = nP\left(\sup_{||s||<\delta'} |U_s^1 \vee U_s^2 - U_{\mathbf{0}}^1 \vee U_{\mathbf{0}}^2| > a_n\delta\right)$$

$$\leq nP\left(\sup_{||s||<\delta'} |U_s^1 - U_{\mathbf{0}}^1| > \frac{a_n\delta}{2}\right) + nP\left(\sup_{||s||<\delta'} |U_s^2 - U_{\mathbf{0}}^2| > \frac{a_n\delta}{2}\right)$$

$$= A_1 + A_2.$$

Since $E|\sup_{|s|=\delta'}|Y(s)|<\infty$ (see Proposition 13 in Kabluchko et al. (2009)), we can apply the dominated convergence theorem to obtain

$$A_1 = nP\left(\Gamma_1 < \frac{2\sup_{||s|| < \delta'} |Y_s^1 - Y_{\mathbf{0}}^1|}{\delta a_n}\right) = n\int \left(1 - e^{-z/\delta a_n}\right) g(Z)dZ \to \frac{2E(\sup_{||s|| < \delta'} |Y_s - Y_{\mathbf{0}}|)}{\delta} \to 0,$$

where $Z = 2 \sup_{||s|| < \delta'} |Y_s - Y_0|$.

To show $A_2 \to 0$, we follow the arguments in Davis and Mikosch (2008).

$$\begin{split} A_2 &= nP\left(\sup_{||s|| < \delta'} \bigvee_{j \geq 2}^{\infty} \Gamma_j^{-1} |Y_s^j - Y_s^j| > \frac{a_n \delta}{2}\right) &\leq n \sum_{j=2}^{\infty} P\left(2\sup_{||s|| < \delta'} |Y_s| > \Gamma_j \delta a_n / 2\right) \\ &= n \int \left(\sum_{j \geq 2}^{\infty} P\left(4y > \Gamma_j \delta a_n\right)\right) P\left(\sup_{||s|| < \delta'} |Y_s| \in dy\right) \\ &= n \int_0^{\infty} \left(\frac{4y}{\delta a_n} - \left(1 - e^{-\frac{4y}{\delta a_n}}\right)\right) P\left(\sup_{||s|| < \delta'} |Y_s| \in dy\right) \end{split}$$

The last line is from $ET[0, \frac{4y}{\delta a_n}] = \sum_{j=1}^{\infty} P\left(\Gamma_j < \frac{4y}{\delta a_n}\right) = \frac{4y}{\delta a_n}$, where $T = \sum_{j=1}^{\infty} \epsilon_{\Gamma_j}$ is a homogeneous point process. The dominated convergence theorem applies as $f_n(y) = n\left(\frac{4y}{\delta a_n} - (1 - e^{-\frac{4y}{\delta a_n}})\right) \le cy$ for some c > 0, all y > 0 and $f_n(y) \to 0$ as $n \to 0$, and $E\sup_{||s|| < \delta'} |Y_s| < \infty$ from Kabluchko et al. (2009).

Now we check conditions (2.11)-(2.16). Recall from (3.12) that $\alpha_{c,c}(h) \leq const \frac{1}{\sqrt{||h||^{\alpha}}} e^{-\theta||h||^{\alpha}/2}$ holds for the process. For convenience in the calculations that follow, set $g(h) = \frac{1}{\sqrt{||h||^{\alpha}}} e^{-\theta||h||^{\alpha}/2}$. We will find the sufficient conditions for (2.11)-(2.16). For (2.11),

$$\int_{\mathbb{R}^2} g(y)dy < \infty \tag{5.20}$$

is sufficient. To see this, infer from (3.11) that

$$m_n P(X_y > \epsilon a_m, X_0 > \epsilon a_m) = m_n \left[1 - 2e^{-1/a_m} + e^{-2\Phi(\sqrt{\delta(h)})/a_m} \right] = \frac{2m_n}{a_m} \bar{\Phi}(\sqrt{\delta(h)}) + O\left(\frac{m_n}{a_m^2}\right).$$

Thus

$$\begin{split} m_n \int_{B[k,r_n]} P(X_y > \epsilon a_m, X_{\mathbf{0}} > \epsilon a_m) dy &= \int_{B[k,r_n]} \frac{2m_n}{a_m} \bar{\Phi}(\sqrt{\delta(y)}) dy + O\left(\frac{r_n^2}{m_n}\right) \\ &\leq const \int_{B[k,\infty]} g(y) dy + o(1), \end{split}$$

where the last inequality is from (3.12).

From (3.12), the condition (2.12) is satisfied if

$$\int_{\mathbb{R}^2 \setminus B[0,r_n)} m_n g(y) dy \to 0. \tag{5.21}$$

Similarly, using (3.12), the second condition in (2.15) is implied if (5.20) holds. The condition (2.16) is checked immediately from (3.1) since

$$\sup_{l} \frac{\alpha_{l,l}(||h||)}{l^2} \le const \ \frac{1}{\sqrt{||h||^{\alpha}}} e^{-\theta||h||^{\alpha}/2} = O(||h||^{-\epsilon}).$$

We check the condition (2.14) with $\delta = 1$ is satisfied if (3.14) assumed, but we skip this as it is tedious. Hence, it suffices to find conditions under which (5.20) and (5.21) hold.

Remark 6. If the process is regularly varying in the space of continuous functions in every compact set, then

LUNC is satisfied. See Hult and Lindskog (2006), Theorem 4.4.

Proposition 5.7. For Example 3.4, the conditions (5.20) - (5.21) hold if $\log m_n = o(r_n^a)$.

Proof. Using change of variables to polar coordinates and $r^a/8 = t$, (5.20) is checked. For $a \in (0,2]$

$$\int_{\mathbb{R}^2} g(y) dy = const \ \int_0^\infty t^{\frac{2}{a} - \frac{3}{2}} e^{-t} dt < \infty.$$

For (5.21), notice that for sufficiently large n, $m_n g(r_n) \leq m_n e^{-\theta r_n^{\alpha}/2} = o(1)$ provided $\log m_n = o(r_n^a)$. This completes the proof.

Finally, we find the condition under which (2.19) holds.

Proposition 5.8. For the Brown-Resnick process, (2.19) holds if $\frac{|S_n|\lambda_n^2}{m_n^3} \to 0$.

Proof. From (3.11),

$$|\rho_{AB,m}(h) - \rho_{AB}(h)| = \frac{1 + o(1)}{\mu(h)} |\tau_{AB,m}(h)\mu(h) - \tau_{AB}(h)p_m(h)| = \frac{1 + o(1)}{\mu(h)} O(m_n/a_m^2) = O(1/m_n).$$

Therefore, (2.19) holds if
$$\frac{|S_n|\lambda_n^2}{m_n^3} \to 0$$
.

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