THE CLUSTER INDEX OF REGULARLY VARYING SEQUENCES WITH APPLICATIONS TO LIMIT THEORY FOR FUNCTIONS OF MULTIVARIATE MARKOV CHAINS

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ABSTRACT. We introduce the cluster index of a multivariate regularly varying stationary sequence and characterize the index in terms of the spectral tail process. This index plays a major role in limit theory for partial sums of regularly varying sequences. We illustrate the use of the cluster index by characterizing infinite variance stable limit distributions and precise large deviation results for sums of multivariate functions acting on a stationary Markov chain under a drift condition.

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1. Introduction

Consider a stationary Markov chain $(\Phi_t)$ and a function $f$ acting on the state space of the Markov chain and mapping into $\mathbb{R}^d$ for some $d \geq 1$. For the resulting stationary process $X_t = f(\Phi_t), t \in \mathbb{Z}$, the corresponding partial sum process is given by

$$ S_0 = 0, \quad S_n = X_1 + \cdots + X_n. $$

We also assume that the finite-dimensional distributions of the process $(X_t)$ are regularly varying with index $\alpha$; see Section 2.1 for a definition. Roughly speaking, this condition ensures that the tails of the finite-dimensional distributions have power law behavior, hence sufficiently high moments of $X$ are infinite. (Here and in what follows, we write $Y$ for a generic element of any stationary sequence $(Y_t)$.) Regular variation of a random vector and, more generally, of a stationary sequence is a condition which determines the extremal dependence structure in a flexible way.

For an iid sequence the condition of regular variation of $X$ with index $\alpha \in (0, 2)$ is necessary and sufficient for the central limit theorem

$$ a_n^{-1} (S_n - b_n) \xrightarrow{d} \xi_\alpha, \quad n \to \infty, $$

where $a_n > 0, b_n \in \mathbb{R}, n \in \mathbb{N}$, are suitable constants and $\xi_\alpha$ has an $\alpha$-stable distribution in $\mathbb{R}^d$; see [49] for the limit theorem and [50] for a description of infinite variance stable laws in $\mathbb{R}^d$. Limit theory with $\alpha$-stable limits for dependent sequences was studied in [28, 29] by using the convergence of characteristic functions and in [10] by using the continuous mapping theorem acting on suitable weakly converging point processes; see also [5] for a functional central limit theorem using the same technique. These results were proved for univariate sequences, but [17] proved that the point process convergence results remain valid in the multivariate case by a slight modification of the proofs in [10].
Using ideas from [23, 29], the authors of [2] studied stable limit theory for general univariate regularly varying sequences; see Theorem 6.1 below. We use this result and the Cramér-Wold device to derive the corresponding limits for all linear combinations \( \theta'S_n, \theta \in \mathbb{R}^{d-1} \), where \( \mathbb{R}^{d-1} \) is the unit sphere in \( \mathbb{R}^d \) with respect to the Euclidean norm. According to Theorem 6.1 the \( \alpha \)-stable limit laws of \( \theta'S_n \) (with suitable normalization and centering) are characterized by the function
\[
b(\theta) = \lim_{k \to \infty} \lim_{x \to \infty} \frac{\mathbb{P}(\theta'S_k > x) - \mathbb{P}(\theta'S_{k-1} > x)}{\mathbb{P}(|X| > x)}, \quad \theta \in \mathbb{R}^{d-1}.
\]
We discuss the so-called cluster index \( b \) in Section 3. The existence of the limits in (1.1) is guaranteed under the conditions of this paper; see Theorem 3.2. Moreover, the cluster index \( b \) determines the \( \alpha \)-stable limit laws in the multivariate case; see Theorem 4.1. In a way, the function \( b \) plays a similar role as the notion of extremal index in limit theory for maxima of dependent sequences; see [35] for this notion.

Regular variation is also the key to precise large deviation theory for the sums \( S_n \). In the univariate iid case, classical work by A.V. and S.V. Nagaev [11, 12] shows that relations of the following type hold
\[
\sup_{x \geq b_n} \left\{ \frac{\mathbb{P}(S_n > x)}{n \mathbb{P}(|X| > x)} - p \right\} \to 0,
\]
where \( p = \lim_{x \to \infty} \mathbb{P}(X > x)/\mathbb{P}(|X| > x) = p \) and \( (b_n) \) is a suitably chosen sequence such that \( b_n \to \infty \) and \( S_n/b_n \overset{P}{\to} 0 \) as \( n \to \infty \). Related work for dependent regularly varying sequences was proved in [29] for linear processes, in [13, 34] for solutions to stochastic recurrence equations and for general regularly varying sequences in [40]; for earlier results see also [28, 29, 16]. The results in these papers are all of the type
\[
\sup_{x \in (b_n, c_n)} \left| \frac{\mathbb{P}(S_n > x)}{n \mathbb{P}(|X| > x)} - b(1) \right| \to 0,
\]
where \( b(1) \) is the limit in (1.1) for \( d = 1 \) and \( (b_n, c_n) \) are suitable regions tending to infinity.

The case of iid multivariate regularly varying \( (X_t) \) was treated in [27], including a corresponding functional large deviation result. In this paper, we get a corresponding large deviation principle for regularly varying multivariate functions acting on a Markov chain (see Theorem 1.3):
\[
\frac{\mathbb{P}(\lambda_n^{-1}S_n \in \cdot)}{n \mathbb{P}(|X| > \lambda_n)} \overset{\nu_n}{\to} \nu_\alpha.
\]
Here \( \overset{\nu_n}{\to} \) denotes vague convergence on some Borel \( \sigma \)-field, \( \lambda_n \to \infty \) is a suitable normalizing sequence and the limit \( \nu_\alpha \) is a measure which is induced by the regular variation of the sums \( S_k, \ k \geq 1 \). As for the case of stable limits, we start by proving the large deviation principle for linear combinations \( \theta'S_n \), exploiting the corresponding result (1.2) with \( b(1) \) replaced by \( b(\theta) \) from (1.1); see Theorem 7.2. This corresponds to (1.3) restricted to half-planes not containing the origin. It is in general not possible to extend the limit relation (1.3) from half-spaces to general Borel sets. This extension is however possible by assuming some additional conditions such as \( \alpha \) is non-integer. As a matter of fact, relation (1.3) cannot be written as a uniform result in the spirit of (1.2), due to its multivariate character.

The paper is organized as follows. In Section 2 we introduce regular variation of a stationary sequence and the drift condition of a Markov chain. In Section 3 we define the cluster index \( b(\theta) \), \( \theta \in \mathbb{R}^{d-1} \), of a regularly varying stationary sequence. We prove the existence of the cluster index for multivariate functions acting on a Markov chain under a drift condition (Theorem 3.2). In Section 4 we formulate the main results of this paper. They include \( \alpha \)-stable limit theory (Theorem 4.1) and precise large deviation principles (Theorem 4.2) for functions of regenerative Markov chains.
In Section 5 we calculate the cluster index for several important time series models, including multivariate autoregressive processes, solutions to stochastic recurrence equations, GARCH(1,1) processes and their sample covariance functions. In the remaining sections we prove the results of Section 4.

2. Preliminaries

2.1. Regular variation of vectors and sequences of random vectors. In what follows, we will use the notion of regular variation as a suitable way of describing heavy tails of random vectors and sequences of random vectors. We commence with a random vector $X$ with values in $\mathbb{R}^d$ for some $d \geq 1$. We say that this vector (and its distribution) are regularly varying with index $\alpha > 0$ if the following relation holds as $x \to \infty$:

$$\frac{\mathbb{P}(x |X| > uX, |X| \in \cdot)}{\mathbb{P}(|X| > x)} \xrightarrow{w} u^{-\alpha} \mathbb{P}(\Theta \in \cdot), \quad u > 0.$$  

(2.1)

Here $\xrightarrow{w}$ denotes weak convergence of finite measures and $\Theta$ is a vector with values in the unit sphere $S^{d-1} = \{ x \in \mathbb{R}^d : |x| = 1 \}$ of $\mathbb{R}^d$. Its distribution is the spectral measure of regular variation and depends on the choice of the norm. However, the definition of regular variation does not depend on any concrete norm; we always refer to the Euclidean norm. An equivalent way to define regular variation of $X$ is to require that there exists a non-null Radon measure $\mu$ on the Borel $\sigma$-field of $\mathbb{R}_0^d = \mathbb{R}^d \setminus \{0\}$ such that

$$n \mathbb{P}(a_n^{-1}X \in \cdot) \xrightarrow{w} \mu_X(\cdot),$$

where the sequence $(a_n)$ can be chosen such that $n \mathbb{P}(|X| > a_n) \sim 1$ and $\xrightarrow{w}$ refers to vague convergence. The limit measure $\mu_X$ necessarily has the property $\mu_X(u \cdot) = u^{-\alpha} \mu_X(\cdot), u > 0$, which explains the relation with the index $\alpha$. We refer to [7] for an encyclopedic treatment of one-dimensional regular variation and [38, 47] for the multivariate case.

Next consider a strictly stationary sequence $(X_t)_{t \in \mathbb{Z}}$ of $\mathbb{R}^d$-valued random vectors with a generic element $X$. It is regularly varying with index $\alpha > 0$ if every lagged vector $(X_{t-1}, \ldots, X_k)$, $k \geq 1$, is regularly varying in the sense of (2.1); see [10]. An equivalent description of a regularly varying sequence $(X_t)$ is achieved by exploiting (2.2): for every $k \geq 1$, there exists a non-null Radon measure $\mu_k$ on the Borel $\sigma$-field of $\mathbb{R}_0^{dk}$ such that

$$n \mathbb{P}(a_n^{-1}(X_1, \ldots, X_k) \in \cdot) \xrightarrow{w} \mu_k,$$

(2.2)

where $(a_n)$ is chosen such that $n \mathbb{P}(|X_0| > a_n) \sim 1$.

A convenient characterization of a regularly varying sequence $(X_t)$ was given in Theorem 2.1 of [6]: there exists a sequence of $\mathbb{R}^d$-valued random vectors $(Y_t)_{t \in \mathbb{Z}}$ such that $\mathbb{P}(|Y_0| > y) = y^{-\alpha}$ for $y > 1$ and for $k \geq 0$,

$$\mathbb{P}(x^{-1}(X_{t-k}, \ldots, X_k) \in \cdot | |X_0| > x) \xrightarrow{w} \mathbb{P}(Y_{t-k}, \ldots, Y_k) \in \cdot), \quad x \to \infty.$$

(2.3)

The process $(Y_t)$ is the tail process of $(X_t)$. Writing $\Theta_t = Y_t/|Y_0|$ for $t \in \mathbb{Z}$, one also has for $k \geq 0$,

$$\mathbb{P}(|X_0|^{-1}(X_{t-k}, \ldots, X_k) \in \cdot | |X_0| > x) \xrightarrow{w} \mathbb{P}((\Theta_{t-k}, \ldots, \Theta_k) \in \cdot), \quad x \to \infty.$$  

(2.4)

We will identify $|Y_0| (Y_t/|Y_0|)_{|t| \leq k} = |Y_0| (\Theta_t)_{|t| \leq k}$, $k \geq 0$. Then $|Y_0|$ is independent of $(\Theta_t)_{|t| \leq k}$ for every $k \geq 0$. We refer to $(\Theta_t)_{t \in \mathbb{Z}}$ as the spectral tail process of $(X_t)$.

We formulate our main condition on the tails of the sequence $(X_t)$:

**Condition (RV)k**: The strictly stationary sequence $(X_t)$ is regularly varying with index $\alpha > 0$ and spectral tail process $(\Theta_t)$.
2.2. The drift condition. Assume that the following drift condition holds for the Markov chain \((\Phi_t)\) for suitable \(p > 0\) and an \(\mathbb{R}^d\)-valued function \(f\) acting on the state space of the Markov chain:

Condition \((\text{DC}_p)\): There exist constants \(\beta \in (0, 1), b > 0\), and a function \(V : \mathbb{R}^d \to (0, \infty)\) such that \(c_1|x|^p \leq V(x) \leq c_2|x|^p, c_1, c_2 > 0\), satisfying for any \(y\) in the state space of the Markov chain,

\[
\mathbb{E}(V(f(\Phi_1)) | \Phi_0 = y) \leq \beta V(f(y)) + b.
\]

We mention that Jensen’s inequality ensures that \((\text{DC}_p)\) implies \((\text{DC}_{p'})\) for \(p' < p\). We exploited condition \((\text{DC}_p)\) in [40], where we proved large deviation principles for regularly varying strictly stationary sequences of random variables, in particular for irreducible Markov chains.

If \((\Phi_t)\) is an irreducible Markov chain then \((\text{DC}_p)\) for any \(p > 0\) implies \(\beta\)-mixing with geometric rate; see [39], p. 371. Moreover, without loss of generality, by considering the Nummelin splitting scheme, see [15] for details, we will assume that \((\Phi_t)\) possesses an atom \(A\). The notions of drift, small set, atom, etc. used throughout are borrowed from [39]. In what follows, we write \(\mathbb{P}_A(\cdot) = \mathbb{P}(\cdot | \Phi_0 \in A)\) and \(\mathbb{E}_A\) for the corresponding expectation.

We always assume the existence of some \(M > 0\) such that \(\{x : V(f(x)) \leq M\}\) is a small set (this is true in all our examples). Then the condition \((\text{DC}_p)\) is equivalent to the existence of constants \(\beta \in (0, 1)\) and \(b > 0\) such that for any \(y\),

\[
\mathbb{E}(V(f(\Phi_1)) | \Phi_0 = y) \leq \beta V(f(y)) + b \mathbb{I}_A(y).
\]

Direct verification of the condition \((\text{DC}_p)\) is in general difficult. We will use the following result which can often be checked much easier.

Lemma 2.1. Assume that the stationary Markov chain \((\Phi_t)\) is aperiodic, irreducible and satisfies the following condition for some \(p > 0\) and integer \(m \geq 1\):

Condition \((\text{DC}_{p,m})\): (a) There exist \(b > 0\) and \(\beta \in (0, 1)\) such that for any \(y\) in the state space of the Markov chain,

\[
\mathbb{E}(V(f(\Phi_m)) | \Phi_0 = y) \leq \beta V(f(y)) + b \mathbb{I}_A(y),
\]

where \(V\) is the function from \((\text{DC}_p)\).

(b) There exist \(c_1, c_2 > 0\) such that for any \(y\) in the state space of the Markov chain

\[
\mathbb{E}(V(f(\Phi_1)) | \Phi_0 = y) \leq c_1 V(f(y)) + c_2.
\]

Then condition \((\text{DC}_p)\) holds.

Proof. Theorem 15.3.3 in [39] says that the drift condition in part (a) of \((\text{DC}_{p,m})\) implies \(V\)-geometric regularity of the \(m\)-skeleton Markov chain \((\Phi_{tm})\). Theorem 15.3.6 in [39] yields the equivalence between \(V\)-geometric regularity and \(g\)-geometric regularity of the original Markov chain for a function \(g\) satisfying \(\sum_{i=1}^m \mathbb{E}(g(\Phi_t) | \Phi_0 = y) = V(f(y))\). Thus the drift condition is satisfied for the original Markov chain and some finite Lyapunov function \(V' \geq g\). Making multiple use of part (b) of \((\text{DC}_{p,m})\), we can show that there exist constants \(c'_1, c'_2 > 0\) satisfying \(\sum_{i=1}^m \mathbb{E}(g(\Phi_t) | \Phi_0 = y) \leq c'_1 V(f(y)) + c'_2\). Thus \((\text{DC}_p)\) follows for a function \(V'(x) = c''_1 V(x) + c''_2\) and suitable constants \(c'_1, c''_2 > 0\). 

Consider the sequence of the hitting times of the atom \(A\) by the Markov chain \((\Phi_t)\), i.e. \(\tau_A(1) = \tau_A = \min\{k > 0 : \Phi_k \in A\}\) and \(\tau_A(j + 1) = \min\{k > \tau_A(j) : \Phi_k \in A\}, j \geq 1\). We will write

\[
S(0) = \sum_{t=1}^{\tau_A} X_t \quad \text{and} \quad S(i) = \sum_{t=\tau_A(i)+1}^{\tau_A(i+1)} X_t, \quad i \geq 1.
\]
According to the theory in \[31\], \((\tau_A(i) - \tau_A(i - 1))_{i \geq 2}\) and \((S(i))_{i \geq 1}\) constitute iid sequences; we will refer to regenerative Markov chains. The drift condition \((\text{DC}_p)\) is tailored for proving the existence of moments of \(S(1)\) under the existence of moments of \(X_t = f(\Phi_t)\) of the same order.

The drift condition \((\text{DC}_p)\) is useful for proving central limit theory and other asymptotic results for functions of Markov chains. As a benchmark result we quote a central limit theorem which is a simple corollary of Proposition 2.1 in Samur \[51\]. To apply this result notice that \((\text{DC}_1)\) implies condition \((\text{D}_2)\) of \[21\] for \(|X_t|\) with \(V = c|f|\) with \(c > 0\) sufficiently small.

**Theorem 2.2.** Assume that the stationary Markov chain \((\Phi_t)\) is aperiodic, irreducible and \((X_t) = (f(\Phi_t))\) satisfies \((\text{DC}_1)\), \(\mathbb{E}|X|^2 < \infty\) and \(\mathbb{E}X = 0\). Then the following statements hold:

1. The partial sum \(S(1)\) has finite second moment.
2. The central limit theorem \(n^{-0.5}S_n \stackrel{d}{\to} \mathcal{N}(0, \Sigma)\) holds with
   \[
   \Sigma = \mathbb{E}_A[S(1)S(1)'] = \lim_{k \to \infty} \mathbb{E} \left[ \left( \sum_{i=0}^{k} X_i \right) \left( \sum_{i=0}^{k} X_i' \right)' - \left( \sum_{i=1}^{k} X_i \right) \left( \sum_{i=1}^{k} X_i' \right)' \right].
   \]

Together with Theorem \[4.1\] that deals with the case of infinite variance stable limits, Theorem \[2.2\] complements the limit theory for partial sums of functions of Markov chains in the case of finite variance summands and Gaussian limits.

3. **The cluster index**

We commence by considering a general \(\mathbb{R}^d\)-valued stationary process \((X_t)\) satisfying \((\text{RV}_\alpha)\) for some \(\alpha > 0\). A continuous mapping argument for regular variation (see e.g. \[24, 25\]) and \[2.3\] ensure the existence of the limits

\[
b_k(\theta) = \lim_{n \to \infty} n \mathbb{P}(\theta' S_k > a_n), \quad k \geq 1, \quad \theta \in \mathbb{S}^{d-1}.
\]

The difference \(b_{k+1}(\theta) - b_k(\theta)\) can be expressed in terms of the spectral tail process \((\Theta_t)\) of \((X_t)\).

**Lemma 3.1.** Let \((X_t)\) be an \(\mathbb{R}^d\)-valued stationary process satisfying \((\text{RV}_\alpha)\) for some \(\alpha > 0\). Then, for any \(k \geq 1\),

\[
b_{k+1}(\theta) - b_k(\theta) = \mathbb{E} \left[ \left( \theta' \sum_{i=0}^{k} \Theta_i \right)_+^\alpha - \left( \theta' \sum_{i=1}^{k} \Theta_i \right)_+^\alpha \right].
\]

**Proof.** We start by observing that each \(b_k(\theta)\) can be expressed in terms of the spectral tail process \((\Theta_t)\). Indeed, \((\text{RV}_\alpha)\) yields for every \(k \geq 1\) and \(\theta \in \mathbb{S}^{d-1}\) that

\[
b_k(\theta) = \lim_{x \to \infty} \frac{\mathbb{P}(\theta' S_k > x)}{\mathbb{P}(|X| > x)} = \lim_{x \to \infty} \frac{\mathbb{P}\left( \bigcup_{j=1}^{k} \{ \theta' S_k > x, \theta' X_j > x/k \} \cap \{ \theta' X_i < x/k, 1 \leq i < j \} \right)}{\mathbb{P}(|X| > x)} = \lim_{x \to \infty} \sum_{j=1}^{k} \left( \frac{\mathbb{P}(\theta' S_k > x, \theta' X_j > x/k)}{\mathbb{P}(|X| > x)} - \frac{\mathbb{P}(\theta' S_k > x, \theta' X_j > x/k, \max_{1 \leq i < j} \theta' X_i > x/k)}{\mathbb{P}(|X| > x)} \right).
\]
By stationarity, the summands in the above expression can be written in the form
\[
\frac{\mathbb{P}(|X_0| > x/k)}{\mathbb{P}(|X_0| > x)} \left[ \mathbb{P}\left( \theta' \sum_{t=1-j}^{k-j} X_t > x, \theta' X_0 > x/k \mid |X_0| > x/k \right) - \mathbb{P}\left( \theta' \sum_{t=1-j}^{k-j} X_t > x, \max_{1-j \leq s < 0} \theta' X_s > x/k \mid |X_0| > x/k \right) \right].
\]

Here we used the fact that \{\theta' X_0 > x/k\} \subset \{|X_0| > x/k\}. Letting \(x \to \infty\) in the above expressions, applying the conditional limits \([24]\) and observing that \(\mathbb{P}(|Y_0| > y) = y^{-\alpha}, y > 1\), we obtain the limiting expressions
\[
k^{\alpha} \left[ \mathbb{P}\left( Y_0 \theta' \sum_{t=1-j}^{k-j} \Theta_t > k, |Y_0| \theta' \Theta_0 > 1 \right) - \mathbb{P}\left( Y_0 \theta' \sum_{t=1-j}^{k-j} \Theta_t > 1, |Y_0| \theta' \Theta_0 > 1, \max_{1-j \leq s < 0} \theta' \Theta_s > 1 \right) \right] = \mathbb{E} \left[ \left( \theta' \sum_{t=1-j}^{k-j} \Theta_t \right)^\alpha + \left( k \theta' \Theta_0 \right)^\alpha \right] - \mathbb{E} \left[ \left( \theta' \sum_{t=1-j}^{k-j} \Theta_t \right)^\alpha + \left( k \theta' \Theta_0 \right)^\alpha \wedge \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right].
\]

Hence \(b_k(\theta)\) has representation
\[
b_k(\theta) = \sum_{j=1}^{k} \mathbb{E} \left[ \left( \theta' \sum_{t=1-j}^{k-j} \Theta_t \right)^\alpha - \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right] + \left( k \theta' \Theta_0 \right)^\alpha - \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right],
\]
and therefore
\[
b_{k+1}(\theta) - b_k(\theta) = \mathbb{E} \left[ \left( \theta' \sum_{t=0}^{k} \Theta_t \right)^\alpha + \left( k \theta' \Theta_0 \right)^\alpha \right] + \sum_{j=1}^{k} \mathbb{E} \left[ \left( \theta' \sum_{t=1-j}^{k-j} \Theta_t \right)^\alpha - \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right] + \left( k \theta' \Theta_0 \right)^\alpha - \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right] - \mathbb{E} \left[ \left( \theta' \sum_{t=1-j}^{k-j} \Theta_t \right)^\alpha - \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right] + \left( k \theta' \Theta_0 \right)^\alpha - \max_{1-j \leq s < 0} \left( k \theta' \Theta_s \right)^\alpha \right].
\]

The expectations in the sum are of the type \(\mathbb{E} f(\Theta_{-s}, \ldots, \Theta_t)\) for integrable \(f\) such that \(f(x_{-s}, \ldots, x_t) = 0\) if \(x_{-s} = 0\), \(s, t \geq 0\). Then, according to Theorem 3.1 (iii) in \([6]\),
\[
\mathbb{E} f(\Theta_{-s}, \ldots, \Theta_t) = \mathbb{E} \left( f(\Theta_{0/\Theta_s}, \ldots, \Theta_{t+s/\Theta_s}) | \Theta_s |^\alpha \right), \ s, t \geq 0.
\]

Application of this formula and the fact that our functions \(f\) are homogeneous of order \(\alpha\) yield
\[
b_{k+1}(\theta) - b_k(\theta) = \mathbb{E} \left[ \left( \theta' \sum_{t=0}^{k} \Theta_t \right)^\alpha + \left( k \theta' \Theta_0 \right)^\alpha \right] + \sum_{j=1}^{k} \mathbb{E} \left[ \left( \theta' \sum_{t=0}^{k} \Theta_t \right)^\alpha - \max_{0 \leq i < j} \left( k \theta' \Theta_i \right)^\alpha \right] + \left( k \theta' \Theta_j \right)^\alpha - \max_{0 \leq i < j} \left( k \theta' \Theta_i \right)^\alpha \right].
\]
Theorem 3.2. Assume that $b_{k+1}(\theta) = \lim_{k \to \infty} (b_{k+1}(\theta) - b_k(\theta))$, $\theta \in \mathbb{R}^{d-1}$. In contrast to the quantities $b_k(\theta)$ the existence of the limits $b(\theta)$ is not straightforward. The following result yields a sufficient condition for the existence of $b$.

**Theorem 3.2.** Assume that $(X_t)$ satisfies (RV$_\alpha$) for some $\alpha > 0$ and that $X_t = f(\Phi_t)$, $t \in \mathbb{Z}$, where $f$ is an $\mathbb{R}^d$-valued function acting on the Markov chain $(\Phi_t)$ satisfying (DC$_p$) for some positive
\( p \in (\alpha - 1, \alpha). \) Then the limits
\[
b(\theta) = \mathbb{E}\left[ \left( \sum_{t \geq 0} \theta^t \Theta_t \right)^\alpha - \left( \sum_{t \geq 1} \theta^t \Theta_t \right)^\alpha \right], \quad \theta \in S^{d-1},
\]
exist and are finite.

**Remark 3.3.** The cluster index \( b \) of \((X_t)\) is a continuous function on \( S^{d-1} \). This is shown in the proof below: \( b \) is the uniform limit of continuous functions on \( S^{d-1} \). The index \( b(\theta) \) is non-negative since it coincides with the Cesàro mean \( \lim_{k \to \infty} k^{-1} b_k(\theta) \). For \( 0 < \alpha \leq 1 \), the sub-additivity of the function \( x \to x_+^{\alpha} \) implies the inequality \( b(\theta) \leq \mathbb{E}[\theta^t \Theta_0]_+^\alpha \). Moreover, if \( \mathbb{E}[\theta^t \Theta_0]_+^\alpha > 0 \) then \( b(\theta) > 0 \) by an application of the mean value theorem when \( 0 < \alpha \leq 1 \). These two properties are shared by the extremal index of a multivariate stationary process. The extremal index admits a similar representation in terms of the spectral tail process, i.e., \( \mathbb{E}[\sup_{t \geq 0} \theta^t \Theta_t]_+^\alpha - \mathbb{E}[\sup_{t \geq 1} \theta^t \Theta_t]_+^\alpha \); see [5].

**Remark 3.4.** The limit \( b \) also exists for various classes of regularly varying stationary processes beyond functions of a Markov chain; see [2, 10] for such examples in the case \( d = 1 \). The cluster index \( b \) plays a crucial role for characterizing weak and large deviation limits for partial sums of the processes \((X_t)\). This was recognized in [2, 10], and we extend some of these results to the multivariate case in Section 4.

**Proof.** We will show that the limit \( b(\theta) \) of (3.1) exists as \( k \to \infty \). We start with the case \( \alpha > 1 \). Then, for \( x, y \in \mathbb{R} \), by the mean value theorem, \( |(x + y)_+^\alpha - x_+^\alpha| \leq (\alpha |y| |x + \xi y|^{\alpha - 1}) \lor |y|^\alpha \) for some \( \xi \in (0, 1) \). Hence, since \( |\theta^t \Theta_0| \leq 1 \) a.s.,
\[
|b_{k+1}(\theta) - b_k(\theta)| \leq \mathbb{E}\left[ \left( \alpha |\theta^t \Theta_0| \sum_{t=1}^k \theta^t \Theta_t + \xi \theta^t \Theta_0 \right)^\alpha \lor |\theta^t \Theta_0|^\alpha \right] \\
\leq \mathbb{E}\left[ \left( \alpha \sum_{t=1}^k \theta^t \Theta_t + \xi \theta^t \Theta_0 \right)^{\alpha - 1} \lor 1 \right] = I_0.
\]

For \( \alpha \in (1, 2] \),
\[
I_0 \leq 1 + \alpha \sum_{t=0}^k \mathbb{E}|\theta^t \Theta_t|^{\alpha - 1}.
\]

We will show that the right-hand side is finite, implying that \( \mathbb{E}\left[ \sum_{t=0}^\infty |\theta^t \Theta_t|^{\alpha - 1} \right] < \infty \) and \( \sum_{t=0}^\infty \theta^t \Theta_t \) converges absolutely a.s. An application of Lebesgue dominated convergence shows that the limit \( b(\theta) \) exists and is finite. For \( \alpha > 2 \), an application of Minkowski’s inequality yields
\[
I_0 \leq 1 + \alpha \left( \sum_{t=0}^k \mathbb{E}|\theta^t \Theta_t|^{\alpha - 1} \right)^{1/(\alpha - 1)} \cdot \left( \sum_{t=0}^k \mathbb{E}|\theta^t \Theta_t|^{\alpha - 1} \right)^{1/(\alpha - 1)}.
\]

We will show that the right-hand side is finite and then the same argument as for \( \alpha \in (1, 2] \) applies. We will achieve the bounds for \( I_0 \) by showing that there exists \( c > 0 \) such that
\[
(3.1) \quad \mathbb{E}|\theta^t \Theta_t|^{\alpha - 1} \leq c |\theta|^t, \quad t \geq 0.
\]

Using the fact that \( \Theta_t = Y_t / Y_0 \) and \( Y_0 \) are independent, for \( t \geq 1 \) and \( s = \alpha - 1 \),
\[
\mathbb{E}|Y_0|^s \mathbb{E}|\theta^t \Theta_t|^s \leq \mathbb{E}|Y_0|^s \mathbb{E}|\Theta_t|^s = \mathbb{E}|Y_t|^s.
\]
By definition of the tail process and Markov’s inequality, for small $\epsilon > 0$ such that $s(1 + \epsilon) < \alpha$,

$$
\mathbb{E}[Y_t]^s = \int_0^\infty \mathbb{P}(|Y_t|^s > y) \, dy \\
= \int_0^\infty \lim_{x \to \infty} \mathbb{P}(|x^{-1} Y_t|^s > y \mid |X_0| > x) \, dy \\
\leq \int_1^\infty y^{-(1+\epsilon)} \, dy \lim_{x \to \infty} \frac{\mathbb{E}[|X_t|^{s(1+\epsilon)} \mathbb{1}_{|X_0|>x}]}{x^{s(1+\epsilon)} \mathbb{P}(|X_0| > x)} \\
+ \int_0^1 y^{-(1-\epsilon)} \, dy \lim_{x \to \infty} \frac{\mathbb{E}[|X_t|^{s(1-\epsilon)} \mathbb{1}_{|X_0|>x}]}{x^{s(1-\epsilon)} \mathbb{P}(|X_0| > x)} \\
= R_1 + R_2.
$$

By virtue of (DC$_p$) for some $p \in (\alpha - 1, \alpha)$, using a recursive argument, we obtain for sufficiently large $y$ and $s(1 + \epsilon) \leq p$,

$$
\mathbb{E}[|X_t|^{s(1+\epsilon)} \mid \Phi_0 = y] \leq \beta^t |f(y)|^{s(1+\epsilon)} + b \sum_{j=1}^t \beta^j.
$$

Using this inequality and Karamata’s theorem (see [7]), for some $c > 0$,

$$
R_1 \leq c \lim_{x \to \infty} \frac{\mathbb{E}[|X_t|^{s(1+\epsilon)} \mathbb{1}_{|X_0|>x}]}{x^{s(1+\epsilon)} \mathbb{P}(|X_0| > x)} \leq c \beta^t \lim_{x \to \infty} \frac{\mathbb{E}[|X_t|^{s(1+\epsilon)} \mathbb{1}_{|X_0|>x}]}{x^{s(1+\epsilon)} \mathbb{P}(|X_0| > x)} \leq c \beta^t.
$$

Similarly, $R_2 \leq c \beta^t$. We conclude that (3.1) holds for $\alpha > 1$.

It remains to consider the case $\alpha \leq 1$. We observe that $|(x + y)^\alpha_+ - x^\alpha_+| \leq |y|^\alpha$ for any $x, y \in \mathbb{R}$.

Hence

$$
|b_{k+1}(\theta) - b_k(\theta)| \leq \mathbb{E}(|\theta|^\alpha \Theta_0|^{\alpha} \leq 1.
$$

It suffices to show that $\sum_{t=0}^\infty |\theta|^{\Theta_t} < \infty$ a.s. This follows if $\sum_{t=0}^\infty \mathbb{E}|\theta|^{\Theta_t}|s < \infty$ for some $s < p$.

The proof is analogous, using (DC$_p$) for some $p < \alpha$. □

4. LIMIT THEORY FOR FUNCTIONS OF REGENERATIVE MARKOV CHAINS

In this section we present the main results of this paper. Throughout we consider an $\mathbb{R}^d$-valued process $X_t = f(\Phi_t)$, $t \in \mathbb{Z}$, where $(\Phi_t)$ is an irreducible aperiodic Markov chain. We present two types of limit results for the partial sums $(S_n)$ of $(X_n)$: central limit theory with infinite stable limits in Theorem 4.1 and precise large deviation results in Theorem 4.3. The proofs of these results are postponed to Sections 5 and 7.

4.1. Stable limit theory. We start with a central limit theorem with stable limit law.

**Theorem 4.1.** Consider an $\mathbb{R}^d$-valued strictly stationary sequence $(X_t) = (f(\Phi_t))$ satisfying the following conditions:

- (RV$_\alpha$) for some $\alpha \in (0, 2)$, $\mathbb{E}X = 0$ if $\alpha > 1$ and $X$ is symmetric if $\alpha = 1$.
- (DC$_p$) for some $p \in ((\alpha - 1) \vee 0, \alpha)$.

Let $(a_n)$ be a sequence of positive numbers such that $n \mathbb{P}(|X_0| > a_n) \sim 1$. Then the following statements hold:
Theorem 4.2. Assume that the strictly stationary $\mathbb{R}^d$-valued sequence $(X_t)$ satisfies \((RV_\alpha)\) for some $\alpha > 0$ and the following point process convergence result holds:

$$N_n = \sum_{i=1}^{n} \delta_{a_n^{-1} X_i} \overset{d}{\rightarrow} N = \sum_{i=1}^{\infty} \sum_{j=1}^{\infty} \delta_{P_i Q_{ij}},$$

where $(P_i)$ are the points of a Poisson random measure on $(0, \infty)$ with intensity $h(y) = \gamma \alpha y^{-\alpha-1}$, $y > 0$, and it is assumed that $\gamma > 0$ and the sequence $(Q_{ij})_{j \geq 1}$, $i = 1, 2, \ldots$, is iid with values $|Q_{ij}| \leq 1$, independent of $(P_i)$ and such that $\sup_{j \geq 1}|Q_{ij}| = 1$.

1. If $\alpha \in (0, 1)$ then

$$a_n^{-1} S_n \overset{d}{\rightarrow} \xi_\alpha = \sum_{i=1}^{\infty} \sum_{j=1}^{\infty} P_i Q_{ij}$$

and $\xi_\alpha$ has an $\alpha$-stable distribution,
(2) If \( \alpha \in [1, 2) \) and for any \( \delta > 0 \),
\[
\lim_{n \to \infty} \limsup_{n \to \infty} P(|S_n(0, \varepsilon) - \mathbb{E}S_n(0, \varepsilon)| > \delta) = 0,
\]
where \( S_n(0, \varepsilon) = a_n^{-1} \sum_{i=1}^{n} X_i \mathbb{1}_{\{|X_i| \leq \varepsilon a_n\}} \), then
\[
a_n^{-1} S_n - \mathbb{E}S_n(0, 1) \xrightarrow{d} \xi,\]
where \( \xi \) is the distributional limit as \( \varepsilon \downarrow 0 \) of
\[
\left( \sum_{i=1}^{\infty} \sum_{j=1}^{\infty} P_i Q_{ij} \mathbb{1}_{(\varepsilon, \infty)}(P_i|Q_{ij}|) - \int_{|x| \leq \varepsilon} x \mu_X(dx) \right)
\]
which exists and has an \( \alpha \)-stable distribution. (Recall that \( \mu_X \) is the limit measure in (2.2).)

The latter result has been the basis for a variety of results for partial sums of strictly stationary processes with infinite variance stable limits; see [17, 37, 39, 52]. The main idea of the proof of Theorem 4.2 is a continuous mapping argument acting on \( N_n \xrightarrow{d} N \), showing that the sums of the points of \( N_n \) converge in distribution to the corresponding sum of the points of \( N \). This method is rather elegant and can be applied to a large variety of strictly stationary regularly varying vector sequences \( (X_t) \). The proofs use advanced point process techniques.

A characterization of the parameters of the distribution of the multivariate limit \( \xi \) in Theorem 4.2 can be given by extending Theorem 3.2 in [16] to the multidimensional case: if
\[
\mathbb{E}(\sum_{j \geq 1} |Q_{1j}|)^\alpha < \infty
\]
then the Lévy spectral measure \( \Gamma_\alpha \) of \( \xi \) is described by
\[
\int_{\mathbb{S}^{d-1}} (\theta')^\alpha_+ \Gamma_\alpha(ds) = \frac{\alpha}{2-\alpha} \mathbb{E}\left[ \left( \sum_{i \geq 1} \theta_i |Q_{1i}| \right)^\alpha_+ \right], \quad \theta \in \mathbb{S}^{d-1}.
\]
This representation is particularly useful for \( \alpha < 1 \). Then (4.3) is always satisfied. Adapting Theorem 4.2 in terms of the tail process as in Basrak et al. [5], an alternative characterization of the Lévy spectral measure \( \Gamma_\alpha \) is the following: if
\[
\mathbb{E}(\sum_{i \geq 0} |\Theta_i|)^\alpha < \infty
\]
then
\[
\int_{\mathbb{S}^{d-1}} (\theta')^\alpha_+ \Gamma_\alpha(ds) = C_\alpha^{-1} \mathbb{E}\left[ \left( \sum_{i \geq 0} \theta_i |\Theta_i| \right)^\alpha_+ \mathbb{1}_{\{\Theta_i = 0, \forall i \leq -1\}} \right], \quad \theta \in \mathbb{S}^{d-1}.
\]
Conditions (4.5) and (4.6) may fail for \( \alpha > 1 \), e.g. for a GARCH(1,1) model; see Section 5.4.

If we assume the conditions of Theorem 4.1, classical computation for \( \alpha \neq 1 \) yields
\[
\mathbb{E}[\exp(i v' \xi)] = \exp\left( -\int_{\mathbb{S}^{d-1}} |v'|(1 - isgn(v') \tan(\pi \alpha/2)) \Gamma_\alpha(d\theta) \right)
= \exp\left( \int_0^{\infty} \mathbb{E}\left[ \exp\left( iv \sum_{t=1}^{\infty} v' \Theta_t \right) - \exp\left( iv \sum_{t=0}^{\infty} v' \Theta_t \right) \right] ax^{-\alpha-1} dx \right).
\]
For \( \alpha \in (0, 1) \), this form of the limiting stable characteristic function was proved in Basrak and Segers [6].

The additional condition (4.4) is not easily checked for dependent sequences. It is implied for stationary \( \rho \)-mixing processes with rate function \( \rho(j) \) satisfying \( \sum_{j \geq 1} \rho(2^j) < \infty \); see [29]. It is
also implied by \((DC_p)\) for functions of an irreducible Markov chain; see \cite{10}. For a (possibly non-
irreducible) Markov chain \((X_t)\), condition \((DC_p)\) is much weaker than this \(\rho\)-mixing condition which is
equivalent to a spectral gap in \(L^2(\mathbb{P})\); see \cite{36}.

In our paper, characteristic function based methods are employed which are close to those used in
classical limit theory for iid sequences; see e.g. \cite{11}. As in the iid case, Theorem \ref{4.1} yields an
explicit form of the characteristic function of the limiting \(\alpha\)-stable random vector. The underlying
extremal dependence structure of \((X_t)\) shows via the cluster index \(b(\theta)\) which appears explicitly in the
characteristic function. We refer the reader to the extensive discussion in \cite{2} on the comparison of
the point process and the characteristic function approaches to stable limit theory. One drawback of
our approach is that, in contrast to the point process approach, we do not have series representations
of \(\xi_\alpha\) in terms of the sequence \((\Theta_t)\).

Recently, the special case of solutions to multivariate stochastic recurrence equations \cite{53} has
attracted attention; see e.g. \cite{15} \cite{12}. In this case, one can exploit the underlying random iterative
contractive structure to derive stable limits without additional restrictions. We mention that drift
conditions such as \((DC_p)\) are automatically satisfied for solutions of stochastic recurrence equations;
see Section \ref{5}

4.2. Precise large deviations for functions of a Markov chain. In this section, we extend
some of the results obtained in \cite{10} for general univariate regularly varying sequences.\footnote{For comparison and since we will use it in the proofs, we quote the main result of \cite{10} as Theorem \ref{7.2}} We again focus on \(\mathbb{R}^d\)-valued sequences \((X_t) = (f(\Phi_t))\) for an underlying aperiodic irreducible Markov chain
\((\Phi_t)\). The case \(\alpha \in (0, 2)\) turns out to be a consequence of Theorem \ref{4.1} the proof is given in
Section \ref{7.1}. The proof in the case \(\alpha > 2\) is more involved and requires different techniques; see Section \ref{7.2}

**Theorem 4.3.** Consider an \(\mathbb{R}^d\)-valued strictly stationary sequence \((X_t) = (f(\Phi_t))\) for an aperiodic
irreducible Markov chain \((\Phi_t)\). Assume that \((X_t)\) satisfies the condition \((RV_\alpha)\) for some \(\alpha > 0\). Let
\((\lambda_n)\) be any sequence such that \(\log(\lambda_n) = o(n)\) and \(\lambda_n/n^{1/\alpha+\varepsilon} \to \infty\) if \(\alpha \in (0, 2)\) and
\(\lambda_n/n^{0.5+\varepsilon} \to \infty\) if \(\alpha > 2\) for any \(\varepsilon > 0\). Assume either

1. \(\alpha \in (0, 2)\) and the conditions of Theorem \ref{4.1} are satisfied, or
2. \(\alpha > 2, \alpha \notin \mathbb{N}\) or \(b(\theta) = b(\theta), \theta \in \mathbb{S}^{-1}\), and \((DC_p)\) holds for every \(p < \alpha\),

then the following large deviation principle holds:

\begin{equation}
\frac{\mathbb{P}(\lambda_n^{-1}S_n \in \cdot)}{n \mathbb{P}(|X| > \lambda_n)} \to \nu_\alpha, \quad n \to \infty,
\end{equation}

where \(\nu_\alpha\) is a Radon measure on the Borel \(\sigma\)-field of \(\mathbb{R}^d_0\) uniquely determined by the relations

\begin{equation}
\nu_\alpha(t \{x : \theta'x > 1\}) = t^{-\alpha} \nu_\alpha(t \{x : \theta'x > 1\}) = t^{-\alpha} b(\theta), \quad \theta \in \mathbb{S}^{-1}, \quad t > 0.
\end{equation}

**Remark 4.4.** The conditions \(\alpha \notin \mathbb{N}\) or \(b(\cdot) = b(\cdot)\) are needed to apply inverse results for regular
variation. For \(\alpha > 2\), we show that the measure \(\nu_\alpha\) on the Borel \(\sigma\)-field of \(\mathbb{R}^d_0\) is uniquely determined
by its values on sets of the form \(t \{x : \theta'x > 1\}, t > 0, \theta \in \mathbb{S}^{-1}\), provided the mentioned additional
conditions are met. In general, such conditions cannot be avoided; \cite{32} \cite{20} give counterexamples
for integer values \(\alpha\). In \cite{3} \cite{8} \cite{33} further conditions on the vector \(X\) are given which allow one to
discover the measure \(\nu_\alpha\) from its knowledge on the sets \(t \{x : \theta'x > 1\}, t > 0, \theta \in \mathbb{S}^{-1}\).

**Remark 4.5.** The proof of Theorem \ref{4.3} shows that \ref{4.1} holds uniformly for certain intervals of
normalizations and for half-spaces not containing the origin. To be precise, the following uniform
relations hold
\[
(4.3) \quad \lim_{n \to \infty} \sup_{x \in \Lambda_n} \left| \frac{P(\theta' S_n > x)}{nP(|X| > x)} - b(\theta) \right| = 0, \quad \theta \in \mathbb{S}^{d-1},
\]
for regions \( \Lambda_n = (b_n, c_n) \). Here \((b_n)\) satisfies \( b_n = n^{0.5+\varepsilon} \) in the case \( \alpha > 2 \) and \( b_n = n^{1/\alpha+\varepsilon} \) in the case \( \alpha \in (0, 2) \) for any \( \varepsilon > 0 \), and \((c_n)\) is chosen such that \( c_n > b_n \) and \( \log c_n = o(n) \). Moreover, for \( (\Theta_t) \) one does not need the additional conditions \( b(-) = b(-) \) and \( \alpha \notin \mathbb{N} \).

5. Examples

Here we consider several examples of regularly varying stationary processes with index \( \alpha > 0 \), where the theory of the previous sections applies. In particular, we will determine the tail process \((\Theta_t)\), the cluster index \( b \) and verify the drift condition \((\text{DC}_p)\) for \( p < \alpha \). All models considered fall in the class of functions acting on an aperiodic irreducible Markov chain.

5.1. Vector-autoregressive process. Consider the vector-autoregressive process of order 1 given by
\[
(5.1) \quad X_t = A X_{t-1} + Z_t, \quad t \in \mathbb{Z},
\]
where \( A \) is a random \( d \times d \) matrix whose eigenvalues are less than 1 in absolute value, and \( A \) is independent of the iid \( \mathbb{R}^d \)-valued sequence \((Z_t)\) which is regularly varying with index \( \alpha > 0 \). Then we also have \( \mathbb{E}\|A\|^s < 1 \) for every \( s > 0 \). Here \( \| \cdot \| \) denotes the operator norm with respect to the Euclidean norm.

Then a stationary solution \((X_t)\) to \( (5.1) \) exists and has representation
\[
X_t = A^t X_0 + \sum_{i=1}^t A^{t-i} Z_i, \quad t \geq 0;
\]
see [11], Chapter 11. Moreover, \( X_0 \) is regularly varying with index \( \alpha \); see [38]. In particular, denoting the limiting measure of the regularly varying vector \( Z_0 \) by \( \mu_Z \), it follows from [38] that
\[
(5.2) \quad \frac{P(x^{-1} X_0 \in \cdot)}{P(|Z_0| > x)} \xrightarrow{v} \sum_{i=0}^{\infty} \mathbb{E}[\mu_Z(\{x \in \mathbb{R}^d : A^i x \in \cdot\})].
\]
Since
\[
(X_1, \ldots, X_h) = (A, \ldots, A^h) X_0 + \left( Z_1, \ldots, \sum_{i=1}^h A^{h-i} Z_i \right),
\]
and \((Z_t)_{t \geq 1}\) is independent of \( X_0 \), regular variation of \((X_1, \ldots, X_h)\) is a consequence of the fact that regular variation is kept under linear transformations. Let \( C \) be a continuity set relative to the limiting measure \( \mu_{h+1} \) of \((X_0, \ldots, X_h)\) and \( I_d \) the identity matrix. Since \( X_0 \) is independent of \((Z_t)_{t \geq 1}\),
\[
P(x^{-1}(X_0, \ldots, X_h) \in C \mid |X_0| > x) = P(x^{-1}(I_d, A, \ldots, A^h) X_0 \in C \mid |X_0| > x) + P(x^{-1}(0, Z_1, \sum_{i=1}^2 A^{2-i} Z_i, \ldots, \sum_{i=1}^h A^{h-i} Z_i \in C) + o(1)
\]
\[
\to P((I_d, A, \ldots, A^h) Y_0 \in C), \quad x \to \infty.
\]
Thus we may identify \((\Theta_t)_{t=0, \ldots, h}\) with \((I_d, A, \ldots, A^h) \Theta_0\). In view of \( (5.2) \),
\[
P(x^{-1} X_0 / |X_0| \in \cdot \mid |X_0| > x) \xrightarrow{w} \frac{\sum_{i=0}^{\infty} \mathbb{E} [\mu_Z(\{x \in \mathbb{R}^d : A^i x / |A^i x| \in \cdot, A^i x| > 1\})]}{\sum_{i=0}^{\infty} \mathbb{E} [\mu_Z(\{x \in \mathbb{R}^d : |A^i x| > 1\})]} = P(\Theta_0 \in \cdot).
\]
Writing \((I_d - A)^{-1} = \sum_{i=0}^{\infty} A^i\) (this series converges since the largest eigenvalue of \(A\) is smaller than 1), we conclude that
\[
b(\theta) = E \left[ \left( \theta'(I_d - A)^{-1} \Theta_0 \right)_+ - \left( \theta' A(I_d - A)^{-1} \Theta_0 \right)_+ \right], \quad \theta \in S^{d-1}.
\]
Next we show \((DC_p)\) for \(p < \alpha\). First assume \(p > 1\). A Taylor series expansion yields
\[
E(|Ax + Z_1|^p - |Ax|^p) \leq pE[|Z_1| |Ax + \xi Z_1|^{p-1} |x|^{p-1} + 1) \leq c(E|Ax|^{p-1} + 1) \leq c(|x|^{p-1} + 1)
\]
for some random variable \(\xi \in (0, 1)\) a.s. Then for some \(\beta \in (E|A|^p, 1)\) and sufficiently large \(|x|\),
\[
E|Ax + Z_1|^p \leq E|Ax|^p + c(1 + |x|^{p-1}) \leq E|A|^p|x|^p(1 + c|x|^{-1}) + c \leq \beta|x|^p + c,
\]
and \((DC_p)\) is satisfied. If \(p \leq 1\) a simpler argument applies with \(\beta = E|A|^p\):
\[
E(|Ax + Z_1|^p) \leq E|Ax|^p + E|Z_1|^p \leq \beta|x|^p + c.
\]
If the Markov chain \((X_t)\) is also aperiodic and irreducible the results in Section 5 are directly applicable with \(f(x) = x\).

5.2. Random affine mapping. Following Kesten [32], we consider the stochastic recurrence equation
\[
X_t = A_t X_{t-1} + B_t, \quad t \in \mathbb{Z},
\]
where \(((A_t, B_t))_{t \in \mathbb{Z}}\) is an iid sequence, \(A_t\) are random \(d \times d\)-matrices and \(B_t\) are \(\mathbb{R}^d\)-valued random vectors. We also assume \(E \log^+ |A| < \infty\), where \(\| \cdot \|\) denotes the operator norm with respect to the Euclidean norm, \(E \log^+ |B| < \infty\), and that the Lyapunov exponent of the stochastic recurrence equation \((5.3)\) is negative. These conditions ensure that an a.s. unique stationary causal solution \((X_t)\) to \((5.3)\) exists; see [9, 11]. Under additional regularity conditions which ensure that the distribution of \(A\) is sufficiently spread out, the equation
\[
\theta(\kappa) = \lim_{n \to \infty} n^{-1} \log E \| A_1 \cdots A_n \|^\kappa = 0, \quad \kappa > 0,
\]
has a unique positive solution \(\alpha\) and \(\theta' X, \theta \in S^{d-1}\), is regularly varying with index \(\alpha\). Under stronger conditions on \(A\), \(\alpha\) can be calculated as the solution to \(E\|A\|^\kappa = 1, \kappa > 0\); see [20, 13] for recent results. Kesten [32] had already given conditions which ensured that at least one of the linear combinations \(\theta' X, \theta \in S^{d-1}\), is regularly varying with index \(\alpha\). In general, one cannot conclude from regular variation of \(\theta' X, \theta \in S^{d-1}\), that \(X\) is regularly varying; see [32] for some counterexamples. In [3, 33, 34] conditions are given which ensure that the regular variation of a vector can be recovered from the regular variation of its linear projections. One of these conditions is that \(\alpha \notin \mathbb{N}; \) see [3] for details. In what follows, we will assume that \(X_t\) is regularly varying with index \(\alpha > 0\) and that the stronger moment conditions \(E\|A\|^{2(\alpha + \epsilon)} < \infty\) and \(E\|B\|^{2(\alpha + \epsilon)} < \infty\) hold for some \(\epsilon > 0\). If \(A_t\) and \(B_t\) are independent the milder moment conditions \(E\|A\|^\alpha < \infty\) and \(E\|B\|^\alpha < \infty\) for some \(\epsilon > 0\) suffice.

Calculation yields
\[
X_t = \Pi_t X_0 + R_t, \quad \text{where} \quad \Pi_t = A_t \cdots A_1, \quad t \geq 1,
\]
where \(E \| R_t \|^{\alpha + \epsilon} < \infty\) and hence
\[
\mathbb{P}(x^{-1}(X_0, \ldots, X_t) \in \cdot \mid \|X_0\| > x) \xrightarrow{w} \mathbb{P}(\|Y_0\| (I_d, \Pi_1, \ldots, \Pi_t) \Theta_0 \in \cdot).
\]
where \(\mathbb{P}(X_0/\|X_0\| \in \cdot \mid \|X_0\| > x) \xrightarrow{w} \mathbb{P}(\Theta_0 = \cdot)\) is independent of \((A_t)_{t \geq 1}\). Therefore \((\Theta_{i=0,\ldots,t}) = (I_d, \Pi_1, \ldots, \Pi_t) \Theta_0\). Writing \(\Pi_0 = I_d\), the identity matrix in \(\mathbb{R}^d\), and \((Z_t)\) for the

\footnote{However, it might be possible to prove regular variation of \(X_t\) by using the structure of [32].}
solution of the stochastic recurrence equation (5.3) in the special case $B = I_d$, we obtain from Theorem 3.2

$$b(\theta) = E\left[\left(\theta' \sum_{t \geq 0} \Pi_t \Theta_0\right)_+^\alpha - \left(\theta' \sum_{t \geq 1} \Pi_t \Theta_0\right)_+^\alpha\right]$$

(5.6)

provided we can show (DC) for the Markov chain $(\Phi_t) = (X_t)$. The formula (5.6) is in agreement with the calculations for $d = 1$ in [2].

Since (5.4) is satisfied we can use Lemma 2.1 for proving (DC). Assuming irreducibility and aperiodicity of the Markov chain $(X_t)$ and exploiting the definition of $\alpha$ as solution to (5.4), one can choose $m \geq 1$ sufficiently large such that $E[|A_1 \cdots A_m|^p < 1$ for any $p \leq \alpha$. Indeed, assume on the contrary that $\rho(p) \geq 0$ for some $p < \alpha$. This contradicts the convexity of $\rho$ which has roots at 0 and $\alpha$. Then the $m$-skeleton of the chain satisfies the drift condition, (DC) follows and Lemma 2.1 yields (DC). Thus we conclude that the results of Section 4 are directly applicable to the Markov chain $(X_t)$ with $f(x) = x$ if it is also aperiodic and irreducible.

5.3. Sample autocovariance function of one-dimensional random affine mapping. Consider the solution $(X_t)$ to the stochastic recurrence equation (5.3) in the case $d = 1$, under irreducibility and aperiodicity. We assume the conditions and use the notation of Section 5.2. In addition, we write

$$\Pi_{s,t} = \begin{cases} A_s \cdots A_t & s \leq t, \\ 1 & \text{otherwise.} \end{cases}$$

In particular, we assume that $(X_t)$ is regularly varying with index $\alpha > 0$ satisfying $E[|A|^\alpha = 1$, $E[|A|^\alpha + \varepsilon < \infty$ and $E[|B|^\alpha + \varepsilon < \infty$ for some $\varepsilon > 0$. For $h \geq 0$, consider the process $\Phi_t = (X_t, X_{t-1}, \ldots, X_{t-h})'$, $t \in \mathbb{Z}$, of lagged vectors. They constitute an $\mathbb{R}^{h+1}$-valued stationary, aperiodic and irreducible Markov chain. Similar arguments as in Section 5.2 show that the chain is regularly varying with index $\alpha > 0$. We consider the following function acting on the Markov chain $(\Phi_t)$:

$$X_t = f(\Phi_t) = (X_t, X_{t-1}, X_{t-2}, \Phi_t), \quad t \in \mathbb{Z}.$$ 

By convention, we will assume that all vectors are understood as column vectors. The sequence $(\Phi_t)$ satisfies the recursion

$$\Phi_t = \begin{pmatrix} A_t & 0 & \cdots & 0 \\ 1 & 0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 \end{pmatrix} \Phi_{t-1} + \begin{pmatrix} B_t \\ 0 \\ \vdots \\ 0 \end{pmatrix} = A_t \Phi_{t-1} + B_t, \quad t \in \mathbb{Z}.$$ 

We will show that $(X_t)$ satisfies (RV) and (DC) for $m$ sufficiently large, $V(x) = |x|^p$ and $p < \alpha/2$. The condition $E[V(f(\Phi_t)) | \Phi_0 = y] \leq c_1 V(f(y)) + c_2$ for some positive $c_1, c_2$ follows immediately from the stochastic recurrence equation

$$X_t = \begin{pmatrix} A_t A_{t-1} & A_t B_t & B_t A_{t-1} \\ 0_{h+1, h+1} & A_t & 0_{h+1, h+1} \\ 0_{h+1, h+1} & 0_{h+1, h+1} & A_t \end{pmatrix} X_{t-1} + \begin{pmatrix} B_t B_t \\ B_t \\ B_t \end{pmatrix} = C_t X_{t-1} + D_t, \quad t \in \mathbb{Z}.$$ 

Condition $E[D]^{(\alpha + \varepsilon)/2} < \infty$ follows by the assumptions. From basic algebra, for $m \geq h$ the matrix products $\prod_{t=1}^m A_t A_t = \Pi_t \prod_{t=1}^m A_t$ can be written as $\Pi_{m,h} \mathbf{B}_h$, where the $(h+1) \times (h+1)$ matrix
\( M_h \) has zero entries but the first column given by \((\Pi_{1,h-1}, \Pi_{2,h-1}, \ldots, 1)\). Products of triangular matrices remain triangular and their diagonal is the product of the diagonals. Thus we obtain

\[
C_m \cdots C_1 = \begin{pmatrix}
\Pi_{m,h}^2 I_h & 0_{1,h+1} & 0_{h+1,h+1} \\
0_{1,h+1} & \Pi_{m,h} & 0_{h+1,h+1} \\
0_{h+1,h+1} & 0_{h+1,h+1} & \Pi_{m,h} I_h \\
\end{pmatrix} \quad \bar{C}_h = D_m \bar{C}_h,
\]

where \( \bar{C}_h \) is an upper triangular block matrix depending only on \((A_t)_{1 \leq t \leq h-1}\). The matrices \( \bar{D}_m \) and \( \bar{C}_h \) are independent and for some \( c > 0 \) we have

\[
E\|\bar{D}_m \bar{C}_h\|^p \leq c E\|\bar{D}_m\|^p E\|\bar{C}_h\|^p.
\]

Since \( p < \alpha/2 \), \( E(|A_0|^{2p}) \to 0 \) and \( E(|A_0|^p) \to 0 \) as \( m \to \infty \). Thus, for \( m \) sufficiently large, \( E\|\bar{D}_m \bar{C}_h\|^p \leq c E\|A_0\|^{2p} + E(|A_0|^p) < 1 \), i.e., condition \((D_{C_{p,m}})\) holds, and Lemma \((2.1)\) applies provided we can also show \((RV_{\alpha/2})\) for \((X_t)\). This is our next goal. Since \( X_t \) and \( \Phi_t \) are regularly varying with index \( \alpha \) we deal with a degenerate case where the limiting measure of regular variation of \( X_t \) is concentrated at zero for the last \( h + 2 \) components. Then, in view of the definition of the cluster index, \( b \) is the same for \((X_t, \Phi_t)\) and \((X_t)\). Therefore we will calculate \( b \) for \((X_t, \Phi_t)\). Abusing notation, we will also use the same notation for the tail process. As in Section \( 5.2 \) we obtain by iteration of the stochastic recurrence equation \( X_t = A_t X_{t-1} + B_t \),

\[
X_t \Phi_t = \Pi_{t-h,t} (\Pi_{t-h+1,t}, \Pi_{t-h+2,t-1}, \ldots, \text{1})^t X_{t-h} + R_{t+h}^{(1)}
\]

\[
\Pi_{1-h,t} (\Pi_{1-h+1,t}, \Pi_{1-h+2,t-1}, \ldots, \text{1})^t X_{t-h} + R_{t+h}^{(2)},
\]

(5.7)

where \( E|R_i^{(1)}|^{(\alpha+c)/2} < \infty \), \( i = 1, 2 \). Then for \( t \geq 0 \),

\[
(X_0, \ldots, X_t \Phi_t) = \begin{pmatrix}
\Pi_{1-h,0}^2 & \Pi_{1-h,1}^2 & \cdots & \Pi_{1-h,t}^2 \\
\Pi_{1-h,0} \Pi_{1-h,1} & \Pi_{1-h,1} \Pi_{1-h,0} & \cdots & \Pi_{1-h,t} \Pi_{1-h,0} \\
\vdots & \vdots & \ddots & \vdots \\
\Pi_{1-h,0} & \Pi_{1-h,1} & \cdots & \Pi_{1-h,t} \\
\end{pmatrix} X_{t-h} + Q_t, \quad t \in \mathbb{Z}.
\]

and \( E|Q_t|^{(\alpha+c)/2} < \infty \). In the remainder of this section we assume that \( P(A = 0) = 0 \); the general case can be treated as well but leads to tedious case studies. An application of Corollary 3.2 in Basrak and Segers \((6)\) yields that for continuity sets \( M \),

\[
P(x^{-1}(X_0, \ldots, X_t \Phi_t) \in M \mid |X_0 \Phi_0| > x) \to P(|Y_0| \in E_t \in M) \),
\]

where

\[
E_t \overset{\text{d}}{=} \frac{1}{|\Pi_h|} \sqrt{\Pi_h^2 + \Pi_{h-1}^2 + \cdots + 1} \begin{pmatrix}
\Pi_h \Pi_h & \Pi_h \Pi_{h+1}^2 & \cdots & \Pi_{t+h} \Pi_{t+h} \\
\Pi_h \Pi_{h-1} & \Pi_{h+1} \Pi_h & \cdots & \Pi_{t+h} \Pi_{t+h-1} \\
\vdots & \vdots & \ddots & \vdots \\
\Pi_h & \Pi_{h+1} \Pi_1 & \cdots & \Pi_{t+h} \Pi_t \\
\end{pmatrix}
\]

and \( E_t \) is independent of \( |Y_0| \). The right-hand side can be identified with \((\Theta_0, \ldots, \Theta_t)\).

An application of Theorem \((4.1)\) now yields a stable limit result for the sample autocovariance function of \((X_t)\). Assume that \((a_n)\) satisfies \( nP(|X_0 \Phi_0| > a_n) \sim 1 \). In view of \((5.7)\) and Breiman’s result (see \((10)\)) we also have

\[
nP(|X_0 \Phi_0| > a_n) \sim nP(X^2 > a_n) E\left[\left(\sqrt{|\Pi_h|} \sqrt{1 + \Pi_h^2 + \cdots + \Pi_{h-1}^2}\right)^{\alpha/2}\right].
\]
In view of Kesten’s result \cite{Kesten1969}, \( \mathbb{P}(|X| > x) \sim c_0 x^{-\alpha} \). Therefore we can choose
\[
a_n = n^{2/\alpha} (c_0 \mathbb{E} \left[ (\Pi_h \sqrt{1 + \Pi_1 + \cdots + \Pi_n \alpha/2}) \right]^{2/\alpha} .
\]
Then we have for \( m \geq 0, \alpha \in (2, 4) \),
\[
\left( a_n^{-1} \sum_{i=1}^{n-h} (X_i X_{t+h} - \mathbb{E}(X_0 X_h)) \right)_{h=0, \ldots, m} \xrightarrow{d} \xi_{\alpha/2} ,
\]
and for \( \alpha \in (0, 2) \),
\[
\left( a_n^{-1} \sum_{i=1}^{n-h} X_i X_{t+h} \right)_{h=0, \ldots, m} \xrightarrow{d} \xi_{\alpha/2} ,
\]
where \( \xi_{\alpha/2} \) is an \( \alpha/2 \)-stable \( \mathbb{R}^{h+1} \)-valued random vector whose characteristic function is given in Theorem 4.1 and \( (\Theta_i) \) is described above. This result was proved in Basrak et al. \cite{Basrak2004}, Theorem 2.13. In the case \( \alpha \in (2, 4) \) the additional condition \( (2.20) \) was needed; the latter condition is hardly verifiable and could be overcome in the present paper by showing condition \( (DC_\rho) \). Moreover, as in \cite{Basrak2004} a straightforward application of the continuous mapping theorem yields a corresponding limit result for the sample autocorrelation function; we omit details. The limit laws in Theorem 2.13 of \cite{Basrak2004} are expressed in terms of the points of the limiting point processes in Theorem 4.1 above, while our limits are expressed in terms of the cluster index \( b \). Neither of the representations of the \( \alpha \)-stable limits are easy due to the complicated dependence structure.

5.4. Sample mean of a GARCH(1, 1) process and its volatility, sample covariance function of a GARCH(1, 1) process. We consider a GARCH(1, 1) process \( X_t = \sigma_t Z_t \), where \( (Z_t) \) is an iid sequence of mean zero unit variance random variables and \( (\sigma_t) \) is a sequence of non-negative random variables such that \( \sigma_t^2 = \alpha_0 + \alpha_1 Z_{t-1}^2 + \beta_1 \). Here \( \alpha_0, \alpha_1, \beta_1 \) are positive constants. The latter equation is of Kesten type \cite{Kesten1969} with \( A_t = \alpha_1 Z_{t-1}^2 + \beta_1 \) and \( B_t = \alpha_0 \). We assume that the conditions of Section 5.2 are satisfied, in particular,
\[
\mathbb{P}(\sigma > x) \sim c_0 x^{-\alpha}, \quad x \to \infty ,
\]
for some constant \( c_0 > 0 \) and tail index \( \alpha > 0 \), satisfying \( \mathbb{E}(\alpha_1 Z_0^2 + \beta_1)^{\alpha/2} = 1 \). We also assume that \( \mathbb{E}|Z|^{\alpha+\varepsilon} < \infty \) for some \( \varepsilon > 0 \). Rewriting \cite{Breiman1987}, we have
\[
\frac{\mathbb{P}(x^{-1}(\sigma_0, \ldots, \sigma_t) \in M)}{\mathbb{P}(\sigma > x)} \sim \frac{\mathbb{P}(x^{-1}\sigma_0(1, \Pi_1^{0.5}, \ldots, \Pi_t^{0.5}) \in M)}{\mathbb{P}(\sigma > x)} \to \int_0^\infty \alpha y^{-\alpha-1} \mathbb{P}(y(1, \Pi_1^{0.5}, \ldots, \Pi_t^{0.5}) \in M) dy .
\]
This shows that regular variation of \( (\sigma_t) \) with index \( \alpha \) follows from the regular variation of \( \sigma \). This property is inherited by the sequence \( (X_t) \). We observe that as \( x \to \infty \),
\[
\frac{\mathbb{P}(|(X_0, \ldots, X_t) - \sigma_0(Z_0, \Pi_0^{0.5} Z_1, \ldots, \Pi_t^{0.5} Z_t)| > x)}{\mathbb{P}(\sigma > x)} \leq \frac{\mathbb{P}(|Z_1|R_0^{0.5} + \cdots + |Z_t|R_t^{0.5} > x)}{\mathbb{P}(\sigma > x)} = o(1) .
\]
In the last step we used the independence of \( Z_t \) and \( R_t \) as well as the moment condition on \( Z \). Condition \( (RV_\alpha) \) for \( (X_t) \) now follows. This property was proved in Mikosch and Stărică \cite{Mikosch2004} under the additional condition that \( Z \) be symmetric. The above calculation shows that this assumption can be avoided.
Next we consider the 2-dimensional Markov chain
\[ \Phi_t = (\sigma_t, X_t)' = \sigma_t(1, Z_t)', \quad t \in \mathbb{Z}. \]
A similar calculation as above shows that this Markov chain satisfies \((\text{RV}_\alpha)\) and for \(h \geq 0\), any continuity set \(N\), observing that \(|\Phi_0| = \sigma_0 \sqrt{Z_0^2 + 1}\),
\[ \mathbb{P}(x^{-1}(\Phi_0, \ldots, \Phi_h) \in N \mid |\Phi_0| > x) \sim \mathbb{P}(x^{-1}\sigma_0((1, Z_0)', \Pi_1^{0.5}(1, Z_1)', \ldots, \Pi_h^{0.5}(1, Z_h)') \in N \mid |\Phi_0| > x) \]
\[ \xrightarrow{\text{a.s.}} \mathbb{P}(|Y_0|((1, Z_0)', \Pi_1^{0.5}(1, Z_1)', \ldots, \Pi_h^{0.5}(1, Z_h)')/(Z_0^2 + 1)^{0.5} \in N). \]
Identifying the limiting vector with \(|Y_0|(\Theta_0, \ldots, \Theta_h)'), we have for any \(\theta \in S,
\[ b(\theta) = \mathbb{E}\left[\left\{ \left( \theta'(1, Z_0)' + \sum_{l \geq 1} \Pi_1^{0.5}\theta'(1, Z_1)' \right)^{\alpha} - \left( \sum_{l \geq 1} \Pi_1^{0.5}\theta'(1, Z_1)' \right)^{\alpha} \right\} / (Z_0^2 + 1)^{\alpha/2}\right]. \]
The Markov chain \((\Phi_t)\) is aperiodic and irreducible under classical conditions on the density of the \(Z\); see e.g. \[ for details. The condition \((\text{DC}_p)\) for \(p < \alpha\) follows by an application of Lemma 2.1 for \(V(x) = |x|^p\). We recall that for \(m \geq 2\), \(\sigma_m^2 = \Pi_{2,m}(o_0 + \alpha_1 X_0 + \beta_1 \sigma_0^2) + \widetilde{R}_m\), where \(E|\widetilde{R}_m|^{(\alpha+c)/2} \leq \infty\) for some \(c > 0\) and \(\widetilde{R}_m\) is independent of \(Z_m\). We have for \(p < \alpha\), some \(c > 0,
\[ E[|\Phi_m|^p \mid \Phi_0 = y] = E\Pi_{2,m}(o_0 + \alpha_1 y_1 + \beta_1 y_2) + \widetilde{R}_m|^{p/2}E(Z^2 + 1)^{p/2} \leq |y|^p E\Pi_{2,m}|^{p/2}E(Z^2 + 1)^{p/2} \max(\alpha_1^p, \beta_1^p) + c. \]
For \(m = 1\), we find constants \(c_1, c_2 > 0\) such that \(E(V(\Phi_1)|\Phi_0 = y) \leq c_1 V(y) + c_2\). Since \(E A^{p/2} < 1\) for \(p < \alpha\), \((\text{DC}_{p,m})\) holds for sufficiently large \(m\) in view of \[ . An application of Lemma 2.1 concludes the proof. Thus we may apply the stable limit theory of Theorem 4.1 with \(f(x) = x\) to \((\Phi_t)\) for \(\alpha < 2\) and the limit law is determined by the cluster index \(b above.
For \(h \geq 0\) consider the Markov chain, recycling the notation \(\Phi_t,\)
\[ \Phi_t = (X_t, \sigma_t, \ldots, X_{t-h}, \sigma_{t-h}), \quad t \in \mathbb{Z}. \]
We also write
\[ \Phi_t^2 = (X_t^2, \sigma_t^2, \ldots, X_{t-h}^2, \sigma_{t-h}^2), \quad t \in \mathbb{Z}, \]
and introduce the function \(f\) acting on \((\Phi_t)\) given by
\[ Y_t = f(\Phi_t) = (X_t(X_{t-1}, \ldots, X_{t-h}), \Phi_t, \Phi_t), \quad t \in \mathbb{Z}. \]
We intend to show \((\text{DC}_{p,m})\) for \(p < \alpha/2\) and some large \(m\). We restrict ourselves to the case \(h = 1;\) the general case is analogous but requires more accounting. We observe that for suitable constants \(c > 0,\)
\[ |f(\Phi_t)|^p = |X_t^2 X_{t-1} + X_t^4 + X_{t-1}^4 + \sigma_t^4 + \sigma_{t-1}^4 + X_t^2 + X_{t-1}^2 + \sigma_t^2 + \sigma_{t-1}^2|^p / 2 \]
\[ \leq c((1 + Z_t^2)(1 + X_{t-1}^4 + \sigma_{t-1}^4) + Z_t^2(1 + \sigma_t^2 + X_{t-1}^2)(1 + X_{t-1}^2) + X_{t-1}^2 + 1 + \sigma_{t-1}^2)^{p/2} \]
Then for suitable constants \(c_1, c_2 > 0,\)
\[ E[|f(\Phi_1)|^p \mid \Phi_0 = y] \leq c(1 + |y_1|^p + |y_2|^p + |y_1|^p + |y_2|^p) \]
\[ \leq c_1|f(y)|^p + c_2. \]
By a similar argument, for sufficiently large $m \geq 1$, suitable constants $c > 0$, recalling that $\sigma^2_t = \Pi_t \sigma^2_0 + R_t$, where $\sigma^2_0$ is independent of $(\Pi_t, R_t)$, and

$$
\mathbb{E}[|f'(\Phi_m)|^p | \Phi_0 = y] \leq c(1 + \mathbb{E}[|\sigma^2_{m-1}|^{p/2} + |\sigma^2_{m-1}|^{p/2} | \Phi_0 = y]) \\
\leq c(1 + \mathbb{E}\Pi_{m}^2 |y|^{2p} + \mathbb{E}\Pi_{m}^2 |y|^p) \\
\leq c(\mathbb{E}[\Pi_{m}^{2p}] + \mathbb{E}[\Pi_{m}^p]) |f(y)|^p + c \\
\leq \beta|f(y)|^p + c,
$$

for some $\beta \in (0, 1)$, sufficiently large $m \geq 1$. Here we used the fact that $\mathbb{E}A^{2p} < 1$ for $p < \alpha/2$. Now we can apply Lemma 2.11 to show $(DC_{p})$ for $p < \alpha/2$.

It remains to show $(RV_{\alpha/2})$ for $(Y_t)$ defined in (5.9). The $\Pi_t$-component of $X_t$ is regularly varying with index $\alpha$. Therefore, without loss of generality and abusing notation, we will consider the sequence

$$
Y_t = f(\Phi_t) = (X_t(X_{t-1}, \ldots, X_{t-h}), \Phi^2_t), \quad t \in \mathbb{Z}.
$$

Similar arguments as in the first part of this subsection and as in Section 5.2 show for $t \geq 0$ that

$$
Y_t = R^{(1)}_t + \\
\sigma^2_t(Z_t \Pi^5_{t-h+1,t-1}(Z_t \Pi^5_{t-h+1,t-1}, \ldots, Z_t \Pi^5_{t-h})), (\Pi^5_{t-h+1,t}, (Z^2_t, 1), \ldots, (Z^2_{t-h}, 1))',
$$

where $\mathbb{E}|R^{(1)}_t|^p < \infty$, $i = 1, 2$ and

$$
(Y_0, \ldots, Y_t)' = D_t \sigma^2_{-h} + \tilde{Q}_t,
$$

where $\mathbb{E}|\tilde{Q}_t|^p < \infty$ and $\mathbb{E}|D_t|^p < \infty$ for some $\varepsilon > 0$ and

$$
D_t = \\
\begin{pmatrix}
Z_0 Z_{t-1} \Pi^5_{t-h,0} \Pi^5_{t-h-1,0} & Z_1 Z_0 A_{t-1} \Pi^5_{t-h,0} \Pi^5_{t-h-1,0} & \cdots & Z_t Z_{t-1} \Pi^5_{t-h,0} \Pi^5_{t-h-1,0} \\
Z_0 Z_{t-1} \Pi^5_{t-h,0} \Pi^5_{t-h-1,0} & Z_1 Z_0 A_{t-1} \Pi^5_{t-h,0} \Pi^5_{t-h-1,0} & \cdots & Z_t Z_{t-1} \Pi^5_{t-h,0} \Pi^5_{t-h-1,0} \\
\cdots & \cdots & \cdots & \cdots \\
Z_0 Z_{t-h} \Pi^5_{t-h,0} & Z_1 Z_{t-h} A_{t-h} \Pi^5_{t-h,0} & \cdots & Z_t Z_{t-h} A_{t-h} \Pi^5_{t-h,0} \\
\Pi_1 \Pi_0 (Z^2_{t-1}, 1) & \Pi_1 \Pi_0 (Z^2_{t-1}, 1) & \cdots & \Pi_1 \Pi_0 (Z^2_{t-1}, 1) \\
\Pi_1 \Pi_0 (Z^2_{t-1}, 1) & \Pi_1 \Pi_0 (Z^2_{t-1}, 1) & \cdots & \Pi_1 \Pi_0 (Z^2_{t-1}, 1) \\
(\Pi^2_{t-h,0}) & (\Pi^2_{t-h,0}) & \cdots & (\Pi^2_{t-h,0}) \\
(\Pi^2_{t-h,0}) & (\Pi^2_{t-h,0}) & \cdots & (\Pi^2_{t-h,0})
\end{pmatrix}.
$$

Notice that $\sigma^2_{-h}$ and $D_t$ are independent and that $\sigma^2_{-h}$ is regularly varying with index $\alpha/2$. Then $(RV_{\alpha/2})$ for $(Y_0, \ldots, Y_t)$ follows by an application of the multivariate Breiman result; see [3]. We omit the calculation of the cluster index; it is similar to its calculation in Section 5.2.

Now we can apply Theorem 4.4 to prove limit theory with $\alpha/2$-stable limits, $\alpha < 4$, for the sample autocovariance function of the GARCH(1, 1) process. The corresponding theory using point process techniques is given in [17, 18]. There the limit theory for the sequences $(|X_t|)$ and $(X^2_t)$ was also provided. The same results can be provided by Theorem 4.1 by calculating the corresponding cluster indices. Applied to the squares $(X^2_t)$ we obtain in particular for $\alpha \in (2, 4)$,

$$
\sqrt{n} a_n^{-1} \frac{1}{n} \sum_{t=1}^{n-h} X^2_t X^2_{t+h} - \left( \frac{1}{n} \sum_{t=1}^{n} X^2_t \right)^2 \overset{d}{\to} \xi_{\alpha/4},
$$

where $a_n$ is the appropriate normalising constant.
where \( \xi_{\alpha/4} \) is an \( \alpha/4 \)-stable random variable whose characteristic function is given in Theorem 4.1 and \((\Theta_t)_{t \geq 0} = (e^{Z^2_t Z^2_{t+h}} \Pi_{t+h}, \Theta_t)_{t \geq 0}\) for some \( c > 0 \). In particular, the \( \Theta_t \)'s are non negative and thus \( b_- = 0 \). Then \( \xi_{\alpha/4} \) is supported on \([-\infty, \infty)\). We omit further details. Relation 5.10 supports the idea of spurious long-range dependence effects observed on real-life log-return data which are often observed to have infinite fourth moments; see 33 for a discussion.

6. Proof of Theorem 4.1

6.1. Proof of part (1). We will use the Cramér-Wold device to show that \( (\alpha_n^{-1}, \Theta_n) \) has a (possibly degenerate) \( \alpha \)-stable limit \( \xi_n(\theta) \) for every \( \theta \in \mathbb{S}^{d-1} \). We will apply Theorem 1 in [2] which we recall for convenience:

**Theorem 6.1.** Assume that \((G_t)\) is a strictly stationary process of random variables, satisfying the following conditions.

1. The regular variation condition \((RV_\alpha)\) holds for some \( \alpha \in (0, 2) \).
2. The mixing condition \((MX)\): There exist \( m = m_n \to \infty \) such that \( k_n = [n/m_n] \to \infty \) and
   \[
   \mathbb{E} e^{itb_n^{-1}S_n(G)} - \left( \mathbb{E} e^{itb_n^{-1}S_m(G)} \right)^{k_n} \to 0, \quad n \to \infty, \quad t \in \mathbb{R},
   \]
   where \( S_n(G) = G_1 + \cdots + G_n \) and \( (b_n) \) is chosen such that \( n \mathbb{P}( |G_1| > b_n) \sim 1 \).
3. The anti-clustering condition
   \[
   \lim_{n \to \infty} \limsup_{m \to \infty} \frac{1}{m} \sum_{j=\ell+1}^m \mathbb{E} \left[ x b_n^{-1} (S_j(G) - S_\ell(G)) \right] = 0, \quad x \in \mathbb{R},
   \]
   holds, where \( m = m_n \) is the same as in \((MX)\) and \( \mathbb{P} = (x \wedge 2) \lor (-2) \).
4. The limits
   \[
   \lim_{\ell \to \infty} (b_+(\ell) - b_+ (\ell - 1)) = c_+ \quad \text{and} \quad \lim_{\ell \to \infty} (b_- (\ell) - b_- (\ell - 1)) = c_-,
   \]
   exist. Here \( b_+(\ell), b_-(\ell) \) are the tail balance parameters given by \( b_+(\ell) = \lim_{m \to \infty} n \mathbb{P}(S_\ell(G) > b_n) \) and \( b_-(\ell) = \lim_{m \to \infty} n \mathbb{P}(S_\ell(G) < -b_n) \).
5. For \( \alpha > 1 \) assume \( \mathbb{E} G_1 = 0 \) and for \( \alpha = 1 \),
   \[
   \lim_{\ell \to \infty} \limsup_{n \to \infty} \mathbb{E} |\sin(b_n^{-1} S_\ell(G))| = 0.
   \]

Then \( c_+ \) and \( c_- \) are non-negative and \( (b_n^{-1} S_n(G)) \) converges in distribution to an \( \alpha \)-stable random variable (possibly zero) with characteristic function \( \psi_{\alpha}(x) = \exp(-|x|^\alpha \chi_{\alpha}(x, c_+, c_-)) \), where for \( \alpha \neq 1 \) the function \( \chi_{\alpha}(x, c_+, c_-), x \in \mathbb{R} \), is given by the formula

\[
\frac{\Gamma(2-\alpha)}{1-\alpha} \left( (c_+ + c_-) \cos(\pi \alpha/2) - i \text{sign}(x) (c_+ - c_-) \sin(\pi \alpha/2) \right),
\]

while for \( \alpha = 1 \) one has

\[
\chi_1(x, c_+, c_-) = 0.5 \pi (c_+ + c_-) + i \text{sign}(x) (c_+ - c_-) \log |x|, \quad x \in \mathbb{R}.
\]

We will verify the conditions of this theorem for the sequence \( G_t = \theta X_t \) for fixed \( \theta \in \mathbb{S}^{d-1} \).

*The regular variation condition \((RV_\alpha)\) for \((G_t)\).* This condition is straightforward from the definition of \((RV_\alpha)\) for \((X_t)\) and the fact that the function \( f(x) = \theta \cdot x, x \in \mathbb{R}^d \), is continuous and homogeneous.
The anti-clustering condition (AC). Without loss of generality we assume that (DC_p) holds for V(y) = |y|^p. We also assume p \leq 1; for p > 1 an application of Jensen’s inequality yields (DC_{p'}) for any p' < p. Since p \leq 1 there exists c > 0 such that y \leq cy^p for y \in [0, 2]. Then one has

\[ T_{\ell m} = \frac{n}{m} \sum_{j=\ell+1}^{m} E \left[ |x b_n^{-1} (S_j(G) - S_\ell(G))| x b_n^{-1} G_1 \right] \]
\[ \leq \frac{c}{m} \sum_{j=\ell+1}^{m} E \left[ |x b_n^{-1} (S_j(G) - S_\ell(G))|^p x b_n^{-1} G_1 \right]. \]

Using (DC_p), a recursive argument yields

\[(6.1) \quad E(|G_k|^p | \Phi_1 = y) \leq \beta^{k-1} |f(y)|^p + b \sum_{j=1}^{k-1} \beta^j, \quad k \geq 2, \]

where \(\beta, b\) appear in (DC_p). Multiple use of this argument and the subadditivity of the function \(z \mapsto z^p\) on \((0, \infty)\) yield for \(\ell < j \leq m\),

\[ E \left[ |x b_n^{-1} (S_j(G) - S_\ell(G))|^p | \Phi_1 \right] \leq c \frac{n}{m} \sum_{i=\ell+1}^{m} \beta^i |X_1|^p + c b_n^{-p} m. \]

Conditioning on \(\Phi_1\), the latter inequality finally yields

\[ \mathbb{E} T_{\ell m} \leq \frac{c}{m} \sum_{j=\ell+1}^{m} E \left[ |x| |b_n^{-p} \sum_{i=1}^{j} \beta^i |X_1|^p x b_n^{-1} |X_1| \right] + \frac{c}{m} b_n^{-p} E |x b_n^{-1} X_1| = I_1 + I_2. \]

We have \(I_2 \leq c b_n^{-p-1} m = o(1)\) if we choose \(m = m_n = \log^2 n\). It remains to prove that \(I_1\) is asymptotically negligible. An application of Karamata’s theorem yields the bound

\[ I_1 \leq c \frac{n}{m} \sum_{j=\ell+1}^{m} \mathbb{P} \left( |X_1| \geq c b_n \sum_{i=\ell}^{j} \beta^i \right)^{-1/(p+1)} \leq c \frac{m}{m} \sum_{j=\ell+1}^{m} \left( \sum_{i=\ell}^{j} \beta^i \right)^{-1/(p+1)} \leq c \beta^{\ell}/(p+1). \]

The right-hand side vanishes as \(\ell \to \infty\). Collecting the above bounds, condition (AC) follows.

The mixing condition (MX). Here we give a significant improvement on Lemma 3 in [2]; in the latter paper it is assumed that \((G_t)\) is strongly mixing. The next result avoids this condition.

**Lemma 6.2.** Consider a strictly stationary real-valued sequence \((G_t)\) satisfying the conditions (RV_\(\alpha\)) for some \(\alpha \in (0, 2)\) and (AC). Then (MX) can be replaced by Condition (MX'): There exists a sequence \((r_n)\) such that \(r_n = o(m_n)\) and

\[ |\varphi_n^{(\ell)}(t) - \varphi_{n, m-\ell}(t)| \to 0, \quad t \in \mathbb{R}. \]

holds for \(\ell = m - r_n\) and \(\ell = r_n\), where

\[ \varphi_n^{(\ell)}(t) = E \left[ \exp \left( it b_n^{-1} \sum_{i=1}^{k_n} \sum_{t=(i-1)m+1}^{|m-\ell|} G_t \right) \right], \]
\[ \varphi_{n,j} = E \left[ \exp \left( it b_n^{-1} \sum_{i=1}^{j} G_t \right) \right], \quad j \geq 1, \varphi_n(t) = \varphi_{n,n}(t), \quad t \in \mathbb{R}. \]
Proof. Notice that condition (MX) can be written in the form \( \varphi_n(t) - \phi^k_{n,m}(t) \to 0 \) as \( n \to \infty \). We have

\[
\varphi_n(t) - \phi^k_{n,m}(t) = [\varphi_n(t) - \phi^{(r)}(t)] + [\phi^{(r)}(t) - \phi^k_{n,m-r}(t)] + [\phi^k_{n,m-r}(t) - \phi^k_{n,m}(t)]
\]

\[= P_1 + P_2 + P_3.\]

In view of (MX)', \( P_2 \to 0 \). Next we deal with \( P_1 \). Assume for simplicity that \( k_n = n/m \) is an integer. We use the classical Bernstein blocks technique, writing

\[
S_n = b_n^{-1} \sum_{i=1}^{k_n} \sum_{t=(i-1)m+1}^{im-r} G_t + b_n^{-1} \sum_{i=1}^{k_n} \sum_{t=im-r+1}^{im} G_t = I_1 + I_2.
\]

We will show that \( \mathbb{E}\exp(itI_2) \to 1 \). Condition (MX)' implies that \( |\mathbb{E}\exp(itI_2) - \varphi^k_{n,r}(t)| \to 0 \) as \( \ell = m - r \geq r \) and \( \ell/n \to 0 \). Moreover, Lemma 3.5 in [14] yields that \( \varphi^k_{n,r}(t) \to 1 \) if and only if \( k(\varphi_{n,r}(t) - 1) \to 0 \). Assuming (RV\( \alpha \)) and (AC), one can follow the proof of Lemma 1 in [2]. We have

\[
\lim_{q \to \infty} \lim_{n \to \infty} \sup \{k(\varphi_{n,r}(t) - 1) - kr(\varphi_{n,q}(t) - \varphi_{n,q-1}(t))\} \to 0, \quad t \in \mathbb{R}.
\]

Under (RV\( \alpha \)), an application of Theorem 3 in Section XVII.5 of Feller gives that \( n(\varphi_{n,q}(t) - 1) \) converges for all \( q \). We deduce that \( n(\varphi_{n,q}(t) - 1) \) converges too. As \( kr/n \to r/m \to 0 \) we conclude that \( kr(\varphi_{n,q}(t) - \varphi_{n,q-1}(t)) \to 0 \) and then \( k_n(\varphi_{n,r}(t) - 1) \to 0 \) which gives the desired result \( \mathbb{E}\exp(itI_2) \to 1 \), equivalently, \( I_2 \xrightarrow{P} 0 \). Since

\[
|P_1| = \mathbb{E}\{\exp(itI_1)(1 - \exp(itI_2))\} \leq \mathbb{E}1 - \exp(itI_2),
\]

dominated convergence yields \( P_1 \to 0 \). Finally,

\[
|P_3| \leq k \left| \left( \varphi_{n,m-r}(t) - 1 \right) - \varphi_{n,m}(t) \right| \to 0.
\]

and the same arguments as above show that \( P_3 \to 0 \). \( \square \)

We finish the proof of (MX) for the sequence \( (G_\ell) \). In view of (DC\( p \)), \( (X_\ell) \), hence \( (G_\ell) \), are \( \beta \)-mixing, hence strongly mixing, with exponential rate \( (\alpha_\ell) \). We will show (MX) by an application of Lemma 6.22. A standard telescoping sum argument shows that

\[
|\varphi^{(t)}_{n}(t) - \phi^k_{n,m-\ell}(t)| \leq c k_n \alpha_\ell.
\]

Since we choose \( m = \log^2 n \) in the proof of (AC), \( k_n \alpha_\ell \leq (n/\log^2 n)\exp(-c \ell_n) \). Thus, choosing \( \ell_n = C \log n \) for some sufficiently large constant \( C > 0 \) we have \( \ell_n = o(m_n) \), \( k_n \alpha_\ell = o(1) \) and we can also find \( r_n = o(\ell_n) \). This proves (MX'), hence (MX).

**Condition (TB).** Note that \( \{|\theta'X| > b_n\} \subset \{|X| > b_n\} \). Then

\[
b_+(\ell) = \lim_{x \to \infty} \frac{\mathbb{P}(S_\ell(G) > x)}{\mathbb{P}(|\theta'X| > x)}
\]

\[
= \lim_{x \to \infty} \frac{\mathbb{P}(\theta'G \geq x)}{\mathbb{P}(|X| > x)} \lim_{x \to \infty} \frac{\mathbb{P}(|X| > x)}{\mathbb{P}(|\theta'X| > x)}
\]

\[
= b_+(\theta) \lim_{x \to \infty} \left( \mathbb{P}(|\theta'X| > x) / \mathbb{P}(|X| > x) \right)^{-1}
\]

\[
= b_+(\theta) \mathbb{P}(\{Y_0|\theta'\Theta_0| > 1\})^{-1}
\]

\[
= b_+(\theta) \mathbb{E}(\theta'\Theta_0)^{\alpha}^{-1}.
\]

Correspondingly, \( b_-(\ell) = b_-(\theta)(\mathbb{E}(\theta'\Theta_0)^{\alpha})^{-1} \). Here we assumed that \( \mathbb{E}(\theta'\Theta_0)^{\alpha} \neq 0 \). Otherwise, \( b_+(\ell) = b_-(\ell) = 0 \).
Thus we may apply Theorem 6.1 to conclude that \( b_n^{-1} \theta' S_n \xrightarrow{d} \xi_\alpha(\theta) \) for an \( \alpha \)-stable random variable \( \xi_\alpha(\theta) \) with characteristic function \( \psi_\alpha(x, \theta) \) given by

\[
\mathbb{E}(|\theta' \Theta_0|^\alpha) \log \psi_\alpha(x, \theta) = -|x|^\alpha \frac{\Gamma(2-\alpha)}{1-\alpha} \left( (b(\theta) + b(-\theta)) \cos(\pi \alpha/2) - i \text{sign}(x)(b(\theta) - b(-\theta)) \sin(\pi \alpha/2) \right), \quad x \in \mathbb{R}. 
\]

The factor \( \mathbb{E}(|\theta' \Theta_0|^\alpha) \) on the left-hand side is due to the normalization \( \langle b_n \rangle \) instead of \( \langle a_n \rangle \). Replacing \( \langle b_n \rangle \) by \( \langle a_n \rangle \), we have for any \( v \in \mathbb{R}^d \) that

\[
\mathbb{E}_v \exp \left\{ -|v|^{\alpha} C_\alpha^{-1} \left( (b(v/|v|) + b(-v/|v|)) - i (b(v/|v|) - b(-v/|v|)) \tan(\pi \alpha/2) \right) \right\},
\]

where \( C_\alpha \) is defined in (4.2). This is the characteristic function of an \( \alpha \)-stable random vector \( \xi_\alpha \). The representation of the Lévy spectral measure \( \Gamma_\alpha \) in (4.1) follows by calculations as in Example 2.3.4 of [50]. Indeed, keeping notations of [50] and identifying the limiting law yields the equations

\[
b(\theta) + b(-\theta) = C_\alpha \sigma_\theta^\alpha = C_\alpha \int_{\mathbb{R}^d} |\theta'|^\alpha \Gamma_\alpha(ds) = C_\alpha \int_{S^{d-1}} (\theta')^\alpha \Gamma_\alpha(ds) + C_\alpha \int_{S^{d-1}} (-\theta')^\alpha \Gamma_\alpha(ds),
\]

and

\[
b(\theta) - b(-\theta) = (b(\theta) + b(-\theta)) \beta_\theta = C_\alpha \int_{S^{d-1}} |\theta'|^\alpha \text{sign}(\theta') \Gamma_\alpha(ds)
\]

\[= C_\alpha \int_{S^{d-1}} (\theta')^\alpha \Gamma_\alpha(ds) - C_\alpha \int_{S^{d-1}} (-\theta')^\alpha \Gamma_\alpha(ds), \quad \theta \in S^{d-1}.
\]

The limiting \( \alpha \)-stable distribution is degenerate if and only if \( b(\theta) = 0 \) for all \( \theta \in S^{d-1} \). This proves part (1) of the theorem.

### Stable limit theory for general regularly varying stationary processes

In this part we want to give some arguments showing that the results of Theorem 4.1 can be applied in much more general context. For this reason, consider a strictly stationary \( \mathbb{R}^d \)-valued regularly varying sequence \( \{X_t, t \in Z\} \), \( t \in Z \), constitutes a Markov chain with infinite-dimensional state space. In this setting, \( \{\text{DC}_p\} \) for \( X_t = f(\Phi_t) \) takes on the form:

**Condition (DC\(_p\))**:

\[
\mathbb{E}(|X_1|^p \mid |X_0, X_{-1}, \ldots| = |x_0, x_{-1}, \ldots|) \leq \beta |x_0|^p + b \quad \text{for some } 0 < \beta < 1 \text{ and } b > 0.
\]

We also need a weak dependence assumption more general than geometric -mixing which, in the irreducible case, is implied by the drift condition.

**Condition (MX\(_m\))**: Consider an integer sequence \( m_n \) such that \( m = m_n \to \infty \) and \( m_n/n = O(1) \) and also write \( k_n = [n/m] \). There exists a sequence \( (r_n) \) such that \( r_n = o(m_n) \) and

\[
\lim_{n \to \infty} |\varphi_n^{(\ell)}(s) - \varphi_{n,m-\ell}(s)| = 0, \quad s \in \mathbb{R}^d,
\]
holds for both \( \ell = \ell_n = m_n - r_n \) and \( \ell = r_n \), where

\[
\varphi_n^{(\ell)}(s) = \mathbb{E} \left[ \exp \left( ia_n^{-1} \sum_{i=1}^{k_n} \sum_{t=(i-1)m+1}^{im-\ell} s' X_t \right) \right],
\]

\[
\varphi_{n,j}(s) = \mathbb{E} \left[ \exp \left( ia_n^{-1} \sum_{t=1}^{j} s' X_t \right) \right], \quad j \geq 1, \quad \varphi_n(s) = \varphi_{n,n}(s), \quad s \in \mathbb{R}^d.
\]
Condition \( (\text{MX}_m) \) is implied by \( \theta \)-weak dependence introduced by Doukhan and Louhichi \[13\]. For every \( m \geq 1 \), equip \( (\mathbb{R}^d)^m \) with the metric \( \| \cdot \|_m = m^{-1} \sum_{i=1}^m \| \cdot \| \). A function \( f : (\mathbb{R}^d)^m \to [-1, 1], m \geq 1 \), is Lipschitz if
\[
\sup_{x \neq y} \frac{|f(x) - f(y)|}{|x - y|_m} = \text{Lip}(f) < \infty.
\]
The \( \theta \)-weak dependence coefficients \( (\theta_r)_{r \geq 0} \) are defined for any \( f \) with \( \text{Lip}(f) = 1 \) and measurable \( g : (\mathbb{R}^d)^v \to [-1, 1], v \geq 1 \), as
\[
\sup_{k, v \geq 1, i_1 < \cdots < i_k \leq 0 \leq j_1 < \cdots < j_v} |\text{cov}(g(X_{i_1}, \ldots, X_{i_k}), f(X_{j_1}, \ldots, X_{j_m}))| = \theta_r.
\]
Condition \( (\text{MX}_m) \) follows if \( \theta_r \to 0 \) for some \( r = r_n = o(m) \) with \( m = m_n \). \( \theta \)-weak dependence covers a wide range of known dependence concepts, including a large variety of mixing conditions; see \[18\].

In the general case, the following analog of Theorem 4.1 holds. The proof follows along the lines of Theorem 4.1. Irreducibility of \( (X_t) \) can be replaced by \( (\text{MX}_m) \). We omit further details.

**Theorem 6.3.** Consider an \( \mathbb{R}^d \)-valued strictly stationary sequence \( (X_t) \) satisfying the following conditions:

- \( (\text{RV}_\alpha) \) for some \( \alpha \in (0, 2) \), \( \mathbb{E}X = 0 \) if \( \alpha > 1 \) and \( X \) is symmetric if \( \alpha = 1 \).
- \( (\text{DC}_p') \) for some \( p \in ((\alpha - 1) \vee 0, \alpha) \).
- \( (\text{MX}_m) \) for \( m_n = o(n^{(p+1)/\alpha - 1}) \).

Let \( (a_n) \) be a sequence of positive numbers such that \( n \mathbb{P}(|X_0| > a_n) \sim 1 \). Then the statement of part (1) of Theorem 4.1 holds.

**6.2. Proof of part (2).** Recall the regenerative structure of the Markov chain \( (X_t) \) from Section 2.2.
We will show that the partial sum \( S(1) \) over a full regenerative cycle is regularly varying with index \( \alpha \). We write
\[
S_n = S(0) + \sum_{t=1}^{N_A(n)} S(t) + \sum_{i=N_A(n)+1}^n X_i,
\]
where \( N_A(n) = \#\{i \geq 0 : \tau_A(i) \leq n\} \), \( n \geq 1 \), is independent of \( (S(i))_{i \geq 1} \). The first and last block sums \( S(0) \) and \( \sum_{i=\tau_A(N_A(n))}^n X_i \) are tight. Therefore
\[
a_n^{-1}S_n = a_n^{-1} \sum_{t=1}^{N_A(n)} S(t) + o_P(1).
\]
By virtue of \( (\text{DC}_p) \) for some \( p > 0 \) the chain \( (X_t) \) is geometrically ergodic. Therefore there exists a constant \( \kappa > 0 \) such that
\[
\sup_{x \in A} \mathbb{E}_x e^{\kappa \tau_A} < \infty,
\]
In the last step we used a maximal inequality of Ottaviani type; see e.g. [44], Chapter 2. The second term on the right-hand side is negligible, as first letting $\beta \to 0$ argument, $N$ Using the Wald identity, we have

$$\sum_{t=1}^{N_A(n)} S(t) = o_P(1).$$

In view of part (1), the sum of iid random vectors $(S(i))$ on the right-hand side has an $\alpha$-stable limit. It follows from [40] that $S(1)$ is regularly varying with index $\alpha$. This concludes the proof.

7. Proof of Theorem 4.3

7.1. Proof of part (1): The case $\alpha \in (0, 2)$. Recall the decomposition [62] of the partial sums $S_n$ in terms of the regenerative cycles of the Markov chain. We start with an auxiliary result which deals with the sums over the first and last blocks.

Lemma 7.1. Assume that (RV$\alpha$) and (DC$\beta$) hold for some $p > \alpha - 1$ provided $\alpha > 1$. Then there exists a constant $c > 0$ such that for any sequence $x = x_n \to \infty$ as $n \to \infty$,

$$\mathbb{P}(\sum_{t=1}^{N_A} |X_t| > x) \leq c \mathbb{P}(|X| > x),$$

$$\mathbb{P}\left(\sum_{t=1}^{N_A} |X_t| > x, \tau_A \leq n\right) = o(n \mathbb{P}(|X| > x)).$$

Proof. We start by proving (7.1). For any random vector $X$ we write $\overline{X} = X \mathbb{1}_{(|X| \leq x)}$. Then

$$\mathbb{P}_{\tau_A} \left(\sum_{t=1}^{\tau_A} |X_t| > x\right) \leq \mathbb{P}_{\tau_A} \left(\sum_{t=1}^{\tau_A} |\overline{X}_t| > x/2\right) + \mathbb{P}_{\tau_A} \left(\bigcup_{t=1}^{\tau_A} \{|X_t| \neq |\overline{X}_t|\}\right) = I_1 + I_2.$$

Using the Wald identity, we have

$$I_2 = \mathbb{E}_{\tau_A} \left(\max_{1 \leq t \leq \tau_A} 1_{\{|X_t| > x\}\}}\right) \leq \mathbb{E}_{\tau_A} \left(\sum_{t=1}^{\tau_A} 1_{\{|X_t| > x\}}\right) = \mathbb{E}_{\tau_A} \mathbb{P}(|X| > x).$$

Write $k_0 = \lceil \alpha \rceil$ and choose $0 < \beta < 1$ such that $\beta k_0 > \alpha$. Since $k_0 (k_0 - 1) \geq \alpha (\alpha - 1)$ for $\alpha > 1$, we will choose $\beta$ such that $p/\beta = k_0 - 1$. Markov’s inequality yields

$$I_1 \leq c \mathbb{E}_{\tau_A} \left(\sum_{t=1}^{\tau_A} |X_t|^{\beta k_0}\right)^{\beta k_0} \leq c \mathbb{E}_{\tau_A} \left(\sum_{t=1}^{\tau_A} |X_t|^{\beta k_0}\right)^{k_0}.$$
Note that \( (|X|^\beta) \) satisfies \((DC_{k_0-1})\). Under the latter condition we may apply Proposition 4.7 of [40] to get \( E_A(\sum_{t=1}^{\tau_A} |X_t|^\beta)^{k_0} \leq c |\Lambda|^{\beta k_0} \). An application of Karamata’s theorem shows that the right-hand side is bounded by \( cP(|X| > x) \). This concludes the proof of (7.2).

Now we turn to the proof of (7.2). Abusing notation, we write \( X = \mathbb{1}_{(|X| \leq x\delta)} \) for any fixed \( \delta \). Then
\[
\mathbb{P}\left( \sum_{t=1}^{\tau_A} |X_t| > x, \tau_A \leq n \right) \leq \mathbb{P}\left( \sum_{t=1}^{\tau_A} |X_t| > x, \tau_A \leq n \right) + \mathbb{P}\left( \bigcup_{t=1}^{\tau_A} \{|X_t| \not= |X_t|\} \right) = I_1 + I_2.
\]

Since \( E\tau_A < \infty \) and \( X \) is regularly varying, we have
\[
I_2 \leq E(\tau_A) \mathbb{P}(|X| > x\delta) = o(n \mathbb{P}(|X| > x)).
\]

Similar arguments as above yield
\[
I_1 \leq c \mathbb{E}(\sum_{t=1}^{\tau_A} |X_t|^\beta \mathbb{1}_{\{\tau_A \leq n\}})^{k_0} \leq c \mathbb{E}(\sum_{t=1}^{\tau_A} |X_t|^\beta \mathbb{1}_{\{\tau_A \geq t\}})^{k_0}.
\]

An argument similar to the one used in the proof of Theorem 4.6 in [40] shows that
\[
\mathbb{E}\left( \sum_{t=1}^{\tau_A} |X_t|^\beta \mathbb{1}_{\{\tau_A \geq t\}} \right)^{k_0} \leq c \mathbb{E}\left( \sum_{t=1}^{\tau_A} |X_t|^\beta \mathbb{1}_{\{\tau_A \geq t\}} \right).
\]

Finally, an application of Pitman’s identity [15], Proposition 4.7 in [40] and Karamata’s theorem yield
\[
\mathbb{E}\left( \sum_{t=1}^{\tau_A} |X_t|^\beta \mathbb{1}_{\{\tau_A \geq t\}} \right) = \mathbb{P}(X_0 \in A) E_A\left( \sum_{k=0}^{\tau_A} \sum_{t=1}^{n} |X_{k+t}|^\beta \mathbb{1}_{\{\tau_A \geq k+t\}} \right) \\
\leq n \mathbb{P}(X_0 \in A) E_A\left( \sum_{t=1}^{\tau_A} |X_t|^\beta \right) \\
\leq c n \mathbb{E}|X|^\beta \sim c n x^\beta \delta^{\beta k_0} \mathbb{P}(|X| > x).
\]

Since \( \beta k_0 > 0 \) and we can make \( \delta \) small as we like, we conclude that \( I_1 = o(n \mathbb{P}(|X| > x)) \). This concludes the proof of (7.2).

Now we are ready to prove part (1). Since \( \tau_A \) has exponential moment, it follows that \( \mathbb{P}(\tau_A > n) = o(\mathbb{P}(|X| > \lambda_n)) \). Therefore we may prove the result on the event \( \{\tau_A \leq n\} \). We write for simplicity \( \mathbb{P}_n(\cdot) = \mathbb{P}(\cdot \cap \{\tau_A \leq n\}) \). In view of Lemma [7.1] and the decomposition [6.2] of \( S_n \) we may neglect the sums over the first and last cycles and it suffices to prove the large deviation principle for the process \( \sum_{t=1}^{N_A(n)} S(t) \) over independent cycles. Observe that
\[
\frac{\mathbb{P}_n(\lambda_n^{-1} \sum_{t=1}^{N_A(n)} S(t) \in \cdot)}{n \mathbb{P}(|X| \geq \lambda_n)} = \frac{\mathbb{P}_n(\lambda_n^{-1} \sum_{t=1}^{N_A(n)} S(t) \in \cdot)}{n \mathbb{P}(|S(1)| \geq \lambda_n)} \frac{\mathbb{P}(|S(1)| > \lambda_n)}{\mathbb{P}(|X| > \lambda_n)}
\]

The same arguments as in the proof of Lemma 4.12 in [40] (here the conditions \( \lambda_n \to \infty \) and \( \lambda_n/n^{\delta+1/\alpha} \to 0 \) for some \( \delta > 0 \) are crucial) show that for any small \( \xi, \varepsilon > 0 \), and any set \( B \) bounded away from zero,
\[
(1 - \varepsilon) \mathbb{P}(\lambda_n^{-1}(1 + \xi)^{-1}(1 + \varepsilon)^{-1} S(1) \in B) \leq \frac{\mathbb{P}_n(\lambda_n^{-1} \sum_{t=1}^{N_A(n)} S(t) \in B)}{n \mathbb{P}(|S(1)| \geq \lambda_n)} + o(1)
\]
\[
\leq \frac{\mathbb{P}_n(\lambda_n^{-1}(1 - \xi)^{-1} S(1) \in B)}{\mathbb{P}(\tau_A) \mathbb{P}(|S(1)| > \lambda_n)} + o(1).
\]
Assume first that the cluster index \( b \) does not vanish. In view of part (2) of Theorem 4.1, we know that \( S(1) \) is regularly varying with index \( \alpha \) and spectral measure \( \mathbb{P} \) given by (4.3), and we also know that
\[
\frac{\mathbb{P}(\lambda_n^{-1} S(1) \in \cdot)}{\mathbb{P}(|S(1)| > \lambda_n)} \overset{\nu}{\to} \mu_{S(1)}(\cdot),
\]
for a non-null Radon measure \( \mu_{S(1)} \). Hence, letting \( \varepsilon \to 0 \) and \( \xi \to 0 \), we conclude that
\[
\frac{\mathbb{P}_n\left(\lambda_n^{-1} \sum_{t=1}^{N_A(n)} S(t) \in \cdot \right)}{n\mathbb{P}(|S(1)| > \lambda_n)} \overset{\nu}{\to} \mu_{S(1)}(\cdot) \cdot \mathbb{E}_A.
\]
It remains to determine the limit of \( \mathbb{P}(|S(1)| > x)/\mathbb{P}(|X| > x) \) as \( x \to \infty \). By virtue of the proof of Theorem 4.1, \( a_n^{-1} \sum_{t=1}^{n/\mathbb{E}_A} S(t) \xrightarrow{d} \xi_\alpha \). Then necessarily
\[
\frac{n}{\mathbb{E}_A} \mathbb{P}(a_n^{-1} S(1) \in \cdot) \overset{\nu}{\to} \nu_\alpha(\cdot),
\]
where \( \nu_\alpha \) is the Lévy measure of \( \xi_\alpha \). Hence
\[
\frac{\mathbb{P}(|S(1)| > a_n)}{\mathbb{P}(|X| > a_n)} \sim \frac{n\mathbb{P}(|S(1)| \geq a_n)}{\mathbb{E}_A} \to \mathbb{E}_A \Gamma_\alpha(S^{d-1}),
\]
where \( \Gamma_\alpha \) is the spectral measure of \( \nu_\alpha \). But from part (2) of Theorem 4.1 we know that \( n\mathbb{P}(|S(1)| \geq a_n) \to \mathbb{E}_A \int_{S^{d-1}} b(\theta) \, d\Gamma_\alpha(\theta) \). This proves the result in the non-generate case \( b \neq 0 \).

In the degenerate case \( b = 0 \), \( \mathbb{P}(|S(1)| \geq x) = o(\mathbb{P}(|X| \geq x)) \) as \( x \to \infty \). By independence of the cycles and since \( \lambda_n/a_n \to \infty \), for any set \( B \) bounded away from zero, some \( \gamma > 0 \),
\[
\mathbb{P}_n\left(\lambda_n^{-1} \sum_{t=1}^{N_A(n)} S(t) \in B \right) \leq \mathbb{P}\left(\left\| \sum_{t=1}^{N_A(n)} S(t) \right\| > \gamma \lambda_n \right) \leq nc \mathbb{P}(|S(1)| \geq \gamma \lambda_n) = o(n \mathbb{P}(|X| > a_n)) = o(1).
\]
The desired result in the degenerate case follows.

7.2. Proof of part (2): The case \( \alpha > 2 \). We only consider the non-degenerate case \( b \neq 0 \). We will apply Theorem 4.6 in 411 for functions of Markov chains in the case \( d = 1 \).

**Theorem 7.2.** Let \( (G_t) = (f(\Phi_t)) \) be a 1-dimensional functional of a strictly stationary \( \mathbb{R} \)-valued irreducible aperiodic Markov chain \( (\Phi_t) \). Write \( S_n(G) = G_1 + \cdots + G_n, n \geq 1 \), for the corresponding random walk. Assume that the following conditions are satisfied.

1. The regular variation condition \((RV_\alpha)\) for some \( \alpha > 2 \) and \( \mathbb{E}G = 0 \).
2. The anti-clustering condition \((AC_\alpha)\):
\[
\lim_{k \to \infty} \limsup_{n \to \infty} \sup_{x \in \Lambda_n} \delta_k^{-\alpha} \sum_{j=k}^{n} \mathbb{P}(|G_j| > x \delta_k) = 0.
\]
for a sequence \( \delta_k = o(k^{-2}) \), \( k \to \infty \), and sets \( \Lambda_n \) such that \( b_n = \inf \Lambda_n \to \infty \) as \( n \to \infty \).
3. The limit \( b_+ = \lim_{k \to \infty} (b_+(k+1) - b_+(k)) \) exists, where the constants \( (b_+(k)) \) are defined in Theorem 4.7.
4. The drift condition \((DC_p)\) for every \( p < \alpha \).

Then the precise large deviation principle
\[
\lim_{n \to \infty} \sup_{x \in \Lambda_n} \left| \frac{\mathbb{P}(S_n(G) > x)}{n \mathbb{P}(|G| > x)} - b_+ \right| = 0,
\]
Lemma A.1. Assume that \( (X_i) \) is an \( \mathbb{R}^d \)-valued strictly stationary sequence which is regularly varying with index \( \alpha > 0 \) and satisfies the one-dimensional large deviation principle

\[
\frac{\mathbb{P}(\theta S_n > \lambda_n)}{\mathbb{P}(|X| > \lambda_n)} \to b(\theta), \quad \theta \in \mathbb{S}^{d-1},
\]

for some sequence \( \lambda_n \to \infty \) such that \( n\mathbb{P}(|X| > \lambda_n) \to 0 \). Moreover, assume that \( \alpha \not\in \mathbb{N} \) or \( b(\cdot) = b(-\cdot) \). Then \( (\ref{A1}) \) holds.

Proof. Define the measures

\[
m_n(\cdot) = \frac{\mathbb{P}(\lambda_n^{-1}S_n \in \cdot)}{\mathbb{P}(|X| > \lambda_n)}, \quad n \geq 1,
\]

on the Borel \( \sigma \)-field of \( \mathbb{R}^d \). We conclude from \( (\ref{A1}) \) that for any Borel set \( B \) bounded away from zero,

\[
\sup_{n \geq 1} m_n(B) < \infty.
\]
This means that \((m_n(B))\) is vaguely tight; see \([31, 46]\). In view of \((1.1)\), any vague subsequential limit \(\mu\) of \((m_n)\) satisfies the relation \((4.2)\). For non-integer \(\alpha\), the latter property combined with the proof of Theorem 1.1 in \([3]\) shows that all vague subsequential limits of \((m_n)\) are identical and uniquely determined by the property \((4.2)\). Hence \((1.1)\) holds and the limit \(\nu_\alpha\) is given by \((4.2)\). A careful study of the proof of Theorem 1.1 in \([3]\) shows that the proof remains valid if the subsequential limits have the property \(\mu(\cdot) = \mu(-\cdot)\) which follows if \(b(\cdot) = b(-\cdot)\). □

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