

Regularly varying measures on metric spaces: Hidden regular variation and hidden jumps

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Abstract: We develop a framework for regularly varying measures on complete separable metric spaces \mathbb{S} with a closed cone C removed, extending material in [15, 24]. Our framework provides a flexible way to consider hidden regular variation and allows simultaneous regular-variation properties to exist at different scales and provides potential for more accurate estimation of probabilities of risk regions. We apply our framework to iid random variables in \mathbb{R}_+^∞ with marginal distributions having regularly varying tails and to càdlàg Lévy processes whose Lévy measures have regularly varying tails. In both cases, an infinite number of regular-variation properties coexist distinguished by different scaling functions and state spaces.

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1. Introduction

This paper discusses a framework for regular variation and heavy tails for distributions of metric-space-valued random elements and applies this framework to regular variation for measures on \mathbb{R}_+^∞ and $\mathbb{D}([0, 1], \mathbb{R})$. (Here and elsewhere we use the notation \mathbb{R}_+ for the space $[0, \infty)$ and $\mathbb{D}([0, 1], \mathbb{R})$ for the space of all real-valued, right-continuous functions with left limits on $[0, 1]$.)

Heavy tails appear in diverse contexts such as risk management; quantitative finance and economics; complex networks of data and telecommunication transmissions; as well as the rapidly expanding field of social networks. Heavy tails are also colloquially called power law tails or Pareto tails, especially in one dimension. The mathematical formalism for discussing heavy tails is the theory of regular variation, originally formulated on \mathbb{R}_+ and extended to more general spaces. See, for instance, [8, 16, 17, 19, 20, 24, 26, 37, 40, 42, 46].

One approach to estimating the probability of a remote risk region relies on asymptotic analysis from the theory of extremes or heavy tail phenomena. Asymptotic methods come with the obligation to choose an asymptotic regime among potential competing regimes. This is often tantamount to choosing a state space for the observed random elements as well as a scaling. For example, in \mathbb{R}_+^2 , for a risk vector $X = (X_1, X_2)$, if we need to estimate $P[X > x] = P[X_1 > x_1, X_2 > x_2]$ for large x , should the state space for asymptotic analysis be $[0, \infty]^2 \setminus \{(0, 0)\}$ or $(0, \infty]^2$? Ambiguity for the choice of asymptotic regime led to the idea of *coefficient of tail dependence* [10, 11, 28, 29, 36, 45], *hidden regular variation* (hrv) [21, 30, 34, 35, 39–41] and the *conditional extreme value* (cev) model [13–15, 22, 43].

Due to the scaling inherent in the definition of regular variation, a natural domain for regularly varying tails is a region closed under scalar multiplication and usually the domain is a cone centered at the origin. Commonly used cones include \mathbb{R}_+ , \mathbb{R}_+^d , or the two sided versions allowing negative values that are natural in finance and economics. However, as argued in [15], there is need for other cones as well, particularly when asymptotic independence or asymptotic full dependence ([42, Chapter 5], [47]) is present. Going beyond finite-dimensional spaces, there is a need for a comprehensive theory covering spaces such as \mathbb{R}_+^∞ and function spaces. Fortunately a good framework for such a theory of regular variation on metric spaces after removal of a point was created in [24]. The need to remove more than a point, perhaps a closed set and certainly a closed cone, was argued in [15]. These ideas build on $w^\#$ -convergence in [12, Section A2.6].

This paper has a number of goals:

1. We follow the lead of [24] and develop a theory of regularly varying measures on complete separable metric spaces \mathbb{S} with a closed cone \mathbb{C} removed. Section 2 develops a topology on the space of measures on $\mathbb{S} \setminus \mathbb{C}$ which are finite on regions at positive distance from \mathbb{C} . This topology allows creation of mapping theorems (Section 2.1) that encourage continuity arguments and is designed to allow simultaneous regular-variation properties to exist at different scales as is considered in hidden regular variation.

2. We apply the general material of Section 2 to two significant applications.

- (a) In Section 4 we focus on \mathbb{R}_+^p and \mathbb{R}_+^∞ , the space of sequences with non-negative components. An iid sequence $X = (X_1, X_2, \dots) \in \mathbb{R}_+^\infty$ such that $P[X_1 > x]$ is regularly varying has a distribution which is regularly varying on $\mathbb{R}_+^\infty \setminus \mathbb{C}_{\leq j}$ for any $j \geq 1$, where $\mathbb{C}_{\leq j}$ are sequences with at most j positive components. Mapping theorems (Section 2.1) allow extension to the regular-variation properties of $S = (X_1, X_1 + X_2, X_1 + X_2 + X_3, \dots)$ in \mathbb{R}_+^∞ minus the set of non-decreasing sequences which are constant after the j th component. See Section 4.5.2. For reasons of simplicity and taste, we restrict discussion to \mathbb{R}_+^∞ but with modest effort, results could be extended to \mathbb{R}^∞ . We also discuss regular variation of the distribution of a sequence of Poisson points in \mathbb{R}_+^∞ (Section 4.5.4).
- (b) The \mathbb{R}_+^∞ discussion of Poisson points in Section 4.5.4 can be leveraged in a natural way to consider (Section 5) regular variation of the distribution of a Lévy process whose Lévy measure ν is regularly varying: $\lim_{t \rightarrow \infty} t\nu(b(t)x, \infty) = x^{-\alpha}$, $x > 0$, for some scaling function $b(t) \rightarrow \infty$. We reproduce the result [23, 25] that the limit measure of regular variation with scaling $b(t)$ on $\mathbb{D}([0, 1], \mathbb{R}) \setminus \{0\}$ concentrates on càdlàg functions with *one* positive jump. This raises the natural question of what happened to the rest of the jumps of the Lévy process that seem to be *hidden* by the scaling $b(t)$. We are able to generalize for any $j \geq 1$ to convergence under the weaker normalization $b(t^{1/j})$ on a smaller space in which the limit measure concentrates on non-decreasing functions with j positive jumps. Again, as in the study of \mathbb{R}_+^∞ , we focus for simplicity only on large positive jumps of the Lévy process.

3. A final goal is to clarify the proper definition of regular variation in metric spaces. For historical reasons, regular variation is usually associated with scalar multiplication but what does this mean in a general metric space? Traditional definitions are in Cartesian coordinates in finite-dimensional spaces and the form of the definition may not survive change of coordinates. For example, in \mathbb{R}_+^p , a random vector X (in Cartesian coordinates) has a regularly varying distribution if for some scaling function $b(t) \rightarrow \infty$ we have $tP[X/b(t) \in \cdot]$ converging to a limit. If we transform to polar coordinates, $X \mapsto (R, \Theta) := (\|X\|, X/\|X\|)$, the limit is taken on $tP[(R/b(t), \Theta) \in \cdot]$, which appears to be subject to a different notion of scaling. The two convergences are equivalent but look different unless one allows for a more flexible definition of scalar multiplication. We discuss requirements for scalar multiplication in Section 2 along with some examples; related discussion is in [2, 3, 32].

The existing theory for regular variation on, say, \mathbb{R}_+^d , uses the set-up of vague convergence. A troubling consequence is the need to use the one-point uncompactification [40, page 170ff] which adds lines through infinity to the state space.

When regular variation is defined on the cone $[0, \infty]^d \setminus \{0\}$, limit measures cannot charge lines through infinity. However, on proper subcones of $[0, \infty]^d \setminus \{0\}$ this is no longer true and this creates some mathematical havoc: Convergence to types arguments can fail and limit measures may not be unique: Examples in [15, Example 5.4] show that under one normalization the limit measure concentrates on lines through infinity and under another it concentrates on finite points. Another difficulty is that the polar coordinate transform $x \mapsto (\|x\|, x/\|x\|)$ cannot be defined on lines through infinity. One way to patch things up is to retain the one-point un-compactification but demand all limit measures have no mass on lines through infinity, but this does not resolve all difficulties since the unit sphere $\{x : \|x\| = 1\}$ defined by the norm $x \mapsto \|x\|$ may not be compact on a subcone such as $(0, \infty]^d$. Another way forward, which we deem cleaner and more suitable to general spaces where compactification is more involved, is not to compactify and just to define tail regions as subsets of the metric space at positive distance from the deleted closed set. This is the approach given in Section 2.

2. Convergence of measures in the space $\mathbb{M}_\mathbb{O}$

Let (\mathbb{S}, d) be a complete separable metric space. The open ball centered at $x \in \mathbb{S}$ with radius r is written $B_{x,r} = \{y \in \mathbb{S} : d(x, y) < r\}$ and these open sets generate \mathcal{S} , the Borel σ -algebra on \mathbb{S} . For $A \subset \mathbb{S}$, let A° and A^- denote the interior and closure of A , respectively, and let $\partial A = A^- \setminus A^\circ$ be the boundary of A . Let \mathcal{C}_b denote the class of real-valued, non-negative, bounded and continuous functions on \mathbb{S} , and let \mathbb{M}_b denote the class of finite Borel measures on \mathcal{S} . A basic neighborhood of $\mu \in \mathbb{M}_b$ is a set of the form $\{\nu \in \mathbb{M}_b : |\int f_i d\nu - \int f_i d\mu| < \epsilon, i = 1, \dots, k\}$, where $\epsilon > 0$ and $f_i \in \mathcal{C}_b$ for $i = 1, \dots, k$. Thus a sub-basis for \mathbb{M}_b are sets of the form $\{\nu \in \mathbb{M}_b : \nu(f) := \int f d\nu \in G\}$ for $f \in \mathcal{C}_b$ and G open in \mathbb{R}_+ . This equips \mathbb{M}_b with the weak topology and convergence $\mu_n \rightarrow \mu$ in \mathbb{M}_b means $\int f d\mu_n \rightarrow \int f d\mu$ for all $f \in \mathcal{C}_b$. See e.g. Sections 2 and 6 in [6] for details.

Fix a closed set $\mathbb{C} \subset \mathbb{S}$ and set $\mathbb{O} = \mathbb{S} \setminus \mathbb{C}$. For example, one possible choice is $\mathbb{C} = \{s_0\}$ for some $s_0 \in \mathbb{S}$ and then $\mathbb{O} = \mathbb{S} \setminus \{s_0\}$. The subspace \mathbb{O} is a metric subspace of \mathbb{S} in the relative topology with σ -algebra $\mathcal{S}_\mathbb{O} = \mathcal{S}(\mathbb{O}) = \{A : A \subset \mathbb{O}, A \in \mathcal{S}\}$.

Let $\mathcal{C}_\mathbb{O} = \mathcal{C}(\mathbb{O})$ denote the real-valued, non-negative, bounded and continuous functions f on \mathbb{O} such that for each f there exists $r > 0$ such that f vanishes on \mathbb{C}^r ; we use the notation $\mathbb{C}^r = \{x \in \mathbb{S} : d(x, \mathbb{C}) < r\}$, where $d(x, \mathbb{C}) = \inf_{y \in \mathbb{C}} d(x, y)$. Similarly, we will write $d(A, \mathbb{C}) = \inf_{x \in A, y \in \mathbb{C}} d(x, y)$ for $A \subset \mathbb{S}$. We say that a set $A \in \mathcal{S}_\mathbb{O}$ is bounded away from \mathbb{C} if $A \subset \mathbb{S} \setminus \mathbb{C}^r$ for some $r > 0$ or equivalently $d(A, \mathbb{C}) > 0$. So $\mathcal{C}_\mathbb{O}$ consists of non-negative continuous functions whose supports are bounded away from \mathbb{C} . Let $\mathbb{M}_\mathbb{O}$ be the class of Borel measures on \mathbb{O} whose restrictions to $\mathbb{S} \setminus \mathbb{C}^r$ are finite for each $r > 0$. When convenient, we also write $\mathbb{M}(\mathbb{O})$ or $\mathbb{M}(\mathbb{S} \setminus \mathbb{C})$. A basic neighborhood of $\mu \in \mathbb{M}_\mathbb{O}$ is a set of the form $\{\nu \in \mathbb{M}_\mathbb{O} : |\int f_i d\nu - \int f_i d\mu| < \epsilon, i = 1, \dots, k\}$, where $\epsilon > 0$ and $f_i \in \mathcal{C}_\mathbb{O}$

for $i = 1, \dots, k$. A sub-basis is formed by sets of the form

$$\{\nu \in \mathbb{M}_{\mathbb{O}} : \nu(f) \in G\}, \quad f \in \mathcal{C}_{\mathbb{O}}, \quad G \text{ open in } \mathbb{R}_+. \quad (2.1)$$

Convergence $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{O}}$ is convergence in the topology defined by this base or sub-base.

For $\mu \in \mathbb{M}_{\mathbb{O}}$ and $r > 0$, let $\mu^{(r)}$ denote the restriction of μ to $\mathbb{S} \setminus \mathbb{C}^r$. Then $\mu^{(r)}$ is finite and μ is uniquely determined by its restrictions $\mu^{(r)}$, $r > 0$. Moreover, convergence in $\mathbb{M}_{\mathbb{O}}$ has a natural characterization in terms of weak convergence of the restrictions to $\mathbb{S} \setminus \mathbb{C}^r$. A similar result when only a point is removed from the space can be found in [4, 24].

Theorem 2.1 (Portmanteau theorem). *Let $\mu, \mu_n \in \mathbb{M}_{\mathbb{O}}$. The following statements are equivalent.*

- (i) $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{O}}$ as $n \rightarrow \infty$.
- (ii) $\int f d\mu_n \rightarrow \int f d\mu$ for each $f \in \mathcal{C}_{\mathbb{O}}$ which is also uniformly continuous on \mathbb{S} .
- (iii) $\limsup_{n \rightarrow \infty} \mu_n(F) \leq \mu(F)$ and $\liminf_{n \rightarrow \infty} \mu_n(G) \geq \mu(G)$ for all closed $F \in \mathcal{S}_{\mathbb{O}}$ and open $G \in \mathcal{S}_{\mathbb{O}}$ and F and G are bounded away from \mathbb{C} .
- (iv) $\lim_{n \rightarrow \infty} \mu_n(A) = \mu(A)$ for all $A \in \mathcal{S}_{\mathbb{O}}$ bounded away from \mathbb{C} with $\mu(\partial A) = 0$.
- (v) $\mu_n^{(r)} \rightarrow \mu^{(r)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$ for all but at most countably many $r > 0$.
- (vi) There exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that $\mu_n^{(r_i)} \rightarrow \mu^{(r_i)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ for each i .

For proofs, see Section 2.4.

Weak convergence is metrizable (for instance by the Prohorov metric; see e.g. p. 72 in [6]) and the close relation between weak convergence and convergence in $\mathbb{M}_{\mathbb{O}}$ in Theorem 2.1(v)–(vi) indicates that the topology in $\mathbb{M}_{\mathbb{O}}$ is metrizable too. With minor modifications of the arguments in [12], pp. 627–628, we may choose the metric

$$d_{\mathbb{M}_{\mathbb{O}}}(\mu, \nu) = \int_0^\infty e^{-r} p_r(\mu^{(r)}, \nu^{(r)}) [1 + p_r(\mu^{(r)}, \nu^{(r)})]^{-1} dr, \quad (2.2)$$

where $\mu^{(r)}, \nu^{(r)}$ are the finite restrictions of μ, ν to $\mathbb{S} \setminus \mathbb{C}^r$, and p_r is the Prohorov metric on $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$.

Theorem 2.2. $(\mathbb{M}_{\mathbb{O}}, d_{\mathbb{M}_{\mathbb{O}}})$ is a separable and complete metric space.

2.1. Mapping theorems

Applications of weak convergence often rely on continuous mapping theorems and we present versions for convergence in $\mathbb{M}_{\mathbb{O}}$. Consider another separable and complete metric space \mathbb{S}' and let $\mathbb{O}', \mathcal{S}_{\mathbb{O}'}, \mathbb{C}', \mathbb{M}_{\mathbb{O}'}$ have the same meaning relative to the space \mathbb{S}' as do $\mathbb{O}, \mathcal{S}_{\mathbb{O}}, \mathbb{C}, \mathbb{M}_{\mathbb{O}}$ relative to \mathbb{S} .

Theorem 2.3 (Mapping theorem). *Let $h : (\mathbb{O}, \mathcal{S}_{\mathbb{O}}) \mapsto (\mathbb{O}', \mathcal{S}_{\mathbb{O}'})$ be a measurable mapping such that $h^{-1}(A')$ is bounded away from \mathbb{C} for any $A' \in \mathcal{S}_{\mathbb{O}'} \cap h(\mathbb{O})$ bounded away from \mathbb{C}' . Then $\hat{h} : \mathbb{M}_{\mathbb{O}} \mapsto \mathbb{M}_{\mathbb{O}'}$ defined by $\hat{h}(\nu) = \nu \circ h^{-1}$ is continuous at μ provided $\mu(D_h) = 0$, where D_h is the set of discontinuity points of h .*

This result is illustrated in Examples 3.3 and 3.4 and is also needed for considering the generalized polar coordinate transformation in Section 4.2.3. It is the basis for the approach to regular variation of Lévy processes in Section 5. Theorem 2.3 is formulated so that h is defined on $\mathbb{O} = \mathbb{S} \setminus \mathbb{C}$, rather than on all of \mathbb{S} . If $\mathbb{S} = \mathbb{R}_+^p$ and $h(x) = (\|x\|, x/\|x\|)$ is the polar coordinate transform, then h is not defined at 0. This lack of definition is not a problem since

$$\begin{aligned} h : \mathbb{O} := \mathbb{R}_+^p \setminus \{0\} &\mapsto \mathbb{O}' := (0, \infty) \times \{x \in \mathbb{R}_+^p : \|x\| = 1\} \\ &= [0, \infty) \times \{x \in \mathbb{R}_+^p : \|x\| = 1\} \setminus (\{0\} \times \{x \in \mathbb{R}_+^p : \|x\| = 1\}). \end{aligned}$$

The proof of Theorem 2.3 is in Section 2.4.4 but it is instructive to quickly consider the special case where $D_h = \emptyset$ so that h is continuous. In this case h induces a continuous mapping $\hat{h} : \mathbb{M}_{\mathbb{O}} \mapsto \mathbb{M}'_{\mathbb{O}'}$ defined by $\hat{h}(\mu) = \mu \circ h^{-1}$. To see this, look at the inverse image of a sub-basis set (2.1): For G open in \mathbb{R}_+ , and $f' \in \mathcal{C}_{\mathbb{O}'}$,

$$\begin{aligned} \hat{h}^{-1}\{\mu' \in \mathbb{M}_{\mathbb{O}'} : \mu'(f') \in G\} &= \{\mu \in \mathbb{M}_{\mathbb{O}} : \mu \circ h^{-1}(f') \in G\} \\ &= \{\mu \in \mathbb{M}_{\mathbb{O}} : \mu(f' \circ h) \in G\}. \end{aligned}$$

Since h is continuous and $f' \circ h \in \mathcal{C}_{\mathbb{O}}$, $\{\mu \in \mathbb{M}_{\mathbb{O}} : \mu(f' \circ h) \in G\}$ is open in $\mathbb{M}_{\mathbb{O}}$.

Here are two variants of the mapping theorem. The first allows application of the operator taking successive partial sums from $\mathbb{R}_+^\infty \mapsto \mathbb{R}_+^\infty$ in Proposition 4.2 and also allows application of the projection map $(x_1, x_2, \dots) \mapsto (x_1, \dots, x_p)$ from $\mathbb{R}_+^\infty \mapsto \mathbb{R}_+^p$ in Proposition 4.3. The second variant allows a quick proof that the polar coordinate transform is continuous on $\mathbb{R}_+^p \setminus \{0\}$ in Corollary 4.3.

Corollary 2.1. *Suppose $h : \mathbb{S} \mapsto \mathbb{S}'$ is uniformly continuous and $\mathbb{C}' := h(\mathbb{C})$ is closed in \mathbb{S}' . Then $\hat{h} : \mathbb{M}_{\mathbb{O}} \mapsto \mathbb{M}_{\mathbb{O}'}$ defined by $\hat{h}(\mu) = \mu \circ h^{-1}$ is continuous.*

Corollary 2.2. *Suppose $h : \mathbb{S} \mapsto \mathbb{S}'$ is continuous and either \mathbb{S} or \mathbb{C} is compact. Then $\hat{h} : \mathbb{M}_{\mathbb{O}} \mapsto \mathbb{M}_{\mathbb{O}'}$ defined by $\hat{h}(\mu) = \mu \circ h^{-1}$ is continuous.*

2.2. Relative compactness in $\mathbb{M}_{\mathbb{O}}$

Proving convergence sometimes requires a characterization of relative compactness. A subset of a topological space is relatively compact if its closure is compact. A subset of a metric space is compact if and only if it is sequentially compact. Hence, $M \subset \mathbb{M}_{\mathbb{O}}$ is relatively compact if and only if every sequence $\{\mu_n\}$ in M contains a convergent subsequence. For $\mu \in M \subset \mathbb{M}_{\mathbb{O}}$ and $r > 0$, let $\mu^{(r)}$ be the restriction of μ to $\mathbb{S} \setminus \mathbb{C}^r$ and $M^{(r)} = \{\mu^{(r)} : \mu \in M\}$. By Theorem 2.1(vi) we have the following characterization of relative compactness.

Theorem 2.4. *A subset $M \subset \mathbb{M}_{\mathbb{O}}$ is relatively compact if and only if there exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that $M^{(r_i)}$ is relatively compact in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ for each i .*

Prohorov’s theorem characterizes relative compactness in the weak topology. This translates to a characterization of relative compactness in $\mathbb{M}_{\mathbb{O}}$.

Theorem 2.5. $M \subset \mathbb{M}_0$ is relatively compact if and only if there exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that for each i

$$\sup_{\mu \in M} \mu(\mathbb{S} \setminus \mathbb{C}^{r_i}) < \infty, \tag{2.3}$$

and for each $\eta > 0$ there exists a compact set $K_i \subset \mathbb{S} \setminus \mathbb{C}^{r_i}$ such that

$$\sup_{\mu \in M} \mu(\mathbb{S} \setminus (\mathbb{C}^{r_i} \cup K_i)) \leq \eta. \tag{2.4}$$

2.3. \mathbb{M} -convergence vs vague convergence

Vague convergence complies with the topology on the space of measures which are finite on compacta. Regular variation for measures on a space such as \mathbb{R}_+^p has traditionally been formulated using vague convergence after compactification of the space. In order to make use of existing regular variation theory on \mathbb{R}_+^p , it is useful to understand how \mathbb{M} -convergence is related to vague convergence.

Let \mathbb{S} be a complete separable metric space and suppose \mathbb{C} is closed in \mathbb{S} . Then $\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$ is the collection of measures finite on $\mathcal{K}(\mathbb{S} \setminus \mathbb{C})$, the compacta of $\mathbb{S} \setminus \mathbb{C}$:

$$\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C}) = \{\mu : \mu(K) < \infty, \forall K \in \mathcal{K}(\mathbb{S} \setminus \mathbb{C})\}.$$

Vague convergence on $\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$ means $\mu \mapsto \mu(f)$ is continuous for all $f \in \mathcal{C}_K^+(\mathbb{S} \setminus \mathbb{C})$, the continuous functions with compact support. The spaces $\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$ and $\mathbb{M}(\mathbb{S} \setminus \mathbb{C})$ are not the same. For example if $\mathbb{S} = [0, \infty)$ and $\mathbb{C} = \{0\}$, $\mu \in \mathbb{M}(\mathbb{S} \setminus \mathbb{C})$ means $\mu(x, \infty) < \infty$ for $x > 0$ but $\mu \in \mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$ means $\mu([a, b]) < \infty$ for $0 < a < b < \infty$. For instance Lebesgue measure is in $\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$ but not in $\mathbb{M}(\mathbb{S} \setminus \mathbb{C})$.

2.3.1. Comparing \mathbb{M} vs \mathbb{M}_+

We have the following comparison.

Lemma 2.1. \mathbb{M} -convergence implies vague convergence and

$$\mathbb{M}(\mathbb{S} \setminus \mathbb{C}) \subset \mathbb{M}_+(\mathbb{S} \setminus \mathbb{C}), \quad \mathcal{C}_K^+(\mathbb{S} \setminus \mathbb{C}) \subset \mathcal{C}(\mathbb{S} \setminus \mathbb{C}). \tag{2.5}$$

Proof. If $f \in \mathcal{C}_K^+(\mathbb{S} \setminus \mathbb{C})$, its compact support $K \subset \mathbb{S} \setminus \mathbb{C}$ must be bounded away from \mathbb{C} and hence $d(K, \mathbb{C}) > 0$ and $f \in \mathcal{C}(\mathbb{S} \setminus \mathbb{C})$. If $\mu \in \mathbb{M}(\mathbb{S} \setminus \mathbb{C})$, and \mathbb{D} satisfies $d(\mathbb{D}, \mathbb{C}) > 0$, then $\mu(\mathbb{D}) < \infty$. If $K \in \mathcal{K}(\mathbb{S} \setminus \mathbb{C})$ then $d(K, \mathbb{C}) > 0$ and so $\mu(K) < \infty$, showing any $\mu \in \mathbb{M}(\mathbb{S} \setminus \mathbb{C})$ is also in $\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$. \square

Remark. Let $\mathbb{S} = [0, \infty)$ and $\mathbb{C} = \{0\}$, and

$$\mu_n = \frac{1}{n} \sum_{i=1}^{n^2} \epsilon_{i/n} \in \mathbb{M}(\mathbb{S} \setminus \mathbb{C}) \subset \mathbb{M}_+(\mathbb{S} \setminus \mathbb{C}).$$

Here and elsewhere, we use the notation ϵ_x for the Dirac measure concentrating mass 1 on the point x so that $\epsilon_x(A) = 1$, if $x \in A$, and $\epsilon_x(A) = 0$, if $x \in A^c$. We have μ_n converging to Lebesgue measure in $\mathbb{M}_+(\mathbb{S} \setminus \mathbb{C})$ but $\{\mu_n\}$ does not converge in $\mathbb{M}(\mathbb{S} \setminus \mathbb{C})$. If f is 0 on $(0, 1)$, linear on $(1, 2)$ and $f \equiv 1$ on $[2, \infty)$, then $f \in \mathcal{C}(\mathbb{S} \setminus \mathbb{C})$ but $\mu_n(f) \geq \frac{1}{n} \sum_{i=2n+1}^{n^2} 1 = n - 2 \rightarrow \infty$.

2.4. Proofs

2.4.1. Preliminaries

We begin with two well known preliminary lemmas in topology. The second one is just a version of Urysohn's lemma [18, 44] for metric spaces.

Lemma 2.2. *Fix a set $B \subset \mathbb{S}$. Then*

- (i) $d(x, B)$ is a uniformly continuous function in x .
- (ii) $d(x, B) = 0$ if and only if $x \in B^-$.

Proof. (i) follows from the following generalization of the triangle inequality. For $x, y \in \mathbb{S}$,

$$d(x, B) \leq d(x, y) + d(y, B).$$

(ii) is an easy deduction from the definition of $d(x, B) = \inf_{z \in B} d(x, z)$. \square

Lemma 2.3. *For any two closed sets $A, B \subset \mathbb{S}$ such that $A \cap B = \emptyset$, there exists a uniformly continuous function f from \mathbb{S} to $[0, 1]$ such that $f \equiv 0$ on A and $f \equiv 1$ on B*

Proof. Define the function f as

$$f(x) = \frac{d(x, A)}{d(x, A) + d(x, B)}.$$

The desired properties of f are easily checked from Lemma 2.2. \square

Lemma 2.4. *If $A \in \mathcal{S}_0$ is bounded away from \mathbb{C} , $A = \cup_{i \in I} A_i$ for an uncountable index set I , disjoint sets $A_i \in \mathcal{S}_0$, and $\mu(A) < \infty$, then $\mu(A_i) > 0$ for at most countably many i .*

Proof. Suppose there exists a countably infinite set I_n such that $\mu(A_i) > 1/n$ for $i \in I_n$. Then

$$\infty = \sum_{i \in I_n} \mu(A_i) = \mu\left(\cup_{i \in I_n} A_i\right) \leq \mu(A)$$

which is a contradiction to the assumption that $\mu(A) < \infty$. The conclusion follows from letting $n \rightarrow \infty$. \square

Lemma 2.5. *For any $\mu \in \mathbb{M}_0$, $\mu(\partial(\mathbb{S} \setminus \mathbb{C}^\delta)) > 0$ for at most countably many $\delta > 0$.*

Proof. Notice first that $\partial(\mathbb{S} \setminus \mathbb{C}^\delta) = \{x \in \mathbb{S} : d(x, \mathbb{C}) = \delta\}$ so $\partial(\mathbb{S} \setminus \mathbb{C}^{\delta_1}) \cap \partial(\mathbb{S} \setminus \mathbb{C}^{\delta_2}) = \emptyset$ for $\delta_1 \neq \delta_2$. The conclusion follows from Lemma 2.4. \square

2.4.2. Proof of Theorem 2.1

We show that (i) \Rightarrow (ii), (ii) \Rightarrow (iii), (iii) \Rightarrow (iv), (iv) \Rightarrow (v), (v) \Rightarrow (vi) and (vi) \Rightarrow (i).

Suppose that (i) holds. Suppose $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}}$ and take $f \in \mathcal{C}_{\mathbb{D}}$. Given $\epsilon > 0$ consider the neighborhood $N_{\epsilon, f}(\mu) = \{\nu : |\int f d\nu - \int f d\mu| < \epsilon\}$. By assumption there exists n_0 such that $n \geq n_0$ implies $\mu_n \in N_{\epsilon, f}(\mu)$, i.e. $|\int f d\mu_n - \int f d\mu| < \epsilon$. Hence $\int f d\mu_n \rightarrow \int f d\mu$.

Suppose that (ii) holds. Take any closed F that is bounded away from \mathbb{C} . Then there exists $r > 0$ such that $F \subset \mathbb{S} \setminus \mathbb{C}^r$. So for all $x \in F$, $d(x, \mathbb{C}) \geq r$. So if we define $F^\epsilon = \{x \in \mathbb{S} : d(x, F) < \epsilon\}$, then each F^ϵ is open, $F \subset F^\epsilon$ and $F^\epsilon \downarrow F$ as $\epsilon \downarrow 0$. Also for $\epsilon < r/2$, we have that for all $x \in F^\epsilon$ $d(x, \mathbb{C}) \geq r - r/2 = r/2$, meaning that $F^\epsilon \subset \mathbb{S} \setminus \mathbb{C}^{r/2}$. For $\epsilon > 0$, $\mathbb{S} \setminus F^\epsilon$ is closed and $F \cap (\mathbb{S} \setminus F^\epsilon) = \emptyset$. So for $0 < \epsilon < r/2$, by Lemma 2.3, there exists a uniformly continuous function f from \mathbb{S} to $[0, 1]$ such that $f \equiv 0$ on $\mathbb{S} \setminus F^\epsilon$ and $f \equiv 1$ on F . Observe that $f \in \mathcal{C}_{\mathbb{D}}$ as $F^\epsilon \subset \mathbb{S} \setminus \mathbb{C}^{r/2}$. So we have

$$\limsup_{n \rightarrow \infty} \mu_n(F) \leq \lim_{n \rightarrow \infty} \int f d\mu_n = \int f d\mu \leq \mu(F^\epsilon).$$

As $\epsilon \downarrow 0$, $F^\epsilon \downarrow F$ and as F is closed, we have $\mu(F^\epsilon) \downarrow \mu(F)$. This leads to

$$\limsup_{n \rightarrow \infty} \mu_n(F) \leq \mu(F).$$

Now take any open G bounded away from \mathbb{C} . Then there exists $r > 0$ such that $G \subset \mathbb{S} \setminus \mathbb{C}^r$. So if we define $G_\epsilon = \mathbb{S} \setminus \{x \in \mathbb{S} \setminus G : d(x, \mathbb{S} \setminus G) < \epsilon\}$, then each G_ϵ is closed, $G_\epsilon \subset G$ and $G_\epsilon \uparrow G$ as $\epsilon \downarrow 0$. So by Lemma 2.3, there exists a uniformly continuous function f from \mathbb{S} to $[0, 1]$ such that $f \equiv 0$ on $\mathbb{S} \setminus G$ and $f \equiv 1$ on G_ϵ . Observe that $f \in \mathcal{C}_{\mathbb{D}}$ as $G \subset \mathbb{S} \setminus \mathbb{C}^r$. So we have

$$\liminf_{n \rightarrow \infty} \mu_n(G) \geq \lim_{n \rightarrow \infty} \int f d\mu_n = \int f d\mu \geq \mu(G_\epsilon).$$

As $\epsilon \downarrow 0$, $G_\epsilon \uparrow G$ and as G is open, we have $\mu(G_\epsilon) \uparrow \mu(G)$. This leads to

$$\liminf_{n \rightarrow \infty} \mu_n(G) \geq \mu(G).$$

This completes the proof of (iii).

Suppose that (iii) holds and take $A \in \mathcal{S}_{\mathbb{D}}$ bounded away from \mathbb{C} with $\mu(\partial A) = 0$.

$$\begin{aligned} \limsup_{n \rightarrow \infty} \mu_n(A) &\leq \limsup_{n \rightarrow \infty} \mu_n(A^-) \leq \mu(A^-) \\ &= \mu(A^\circ) \leq \liminf_{n \rightarrow \infty} \mu_n(A^\circ) \leq \liminf_{n \rightarrow \infty} \mu_n(A). \end{aligned}$$

Hence, $\lim_{n \rightarrow \infty} \mu_n(A) = \mu(A)$, so that (iv) holds.

Suppose that (iv) holds and take $r > 0$ such that $\mu(\partial(\mathbb{S} \setminus \mathbb{C}^r)) = 0$. By Lemma 2.5, all but at most countably many $r > 0$ satisfy this property. As $\mathbb{S} \setminus \mathbb{C}^r$

is trivially bounded away from \mathbb{C} , we have that $\mu_n(\mathbb{S} \setminus \mathbb{C}^r) \rightarrow \mu(\mathbb{S} \setminus \mathbb{C}^r)$. Now any $A \subset \mathbb{S} \setminus \mathbb{C}^r$ is also bounded away from \mathbb{C} and as $\mathbb{S} \setminus \mathbb{C}^r$ is closed, $\partial_{\mathbb{S} \setminus \mathbb{C}^r} A = \partial A$, where the first expression denotes the boundary of A when considered as a subset of $\mathbb{S} \setminus \mathbb{C}^r$. So for any subset $A \subset \mathbb{S} \setminus \mathbb{C}^r$ with $\mu(\partial_{\mathbb{S} \setminus \mathbb{C}^r} A) = 0$, we have by (iv) that $\mu_n(A) \rightarrow \mu(A)$ and hence $\mu_n^{(r)}(A) \rightarrow \mu^{(r)}(A)$. The Portmanteau theorem for weak convergence implies $\mu_n^{(r)} \rightarrow \mu^{(r)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$. This completes the proof of (v).

Suppose that (v) holds. Since, $\mu_n^{(r)} \rightarrow \mu^{(r)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$ for all but at most countably many $r > 0$ we can always choose a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that $\mu_n^{(r_i)} \rightarrow \mu^{(r_i)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ for each i .

Suppose that (vi) holds. Take $\epsilon > 0$ and a neighborhood $N_{\epsilon, f_1, \dots, f_k}(\mu) = \{\nu : |\int f_j d\nu - \int f_j d\mu| < \epsilon, j = 1, \dots, k\}$ where each $f_j \in \mathcal{C}_{\mathbb{D}}$ for $j = 1, 2, \dots, k$. Let $r > 0$ be such that $\mu_n^{(r)} \rightarrow \mu^{(r)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$ and each f_j vanishes on \mathbb{C}^r . Let n_j be an integer such that $n \geq n_j$ implies $|\int f_j d\mu_n^{(r)} - \int f_j d\mu^{(r)}| < \epsilon$. Hence, $n \geq \max(n_1, \dots, n_k)$ implies that $|\int f_j d\mu_n^{(r)} - \int f_j d\mu^{(r)}| < \epsilon$ for all $j = 1, 2, \dots, k$. As each f_j vanishes outside \mathbb{C}^r , we also have that $|\int f_j d\mu_n - \int f_j d\mu| < \epsilon$ for all $j = 1, 2, \dots, k$. So $\mu_n \in N_{\epsilon, f_1, \dots, f_k}(\mu)$. Hence $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}}$. \square

2.4.3. Proof of Theorem 2.2

The proof consists of minor modifications of arguments that can be found in [12], pp. 628–630. Here we change from r to $1/r$. For the sake of completeness we have included a full proof.

We show that (i) $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}}$ if and only if $d_{\mathbb{M}_{\mathbb{D}}}(\mu_n, \mu) \rightarrow 0$, and (ii) $(\mathbb{M}_{\mathbb{D}}, d_{\mathbb{M}_{\mathbb{D}}})$ is separable and complete.

(i) Suppose that $d_{\mathbb{M}_{\mathbb{D}}}(\mu_n, \mu) \rightarrow 0$. The integral expression in (2.2) can be written $d_{\mathbb{M}_{\mathbb{D}}}(\mu_n, \mu) = \int_0^{\infty} e^{-r} g_n(r) dr$, so that for each n , $g_n(r)$ decreases with r and is bounded by 1. Helly's selection theorem (p. 336 in [7]), applied to $1 - g_n$, implies that there exists a subsequence $\{n'\}$ and a nonincreasing function g such that $g_{n'}(r) \rightarrow g(r)$ for all continuity points of g . By dominated convergence, $\int_0^{\infty} e^{-r} g(r) dr = 0$ and since g is monotone this implies that $g(r) = 0$ for all finite $r > 0$. Since this holds for all convergent subsequences $\{g_{n'}(r)\}$, it follows that $g_n(r) \rightarrow 0$ for all continuity points r of g , and hence, for such r , $p_r(\mu_n^{(r)}, \mu^{(r)}) \rightarrow 0$ as $n \rightarrow \infty$. By Theorem 2.1(vi), $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}}$.

Suppose that $\mu_n \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}}$. Theorem 2.1(v) implies that $\mu_n^{(r)} \rightarrow \mu^{(r)}$ in $\mathbb{M}_b(\mathbb{D} \setminus \mathbb{C}^r)$ for all but at most countably many $r > 0$. Hence, for such r , $p_r(\mu_n^{(r)}, \mu^{(r)})[1 + p_r(\mu_n^{(r)}, \mu^{(r)})]^{-1} \rightarrow 0$, which by the dominated convergence theorem implies that $d_{\mathbb{M}_{\mathbb{D}}}(\mu_n, \mu) \rightarrow 0$.

(ii) Separability: For $r > 0$ let D_r be a countable dense set in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$ with the weak topology. Let D be the union of D_r for rational $r > 0$. Then D is countable. Let us show D is dense in $\mathbb{M}_{\mathbb{D}}$. Given $\epsilon > 0$ and $\mu \in \mathbb{M}_{\mathbb{D}}$ pick $r' > 0$ such that $\int_0^{r'} e^{-r} dr < \epsilon/2$. Take $\mu_{r'} \in D_{r'}$ such that $p_{r'}(\mu_{r'}, \mu^{(r')}) < \epsilon/2$. Then $p_r(\mu_{r'}^{(r)}, \mu^{(r)}) < \epsilon/2$ for all $r > r'$. In particular, $d_{\mathbb{M}_{\mathbb{D}}}(\mu_{r'}, \mu) < \epsilon$.

Completeness: Let $\{\mu_n\}$ be a Cauchy sequence for $d_{\mathbb{M}_0}$. Then $\{\mu_n^{(r)}\}$ is a Cauchy sequence for p_r for all but at most countably many $r > 0$. Since \mathbb{S} is separable and complete, its closed subspace $\mathbb{S} \setminus \mathbb{C}^r$ is separable and complete. Therefore, $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^r)$ is complete, which implies that $\{\mu_n^{(r)}\}$ has a limit μ_r . These limits are consistent in the sense that $\mu_{r'}^{(r)} = \mu_r$ for $r' < r$. On \mathcal{S}_0 set $\mu(A) = \lim_{r \rightarrow 0} \mu_r(A \cap \mathbb{S} \setminus \mathbb{C}^r)$. Then μ is a measure. Clearly, $\mu \geq 0$ and $\mu(\emptyset) = 0$. Moreover, μ is countably additive: for disjoint $A_n \in \mathcal{S}_0$ the monotone convergence theorem implies that

$$\begin{aligned} \mu(\cup_n A_n) &= \lim_{r \rightarrow 0} \mu_r(\cup_n A_n \cap [\mathbb{S} \setminus \mathbb{C}^r]) \\ &= \lim_{r \rightarrow 0} \sum_n \mu_r(A_n \cap [\mathbb{S} \setminus \mathbb{C}^r]) = \sum_n \mu(A_n), \end{aligned}$$

proving the result. □

2.4.4. Proof of Theorem 2.3

Firstly, $D_h \in \mathcal{S}_0$ [6, p. 243]. Take $A' \in \mathcal{S}_{0'}$ bounded away from \mathbb{C} with $\mu \circ h^{-1}(\partial A') = 0$. Since $\partial h^{-1}(A') \subset h^{-1}(\partial A') \cup D_h$ (see e.g. (A2.3.2) in [12]), we have $\mu(\partial h^{-1}(A')) \leq \mu h^{-1}(\partial A') + \mu(D_h) = 0$. Since $\mu_n \rightarrow \mu$ in \mathbb{M}_0 , $\mu(\partial h^{-1}(A')) = 0$, and $h^{-1}(A')$ is bounded away from \mathbb{C} , it follows from Theorem 2.1(iv) that $\mu_n h^{-1}(A) \rightarrow \mu h^{-1}(A)$. Hence, $\mu_n h^{-1} \rightarrow \mu h^{-1}$ in $\mathbb{M}_{0'}$. □

2.4.5. Proof of Corollary 2.1

Take $A' \subset \mathbb{S}' \setminus \mathbb{C}'$ such that $d'(A', \mathbb{C}') > 0$. We claim this implies $d(h^{-1}(A'), \mathbb{C}) > 0$. Otherwise, if $d(h^{-1}(A'), \mathbb{C}) = 0$, there exist $x_n \in h^{-1}(A')$ and $y_n \in \mathbb{C}$ such that $d(x_n, y_n) \rightarrow 0$. Then $h(x_n) \in A'$, $h(y_n) \in h(\mathbb{C}) = \mathbb{C}'$ and if h is uniformly continuous, then $d'(h(x_n), h(y_n)) \rightarrow 0$ so that $d'(A', \mathbb{C}') = 0$, giving us a contradiction. □

2.4.6. Proof of Corollary 2.2

The proof of Corollary 2.1 shows that it suffices if either $\{x_n\}$ or $\{y_n\}$ has a limit point. In the former case, if $x_{n'} \rightarrow x$ for some subsequence $n' \rightarrow \infty$, then $d(x, y_{n'}) \rightarrow 0$ and $y_{n'} \rightarrow x \in \mathbb{C}$ and $h(y_{n'}) \rightarrow h(x)$ so $d'(A', \mathbb{C}') = 0$ again giving a contradiction. Note if \mathbb{S} is compact than $\{x_n\}$ has a limit point. On the other hand, if $\{y_n\}$ has a limit point then there exists an infinite subsequence $\{n'\}$ and $y_{n'} \rightarrow y \in \mathbb{C}$ so that $d(x_{n'}, y) \rightarrow 0$. Thus if h is continuous, $h(x_{n'}) \rightarrow h(y) \in h(\mathbb{C}) = \mathbb{C}'$ which contradicts $d'(A', \mathbb{C}') > 0$. Note if \mathbb{C} is compact, then $\{y_n\}$ has a limit point and in particular if $\mathbb{C} = \{s_0\}$. Thus we have the second variant. □

2.4.7. Proof of Theorem 2.4

Suppose $M \subset \mathbb{M}_{\mathbb{O}}$ is relatively compact. Let $\{\mu_n\}$ be a subsequence in M . Then there exists a convergent subsequence $\mu_{n_k} \rightarrow \mu$ for some $\mu \in M^-$. By Theorem 2.1(v), there exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that $\mu_{n_k}^{(r_i)} \rightarrow \mu^{(r_i)}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$. Hence, $M^{(r_i)}$ is relatively compact in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ for each such r_i .

Conversely, suppose there exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that $M^{(r_i)} \subset \mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ is relatively compact for each i , and let $\{\mu_n\}$ be a sequence of elements in M . We use a diagonal argument to find a convergent subsequence. Since $M^{(r_1)}$ is relatively compact there exists a subsequence $\{\mu_{n_1(k)}\}$ of $\{\mu_n\}$ such that $\mu_{n_1(k)}^{(r_1)}$ converges to some μ_{r_1} in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_1})$. Similarly since $M^{(r_2)}$ is relatively compact and $\{\mu_{n_1(k)}\} \subset M$ there exists a subsequence $\{\mu_{n_2(k)}\}$ of $\{\mu_{n_1(k)}\}$ such that $\mu_{n_2(k)}^{(r_2)}$ converges to some μ_{r_2} in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_2})$. Continuing like this; for each $i \geq 3$ let $n_i(k)$ be a subsequence of $n_{i-1}(k)$ such that $\mu_{n_i(k)}^{(r_i)}$ converges to some μ_{r_i} in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$. Then the diagonal sequence $\{\mu_{n_k(k)}\}$ satisfies $\mu_{n_k(k)}^{(r_i)} \rightarrow \mu_{r_i}$ in $\mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ for each $i \geq 1$. Take $f \in \mathcal{C}_{\mathbb{O}}$. There exists some $i_0 \geq 1$ such that f vanishes on $\mathbb{S} \setminus \mathbb{C}^{r_i}$ for each $i \geq i_0$. In particular $f \in \mathcal{C}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ for each $i \geq i_0$ and

$$\int f d\mu_{r_i} = \lim_k \int f d\mu_{n_k(k)}^{(r_i)} = \lim_k \int f d\mu_{n_k(k)}^{(r_{i_0})} = \int f d\mu_{r_{i_0}}.$$

Hence, we can define $\mu' : \mathcal{C}_{\mathbb{O}} \rightarrow [0, \infty]$ by $\mu'(f) = \lim_{i \rightarrow \infty} \int f d\mu_{r_i}$. This μ' induces a measure μ in $\mathbb{M}_{\mathbb{O}}$. Indeed, for $A \in \mathcal{S}_{\mathbb{O}}$ we can find a sequence $f_n \in \mathcal{C}_{\mathbb{O}}$ such that $0 \leq f_n \uparrow 1_A$ and set $\mu(A) = \lim_n \mu'(f_n)$. If $A \in \mathbb{S} \setminus \mathbb{C}^r$ for some $r > 0$, then there exists $f_n \in \mathcal{C}_{\mathbb{O}}$ such that $f_n \downarrow 1_A$ and hence $\mu(A) \leq \mu'(f_n) < \infty$. Thus, μ is finite on sets $A \in \mathbb{S} \setminus \mathbb{C}^r$ for some $r > 0$. To show that μ is countably additive, let A_1, A_2, \dots be disjoint sets in $\mathcal{S}_{\mathbb{O}}$ and $0 \leq f_{nk} \uparrow 1_{A_k}$ for each k . Then $\sum_k f_{nk} \uparrow 1_{\cup_k A_k}$ and, by Fubini's theorem and the monotone convergence theorem, it holds that

$$\mu(\cup_k A_k) = \lim_n \mu' \left(\sum_k f_{nk} \right) = \sum_k \lim_n \mu'(f_{nk}) = \sum_k \mu(A_k).$$

By construction $\int f d\mu = \mu'(f)$ for each $f \in \mathcal{C}_{\mathbb{O}}$. Hence, $\int f d\mu_{n_k(k)} \rightarrow \int f d\mu$ for each $f \in \mathcal{C}_{\mathbb{O}}$, and we conclude that M is relatively compact in $\mathbb{M}_{\mathbb{O}}$. \square

2.4.8. Proof of Theorem 2.5

Suppose $M \subset \mathbb{M}_{\mathbb{O}}$ is relatively compact. By Theorem 2.4, there exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that $M^{(r_i)} \subset \mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ is relatively compact for each r_i . Prohorov's theorem (Theorem A2.4.1 in [12]) implies that (2.3) and (2.4) hold.

Conversely, suppose there exists a sequence $\{r_i\}$ with $r_i \downarrow 0$ such that (2.3) and (2.4) hold. Then, by Prohorov's theorem, $M^{(r_i)} \subset \mathbb{M}_b(\mathbb{S} \setminus \mathbb{C}^{r_i})$ is relatively compact for each i . By Theorem 2.4, $M \subset \mathbb{M}_{\mathbb{O}}$ is relatively compact. \square

3. Regularly varying sequences of measures

3.1. Scaling

The usual notion of regular variation involves comparisons along a ray and requires a concept of scaling or multiplication. We approach the scaling idea in a general complete, separable metric space \mathbb{S} by postulating what is required for a pleasing theory. Given any real number $\lambda > 0$ and any $x \in \mathbb{S}$, we assume there exists a mapping $(\lambda, x) \mapsto \lambda x$ from $(0, \infty) \times \mathbb{S}$ into \mathbb{S} satisfying:

- (A1) the mapping $(\lambda, x) \mapsto \lambda x$ is continuous,
- (A2) $1x = x$ and $\lambda_1(\lambda_2 x) = (\lambda_1 \lambda_2)x$.

Assumptions (A1) and (A2) allow definition of a cone $\mathbb{C} \subset \mathbb{S}$ as a set satisfying $x \in \mathbb{C}$ implies $\lambda x \in \mathbb{C}$ for any $\lambda > 0$. For this section, fix a closed cone $\mathbb{C} \subset \mathbb{S}$ and then $\mathbb{O} := \mathbb{S} \setminus \mathbb{C}$ is an open cone. We require that

- (A3) $d(x, \mathbb{C}) < d(\lambda x, \mathbb{C})$ if $\lambda > 1$ and $x \in \mathbb{O}$.

3.1.1. Examples to fix ideas

To emphasize the flexibility allowed by our assumptions, consider the following circumstances all of which satisfy (A1)–(A3).

1. Let $\mathbb{S} = \mathbb{R}^2$ and $\mathbb{C} = (\{0\} \times \mathbb{R}) \cup (\mathbb{R} \times \{0\})$ and for $\gamma_1 > 0, \gamma_2 > 0$ define $(\lambda, (x_1, x_2)) \mapsto (\lambda^{1/\gamma_1} x_1, \lambda^{1/\gamma_2} x_2)$.
2. Set $\mathbb{S} = \mathbb{R}^2$ and $\mathbb{C} = \mathbb{R} \times \{0\}$. Define $(\lambda, (x_1, x_2)) \mapsto (x_1, \lambda x_2)$.
3. Set $\mathbb{S} = [0, \infty) \times \{x \in \mathbb{R}_+^2 : \|x\| = 1\}$ and $\mathbb{C} = \{0\} \times \{x \in \mathbb{R}_+^2 : \|x\| = 1\}$. For $\lambda > 0$, define $(\lambda, (r, a)) \mapsto (\lambda r, a)$.

3.2. Regular variation

Recall from e.g. [8] that a positive measurable function $c(\cdot)$ defined on $(0, \infty)$ is regularly varying with index $\rho \in \mathbb{R}$ if $\lim_{t \rightarrow \infty} c(\lambda t)/c(t) = \lambda^\rho$ for all $\lambda > 0$. Similarly, a sequence $\{c_n\}_{n \geq 1}$ of positive numbers is regularly varying with index $\rho \in \mathbb{R}$ if $\lim_{n \rightarrow \infty} c_{[\lambda n]}/c_n = \lambda^\rho$ for all $\lambda > 0$. Here $[\lambda n]$ denotes the integer part of λn .

Definition 3.1. A sequence $\{\nu_n\}_{n \geq 1}$ in $\mathbb{M}_\mathbb{O}$ is regularly varying if there exists an increasing sequence $\{c_n\}_{n \geq 1}$ of positive numbers which is regularly varying and a nonzero $\mu \in \mathbb{M}_\mathbb{O}$ such that $c_n \nu_n \rightarrow \mu$ in $\mathbb{M}_\mathbb{O}$ as $n \rightarrow \infty$.

The choice of terminology is motivated by the fact that $\{\nu_n(A)\}_{n \geq 1}$ is a regularly varying sequence for each set $A \in \mathcal{S}_\mathbb{O}$ bounded away from \mathbb{C} , $\mu(\partial A) = 0$ and $\mu(A) > 0$. We will now define regular variation for a single measure in $\mathbb{M}_\mathbb{O}$.

Definition 3.2. A measure $\nu \in \mathbb{M}_\mathbb{O}$ is regularly varying if the sequence

$$\{\nu(n \cdot), n \geq 1\}$$

in $\mathbb{M}_\mathbb{O}$ is regularly varying.

There are many equivalent formulations of regular variation for a measure $\nu \in \mathbb{M}_0$. Some are natural for statistical inference. Consider the following statements.

- (i) There exist a nonzero $\mu \in \mathbb{M}_0$ and a regularly varying sequence $\{c_n\}_{n \geq 1}$ of positive numbers such that $c_n \nu(n \cdot) \rightarrow \mu(\cdot)$ in \mathbb{M}_0 as $n \rightarrow \infty$.
- (ii) There exist a nonzero $\mu \in \mathbb{M}_0$ and a regularly varying function c such that $c(t) \nu(t \cdot) \rightarrow \mu(\cdot)$ in \mathbb{M}_0 as $t \rightarrow \infty$.
- (iii) There exist a nonzero $\mu \in \mathbb{M}_0$ and a set $E \in \mathcal{S}_0$ bounded away from \mathbb{C} such that $\nu(tE)^{-1} \nu(t \cdot) \rightarrow \mu(\cdot)$ in \mathbb{M}_0 as $t \rightarrow \infty$.
- (iv) There exist a nonzero $\mu \in \mathbb{M}_0$ and an increasing sequence $\{b_n\}_{n \geq 1}$ of positive numbers such that $n \nu(b_n \cdot) \rightarrow \mu(\cdot)$ in \mathbb{M}_0 as $n \rightarrow \infty$.
- (v) There exist a nonzero $\mu \in \mathbb{M}_0$ and an increasing positive function b such that $t \nu(b(t) \cdot) \rightarrow \mu(\cdot)$ in \mathbb{M}_0 as $t \rightarrow \infty$.

Theorem 3.1. *The statements (i)-(v) are equivalent and each statement implies that the limit measure μ has the homogeneity property*

$$\mu(\lambda A) = \lambda^{-\alpha} \mu(A) \quad (3.1)$$

for some $\alpha \geq 0$ and all $A \in \mathcal{S}_0$ and $\lambda > 0$.

The proof follows in Section 3.4.2 below.

Notice that a regularly varying measure does not correspond to a single scaling parameter α unless the multiplication operation with scalars is fixed.

3.3. More examples

We amplify the discussion of Section 3.1.1.

3.3.1. Continuation of Section 3.1.1

Example 3.1. Consider again the context of Section 3.1.1, item 1 where $\mathbb{S} = \mathbb{R}^2$ and let $\mathbb{C} = (\{0\} \times \mathbb{R}) \cup (\mathbb{R} \times \{0\})$. Consider two independent Pareto random variables: Let X_1 be $\text{Pa}(\gamma_1)$ and X_2 be $\text{Pa}(\gamma_2)$. Define $(\lambda, (x_1, x_2)) \mapsto (\lambda^{1/\gamma_1} x_1, \lambda^{1/\gamma_2} x_2)$. For $a, b > 0$

$$\begin{aligned} t^2 P[t^{-1}(X_1, X_2) \in (a, \infty) \times (b, \infty)] &= tP[X_1 > t^{1/\gamma_1} a] tP[X_2 > t^{1/\gamma_2} b] \\ &= a^{-\gamma_1} b^{-\gamma_2}. \end{aligned}$$

According to our definition, the distribution of (X_1, X_2) is regularly varying on $\mathbb{S} \setminus \mathbb{C}$. The limit measure therefore has the scaling property that for $\lambda > 0$,

$$\begin{aligned} \mu(\lambda[(a, \infty) \times (b, \infty)]) &= \mu((\lambda^{1/\gamma_1} a, \infty) \times (\lambda^{1/\gamma_2} b, \infty)) \\ &= \lambda^{-2} a^{-\gamma_1} b^{-\gamma_2} = \lambda^{-2} \mu((a, \infty) \times (b, \infty)). \end{aligned}$$

Example 3.2. Recall Section 3.1.1, item 2 where $\mathbb{S} = \mathbb{R}^2$ and $\mathbb{C} = \mathbb{R} \times \{0\}$ with $(\lambda, (x_1, x_2)) \mapsto (x_1, \lambda x_2)$. Suppose X_1, X_2 are independent with X_1 being $N(0, 1)$ and X_2 being $\text{Pa}(\gamma)$. For $a, b > 0$

$$t^\gamma P[t^{-1}(X_1, X_2) \in (a, \infty) \times (b, \infty)] = P[X_1 > a]t^\gamma P[X_2 > tb] = (1 - \Phi(a))b^{-\gamma},$$

implying that the distribution of (X_1, X_2) is regularly varying. For $\lambda > 0$, the limit measure has the scaling property,

$$\begin{aligned} \mu(\lambda[(a, \infty) \times (b, \infty)]) &= \mu((a, \infty) \times (\lambda b, \infty)) \\ &= \lambda^{-\gamma}(1 - \Phi(a))b^{-\gamma} = \lambda^{-\gamma}\mu((a, \infty) \times (b, \infty)). \end{aligned}$$

3.3.2. Examples using the mapping Theorem 2.3

Example 3.3 (Cf. [31], Theorem 2.1, page 677). Suppose $\mathbb{S} = [0, \infty)^2$, $\mathbb{C} = [0, \infty) \times \{0\}$ so that

$$\mathbb{O} = \mathbb{S} \setminus \mathbb{C} = [0, \infty)^2 \setminus [0, \infty) \times \{0\} = [0, \infty) \times (0, \infty) =: \mathbb{D}_\square.$$

Define $h : \mathbb{D}_\square \mapsto \mathbb{D}_\square$ by $h(x, y) = (xy, y)$. If $A' \subset \mathbb{D}_\square$ and $d(A', [0, \infty) \times \{0\}) > 0$, then $\inf\{y : (x, y) \in A' \text{ for some } x\} > 0$ and so $h^{-1}(A') = \{(x, y) : h(x, y) \in A'\} = \{(x, y) : (xy, y) \in A'\}$ is also at positive distance from $[0, \infty) \times \{0\}$. So the hypotheses of Theorem 2.3 are satisfied with h and if $\mu_t \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}_\square}$ then it follows that $\mu_t \circ h^{-1} \rightarrow \mu \circ h^{-1}$ in $\mathbb{M}_{\mathbb{D}_\square}$. In particular, suppose for a random vector (X, Y) and scaling function $b(t) \rightarrow \infty$,

$$tP\left[\left(X, \frac{Y}{b(t)}\right) \in \cdot\right] \rightarrow \mu(\cdot) \quad \text{in } \mathbb{M}_{\mathbb{D}_\square}. \tag{3.2}$$

This is regular variation of the distribution of (X, Y) on \mathbb{D}_\square with the scaling function defined as $(\lambda, (x, y)) \mapsto (x, \lambda y)$. The mapping Theorem 2.3 gives

$$tP\left[\left(\frac{XY}{b(t)}, \frac{Y}{b(t)}\right) \in \cdot\right] \rightarrow \mu'(\cdot) \quad \text{in } \mathbb{M}_{\mathbb{D}_\square}, \tag{3.3}$$

where $\mu' = \mu \circ h^{-1}$, which is regular variation with respect to the traditional scaling $(\lambda, (x, y)) \mapsto (\lambda x, \lambda y)$.

Conversely, define $g : \mathbb{D}_\square \mapsto \mathbb{D}_\square$ by $g(x, y) = (x/y, y)$. One observes g is continuous and obeys the bounded away condition and so $\mu_t \rightarrow \mu$ in $\mathbb{M}_{\mathbb{D}_\square}$ implies $\mu_t \circ h^{-1} \rightarrow \mu \circ h^{-1}$ in $\mathbb{M}_{\mathbb{D}_\square}$.

The summary is that (3.2) and (3.3) are equivalent.

Example 3.4 (Polar coordinates). Set

$$\mathbb{S} = [0, \infty)^2, \mathbb{C} = \{0\}, \mathbb{O} = [0, \infty)^2 \setminus \{0\}$$

with scaling function $(\lambda, x) = (\lambda, (x_1, x_2)) \mapsto \lambda x = (\lambda x_1, \lambda x_2)$. For some choice of norm $x \mapsto \|x\|$ define $\mathbb{N} = \{x \in \mathbb{S} : \|x\| = 1\}$. Also define

$$\mathbb{S}' = [0, \infty) \times \mathbb{N}, \mathbb{C}' = \{0\} \times \mathbb{N}, \mathbb{O}' = (0, \infty) \times \mathbb{N},$$

and scaling operation on \mathbb{O}' is $(\lambda, (r, a)) \mapsto (\lambda r, a)$. The map

$$h(x) = (\|x\|, x/\|x\|)$$

from $\mathbb{O} \mapsto \mathbb{O}'$ is continuous. Let d and d' be the metrics on \mathbb{S} and \mathbb{S}' . Suppose X has a regularly varying distribution on \mathbb{O} so that for some $b(t) \rightarrow \infty$,

$$tP[X/b(t) \in \cdot] \rightarrow \mu(\cdot) \tag{3.4}$$

in $\mathbb{M}_{\mathbb{O}}$ for some limit measure μ . We show $h(X) =: (R, \Theta)$ has a regularly varying distribution on \mathbb{O}' . We apply Theorem 2.3 so suppose $A' \subset \mathbb{O}'$ satisfies $d'(A', \{0\} \times \mathbb{N}) > 0$, that is, A' is bounded away from the deleted portion of \mathbb{S}' . Then $\inf\{r > 0 : (r, a) \in A'\} = \delta > 0$ and $h^{-1}(A') = \{x \in \mathbb{O} : (\|x\|, x/\|x\|) \in A'\}$ satisfies $\inf\{\|x\| : x \in h^{-1}(A')\} = \delta' > 0$. So the hypotheses of Theorem 2.3 are satisfied and allow the conclusion that

$$tP\left[\left(\frac{R}{b(t)}, \Theta\right) \in \cdot\right] \rightarrow \mu \circ h^{-1}(\cdot) \quad \text{in } \mathbb{M}_{\mathbb{O}'}. \tag{3.5}$$

Conversely, given regular variation on \mathbb{O}' as in (3.5), define $g : \mathbb{O}' \mapsto \mathbb{O}$ by $g(r, a) = ra$. Mimic the verification above to conclude (3.5) implies (3.4).

Example 3.5. Examples 3.3 and 3.4 typify the following paradigm. Consider two triples $(\mathbb{S}, \mathbb{C}, \mathbb{O})$ and $(\mathbb{S}', \mathbb{C}', \mathbb{O}')$, and a homeomorphism $h : \mathbb{O} \rightarrow \mathbb{O}'$ with the property that $h^{-1}(A')$ is bounded away from \mathbb{C} if A' is bounded away from \mathbb{C}' . The multiplication by a scalar $(\lambda, x) \mapsto \lambda x$ in \mathbb{O} gives rise to the multiplication by a scalar $(\lambda', x') \mapsto \lambda' x' := h(\lambda' h^{-1}(x'))$ in \mathbb{O}' . Notice that $(\lambda', x') \mapsto \lambda' x'$ is continuous, $1x' = x'$, and $\lambda'_1(\lambda'_2(x')) = (\lambda'_1 \lambda'_2)x'$. We also need to check that $d'(\lambda' x', \mathbb{C}') > d'(x', \mathbb{C}')$ if $\lambda' > 1$.

3.4. Proofs

3.4.1. Preliminaries

For $A \in \mathcal{S}_{\mathbb{O}}$, write $S(A) = \{\lambda x : x \in A, \lambda \geq 1\}$.

Lemma 3.1. *Let $\mu \in \mathbb{M}_{\mathbb{O}}$ be nonzero. There exists $x \in \mathbb{O}$ and $\delta > 0$ such that $S(B_{x,\delta})$ is bounded away from \mathbb{C} , $\mu(S(B_{x,\delta})) > 0$, and $\mu(\partial r S(B_{x,\delta})) = 0$ for $r \geq 1$ in some set of positive measure containing 1.*

Proof. The first two properties are obvious. In order to prove the final claim, set $\gamma(r) = d(\partial r S(B_{x,\delta}), \mathbb{C})$. Notice that $\gamma(r) = d(\partial r B_{x,\delta}, \mathbb{C})$ and that $\partial r S(B_{x,\delta}) \subset \mathbb{O} \setminus \mathbb{C}^{\gamma(r)}$. Choose $x \in \mathbb{O}$ and $\delta > 0$ such that $\mu(\partial S(B_{x,\delta})) = 0$ and $\mu(\partial(\mathbb{O} \setminus \mathbb{C}^{\gamma(1)})) = 0$, and such that $\gamma(r') > \gamma(1)$ for some $r' > 1$. The existence of such x and δ follows from Lemma 2.5. Lemma 2.5 also implies that $\mu(\partial(\mathbb{O} \setminus \mathbb{C}^{\gamma})) = 0$ for all but at most countably many $\gamma \in [\gamma(1), \gamma(r')]$. Since $\gamma(r)$ is a nondecreasing and continuous function and $\gamma(r') > \gamma(1)$, there exists a set $R \subset [1, r']$ of positive Lebesgue measure, with $1 \in R$, such that $\mu(\partial(\mathbb{O} \setminus \mathbb{C}^{\gamma(r)})) = 0$ for $r \in R$. \square

Given $\mu \in \mathbb{M}_\mathbb{O}$, let \mathcal{A}_μ denote the set of μ -continuity sets $A \in \mathcal{S}_\mathbb{O}$ bounded away from \mathbb{C} satisfying $S(A) = A$.

Lemma 3.2. *If $\mu_n(A) \rightarrow \mu(A)$ for all $A \in \mathcal{A}_\mu$, then $\mu_n \rightarrow \mu$ in $\mathbb{M}_\mathbb{O}$.*

Proof. Let \mathcal{D}_μ denote the π -system of finite differences of sets of the form $A_1 \setminus A_2$ for $A_1, A_2 \in \mathcal{A}_\mu$ with $A_2 \subset A_1$. Take $x \in \mathbb{O}$ and $\epsilon > 0$ such that $B_{x,\epsilon}$ is bounded away from \mathbb{C} . The sets $\partial S(B_{x,r})$, for $r \in (0, \epsilon)$, are disjoint. Similarly, the sets $\partial B_{x,r}$, for $r \in (0, \epsilon)$, are disjoint. Therefore, $\mu(\partial S(B_{x,r})) = \mu(\partial B_{x,r}) = 0$ for all but at most countably many $r \in (0, \epsilon)$. Moreover, $B_{x,r} = S(B_{x,r}) \setminus (S(B_{x,r}) \setminus B_{x,r})$, so $B_{x,r} \in \mathcal{D}_\mu$ for all but at most countably many $r \in (0, \epsilon)$. Hence, there exists $A \in \mathcal{D}_\mu$ such that $x \in A^\circ \subset A \subset B_{x,\epsilon}$. Moreover, for any x in an open set G bounded away from \mathbb{C} , there exists $A \in \mathcal{D}_\mu$ such that $x \in A^\circ \subset A \subset G$. Since \mathbb{O} is separable we find (as in the proof of Theorem 2.3 in [6]) that there is a countable subcollection $\{A_{x_i}^\circ\}$ of $\{A_x^\circ : x \in G\}$, $A_x^\circ \in \mathcal{D}_\mu$, that covers G and that $G = \cup_i A_{x_i}^\circ$. The inclusion-exclusion argument in the proof of Theorem 2.2 in [6] implies that $\liminf_n \mu_n(G) \geq \mu(G)$ for all open sets G bounded away from \mathbb{C} . Any closed F bounded away from \mathbb{C} is a subset of an open μ -continuity set $A = \mathbb{O} \setminus C^r$ for some $r > 0$. Notice that $A \in \mathcal{A}_\mu$. Therefore

$$\mu(A) - \limsup_n \mu_n(F) = \liminf_n \mu_n(A \setminus F) \geq \mu(A \setminus F) = \mu(A) - \mu(F),$$

i.e. $\limsup_n \mu_n(F) \leq \mu(F)$. The conclusion follows from Theorem 2.1. \square

3.4.2. Proof of Theorem 3.1

The proof is structured as follows. We first prove that (iii) implies the homogeneity property in (3.1) of the limit measure μ . Then we prove that the statements (i)–(v) are equivalent and that the limit measures are the same up to a constant factor.

Suppose that (iii) holds and take $E' = S(B_{x,\delta})$ satisfying the conditions in Lemma 3.1. Then, for $\lambda \geq 1$ in a set of positive measure containing 1,

$$\frac{\nu(t\lambda E')}{\nu(tE')} = \frac{\nu(t\lambda E')}{\nu(tE)} \frac{\nu(tE)}{\nu(tE')} \rightarrow \frac{\mu(\lambda E')}{\mu(E')} \in (0, \infty),$$

as $t \rightarrow \infty$. It follows from Theorem 1.4.1 in [8] that $t \mapsto \nu(tE')$ is regularly varying and that $\mu(\lambda E') = \lambda^{-\alpha} \mu(E')$ for some $\alpha \in \mathbb{R}$ and all $\lambda > 0$. Property (A3) implies that $\nu(t\lambda E')/\nu(tE') \leq 1$ for $\lambda \geq 1$, so $\alpha \geq 0$. Moreover, $\mu(\partial(\lambda E')) = 0$ for all $\lambda > 0$ and $\nu(tE')^{-1} \nu(t \cdot) \rightarrow \mu(E')^{-1} \mu(\cdot)$ in $\mathbb{M}_\mathbb{O}$ as $t \rightarrow \infty$. In particular, if $A \in \mathcal{A}_\mu$, then for any $\lambda > 0$,

$$\frac{\nu(t\lambda A)}{\nu(tE')} = \frac{\nu(t\lambda A)}{\nu(t\lambda E')} \frac{\nu(t\lambda E')}{\nu(tE')} \rightarrow \lambda^{-\alpha} \frac{\mu(A)}{\mu(E')},$$

as $t \rightarrow \infty$. Hence, for $A \in \mathcal{A}_\mu$ $\mu(\lambda A) = \lambda^{-\alpha} \mu(A)$ for all $\lambda > 0$. Lemma 3.2 implies that $\mu(\lambda A) = \lambda^{-\alpha} \mu(A)$ for all $A \in \mathcal{S}_\mathbb{O}$ and $\lambda > 0$.

Suppose that (i) holds and set $c(t) = c_{[t]}$. For each $A \in \mathcal{A}_\mu$ and $t \geq 1$ it holds that

$$\frac{c_{[t]}}{c_{[t+1]}} c_{[t+1]} \nu(([t] + 1)A) \leq c(t) \nu(tA) \leq c_{[t]} \nu([t]A). \tag{3.6}$$

Since $\{c_n\}_{n \geq 1}$ is regularly varying it holds that $\lim_{n \rightarrow \infty} c_n/c_{n+1} = 1$. Hence, $\lim_{t \rightarrow \infty} c(t) \nu(tA) = \mu(A)$ for all $A \in \mathcal{A}_\mu$. It follows from Lemma 3.2 that (ii) holds.

Suppose that (ii) holds. Then $c_{[t]} \nu([t] \cdot) \rightarrow \mu(\cdot)$ in $\mathbb{M}_\mathbb{O}$. Moreover, $\{c_{[t]}\}$ is a regularly varying sequence since $c(t)$ is a regularly varying function. Therefore, (ii) implies (i).

Suppose that (ii) holds. Take a set $E \in \mathcal{S}_\mathbb{O}$ bounded away from \mathbb{C} such that $\nu(tE), \mu(E) > 0$ and $\mu(\partial E) = 0$. Then

$$\frac{\nu(t \cdot)}{\nu(tE)} = \frac{c(t) \nu(t \cdot)}{c(t) \nu(tE)} \rightarrow \frac{\mu(\cdot)}{\mu(E)} \quad \text{in } \mathbb{M}_\mathbb{O}$$

as $t \rightarrow \infty$. Hence, (iii) holds.

Suppose that (iii) holds. It was already proved that statement (iii) implies that $t \mapsto \nu(tE)$ is regularly varying with index $-\alpha \leq 0$. Setting $c(t) = 1/\nu(tE)$ implies that $c(t)$ is regularly varying with index α and that $c(t) \nu(t \cdot) \rightarrow \mu(\cdot)$ in $\mathbb{M}_\mathbb{O}$. This proves that (iii) implies (ii). Up to this point we have proved that statements (i)–(iii) are equivalent.

Suppose that (iv) holds. Set $b(t) = b_{[t]}$ and take $A \in \mathcal{A}_\mu$. Then

$$\frac{[t]}{[t+1]} [t+1] \nu(b_{[t+1]}A) \leq t \nu(b(t)A) \leq \frac{[t+1]}{[t]} [t] \nu(b_{[t]}A),$$

from which it follows that $\lim_{t \rightarrow \infty} t \nu(b(t)A) = \mu(A)$. Lemma 3.2 implies that (v) holds. If (v) holds, then it follows immediately that also (iv) holds. Hence, statements (iv) and (v) are equivalent.

Suppose that (iv) holds. Take E such that $\mu(\partial E) = 0$ and $\mu(E) > 0$. For $t > b_1$, let $k = k(t)$ be the largest integer with $b_k \leq t$. Then $b_k \leq t < b_{k+1}$ and $k \rightarrow \infty$ as $t \rightarrow \infty$. Hence, for $A \in \mathcal{A}_\mu$,

$$\frac{k}{k+1} \frac{(k+1) \nu(b_{k+1}A)}{k \nu(b_k E)} \leq \frac{\nu(tA)}{\nu(tE)} \leq \frac{k+1}{k} \frac{k \nu(b_k A)}{(k+1) \nu(b_{k+1} E)},$$

from which it follows that $\lim_{t \rightarrow \infty} \nu(tA)/\nu(tE) = \mu(A)/\mu(E)$. It follows from Lemma 3.2 that (iii) holds. Hence, each of the statements (iv) and (v) implies each of the statements (i)–(iii).

Suppose that (iii) holds. Then $c(t) := 1/\nu(tE)$ is regularly varying at infinity with index $\alpha \geq 0$. If $\alpha > 0$, then $c(c^{-1}(t)) \sim t$ as $t \rightarrow \infty$ by Proposition B.1.9 (10) in [17] and therefore

$$\lim_{t \rightarrow \infty} t \nu(c^{-1}(t)A) = \lim_{t \rightarrow \infty} c(c^{-1}(t)) \nu(c^{-1}(t)A) = \mu(A)$$

for all $A \in \mathcal{S}_0$ bounded away from \mathbb{C} with $\mu(\partial A) = 0$. If $\alpha = 0$, then Proposition 1.3.4 in [8] says that there exists a continuous and increasing function \tilde{c} such that $\tilde{c}(t) \sim c(t)$ as $t \rightarrow \infty$. In particular, $\tilde{c}(\tilde{c}^{-1}(t)) = t$ and

$$t\nu(\tilde{c}^{-1} \cdot) = \tilde{c}(\tilde{c}^{-1}(t))\nu(\tilde{c}^{-1} \cdot) \rightarrow \mu(\cdot)$$

in \mathbb{M}_0 as $t \rightarrow \infty$. Hence, (v) holds. □

4. \mathbb{R}_+^∞ and \mathbb{R}_+^p

This section considers regular variation for measures on the metric spaces \mathbb{R}_+^∞ and \mathbb{R}_+^p for $p \geq 1$ and applies the theory of Sections 2 and 3. We begin in Section 4.1 with notation and specification of metrics and then address in Section 4.2 continuity properties for the following maps

- CUMSUM : $(x_1, x_2, \dots) \mapsto (x_1, x_1 + x_2, x_1 + x_2 + x_3, \dots)$.
- PROJ_p : $(x_1, x_2, \dots) \mapsto (x_1, \dots, x_p)$.
- POLAR : $(x_1, \dots, x_p) \mapsto (\|x\|, (x_1, \dots, x_p)/\|x\|)$, where $x = (x_1, \dots, x_p)$ and the norm is Euclidean norm on \mathbb{R}_+^p . We also define the generalized polar coordinate transformation in (4.3) which is necessary for estimating the tail measure of regular variation when $\{x \in \mathbb{R}_+^p : \|x\| = 1\}$, the Euclidean unit sphere, is not bounded away from \mathbb{C} in the space $\mathbb{R}_+^\infty \setminus \mathbb{C}$.

Section 4.3 reduces the convergence question to finite dimensions by giving criteria for reduction of convergence of measures in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$ to convergence of projected measures in $\mathbb{M}(\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}))$. Section 4.4 returns to the comparison of vague convergence with \mathbb{M} -convergence initiated in Section 2.3 with the goal of using existing results based on regular variation in a compactified version $\mathbb{R}_+^p \setminus \{0\}$ in our present context, rather than proving things from scratch. Section 4 concludes with Section 4.5, a discussion of regular variation of measures on $\mathbb{R}_+^\infty \setminus \mathbb{C}$ giving particular attention to hidden-regular-variation properties of the distribution of $X = (X_1, X_2, \dots)$, a sequence of iid non-negative random variables whose marginal distributions have regularly varying tails. This discussion extends naturally to hidden-regular-variation properties of an infinite sequence of non-negative decreasing Poisson points whose mean measure has a regularly varying tail. Results for the Poisson sequence provide the basis of our approach in Section 5 to regular variation of the distribution of a Lévy process whose Lévy measure is regularly varying.

4.1. Preliminaries

We write $x \in \mathbb{R}_+^\infty$ as $x = (x_1, x_2, \dots)$. For $p \geq 1$, the projection into \mathbb{R}_+^p is written as $\text{PROJ}_p(x) = x|_p = (x_1, \dots, x_p)$. To avoid confusion, we sometimes write $0_\infty = (0, 0, \dots) \in \mathbb{R}_+^\infty$ and $0_p = (0, \dots, 0) \in \mathbb{R}_+^p$, the vector of p zeros. We also augment a vector in \mathbb{R}_+^p to get a sequence in \mathbb{R}_+^∞ and write, for instance,

$$(x_1, \dots, x_p, 0_\infty) = (x_1, \dots, x_p, 0, 0, \dots).$$

4.1.1. Metrics

All metrics are equivalent on \mathbb{R}_+^p . The usual metric on \mathbb{R}_+^∞ is

$$d_\infty(x, y) = \sum_{i=1}^{\infty} \frac{|x_i - y_i| \wedge 1}{2^i},$$

and we also need

$$d'_\infty(x, y) = \sum_{p=1}^{\infty} \frac{(\sum_{l=1}^p |x_l - y_l|) \wedge 1}{2^p} = \sum_{p=1}^{\infty} \frac{\|x|_p - y|_p\|_1 \wedge 1}{2^p},$$

where $\|\cdot\|_1$ is the usual L_1 norm on \mathbb{R}_+^p .

Proposition 4.1. *The metrics d_∞ and d'_∞ are equivalent on \mathbb{R}_+^∞ and*

$$d_\infty(x, y) \leq d'_\infty(x, y) \leq 2d_\infty(x, y).$$

Proof. First of all,

$$d'_\infty(x, y) = \sum_{i=1}^{\infty} \frac{(\sum_{l=1}^i |x_l - y_l|) \wedge 1}{2^i} \geq \sum_{i=1}^{\infty} \frac{|x_i - y_i| \wedge 1}{2^i} = d_\infty(x, y).$$

Furthermore, observe that

$$\begin{aligned} d'_\infty(x, y) &= \sum_{i=1}^{\infty} \frac{(\sum_{l=1}^i |x_l - y_l|) \wedge 1}{2^i} \leq \sum_{i=1}^{\infty} \frac{\sum_{l=1}^i (|x_l - y_l| \wedge 1)}{2^i} \\ &= \sum_{l=1}^{\infty} \sum_{i=l}^{\infty} 2^{-i} (|x_l - y_l| \wedge 1) = \sum_{l=1}^{\infty} 2 \cdot 2^{-l} (|x_l - y_l| \wedge 1) = 2d_\infty(x, y), \end{aligned}$$

which proves the other inequality. \square

4.2. Continuity of maps

With a view toward applying Corollary 2.1, we consider the continuity of several maps.

4.2.1. Cumsum

We begin with the map $\text{CUMSUM} : \mathbb{R}_+^\infty \mapsto \mathbb{R}_+^\infty$ defined by

$$\text{CUMSUM}(x) = (x_1, x_1 + x_2, x_1 + x_2 + x_3, \dots).$$

Proposition 4.2. *The map $\text{CUMSUM} : \mathbb{R}_+^\infty \mapsto \mathbb{R}_+^\infty$ is uniformly continuous and, in fact, is Lipschitz in the d_∞ metric.*

Proof. We write

$$\begin{aligned} d_\infty(\text{CUMSUM}(x), \text{CUMSUM}(y)) &= \sum_{i=1}^{\infty} \frac{|\sum_{l=1}^i x_l - \sum_{l=1}^i y_l| \wedge 1}{2^i} \\ &\leq \sum_{i=1}^{\infty} \frac{(\sum_{l=1}^i |x_l - y_l|) \wedge 1}{2^i} = d'_\infty(x, y) \leq 2d_\infty(x, y). \end{aligned}$$

□

We can now apply Corollary 2.1.

Corollary 4.1. *Let $\mathbb{S} = \mathbb{S}' = \mathbb{R}_+^\infty$ and suppose both \mathbb{C} and $\text{CUMSUM}(\mathbb{C})$ are closed in \mathbb{R}_+^∞ . If for $n \geq 0$, $\mu_n \in \mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$ and $\mu_n \rightarrow \mu_0$ in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$, then $\mu_n \circ \text{CUMSUM}^{-1} \rightarrow \mu_0 \circ \text{CUMSUM}^{-1}$ in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \text{CUMSUM}(\mathbb{C}))$.*

For example, if $\mathbb{C} = \{0_\infty\}$, then $\text{CUMSUM}(\mathbb{C}) = \{0_\infty\}$. For additional examples, see (4.6).

4.2.2. Projection

For $p \geq 1$, recall $\text{PROJ}_p(x) = x_{|p} = (x_1, \dots, x_p)$ from $\mathbb{R}_+^\infty \mapsto \mathbb{R}_+^p$.

Proposition 4.3. *$\text{PROJ}_p : \mathbb{R}_+^\infty \mapsto \mathbb{R}_+^p$ is uniformly continuous.*

Proof. Let $d_p(x_{|p}, y_{|p}) = \sum_{i=1}^p |x_i - y_i|$ be the usual L_1 metric. Given $0 < \epsilon < 1$, we must find $\delta > 0$ such that $d_\infty(x, y) < \delta$ implies $d_p(x_{|p}, y_{|p}) < \epsilon$. We try $\delta = 2^{-p}\epsilon$. Then

$$\begin{aligned} \delta = 2^{-p}\epsilon > d_\infty(x, y) &= \sum_{i=1}^{\infty} \frac{|x_i - y_i| \wedge 1}{2^i} \geq \sum_{i=1}^p \frac{|x_i - y_i| \wedge 1}{2^i} \\ &\geq 2^{-p} \sum_{i=1}^p |x_i - y_i| \wedge 1. \end{aligned}$$

Therefore $\epsilon \geq \sum_{i=1}^p |x_i - y_i| \wedge 1$, so that $\epsilon \geq \sum_{i=1}^p |x_i - y_i| = d_p(x_{|p}, y_{|p})$. □

Apply Corollary 2.1:

Corollary 4.2. *Let $\mathbb{S} = \mathbb{S}' = \mathbb{R}_+^\infty$ and suppose \mathbb{C} is closed in \mathbb{R}_+^∞ and $\text{PROJ}_p(\mathbb{C})$ is closed in \mathbb{R}_+^p . If for $n \geq 0$, $\mu_n \in \mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$ and $\mu_n \rightarrow \mu_0$ in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$, then $\mu_n \circ \text{PROJ}_p^{-1} \rightarrow \mu_0 \circ \text{PROJ}_p^{-1}$ in $\mathbb{M}(\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}))$.*

4.2.3. Polar coordinate transformations

The polar coordinate transformation in \mathbb{R}_+^p is heavily relied upon when making inferences about the limit measure of regular variation. Transforming from Cartesian to polar coordinates disintegrates the transformed limit measure into a product measure, one of whose factors concentrates on the unit sphere. This

factor is called the angular measure. Estimating the angular measure and then transforming back to Cartesian coordinates provides the most reliable inference technique for tail probability estimation in \mathbb{R}_+^p using heavy-tail asymptotics. See [40, pages 173ff, 313]. When removing more than $\{0\}$ from \mathbb{R}_+^p , the unit sphere may no longer be bounded away from what is removed and an alternative technique we call the *generalized polar coordinate transformation* can be used.

We continue the discussion of Example 3.4 that relied on the mapping Theorem 2.3. Here we rely on Corollary 2.2. Pick a norm $\|\cdot\|$ on \mathbb{R}_+^p and define $\mathfrak{N} = \{x \in \mathbb{R}_+^p : \|x\| = 1\}$. The conventional polar coordinate transform $\text{POLAR} : \mathbb{R}_+^p \setminus \{0_p\} \mapsto (0, \infty) \times \mathfrak{N}$ is traditionally defined as

$$\text{POLAR}(x) = (\|x\|, x/\|x\|). \quad (4.1)$$

Compared with the notation of Corollary 2.2, we have $\mathbb{S} = \mathbb{R}_+^p$, $\mathbb{C} = \{0_p\}$ which is compact in \mathbb{S} , $\mathbb{S}' = [0, \infty) \times \mathfrak{N}$, $\mathbb{C}' = \{0\} \times \mathfrak{N}$ which is closed in \mathbb{S}' . Since POLAR is continuous on the domain, we get from Corollary 2.2 the following.

Corollary 4.3. *Suppose $\mu_n \rightarrow \mu_0$ in $\mathbb{M}(\mathbb{R}_+^p \setminus \{0_p\})$. Then*

$$\mu_n \circ \text{POLAR}^{-1} \rightarrow \mu_0 \circ \text{POLAR}^{-1}$$

in $\mathbb{M}((0, \infty) \times \mathfrak{N})$.

When removing more from the state space than just $\{0_p\}$, the conventional polar coordinate transform (4.1) is not useful if \mathfrak{N} is not compact, or at least bounded away from what is removed. For example, if $\mathbb{S} \setminus \mathbb{C} = (0, \infty)^p$, \mathfrak{N} is not compact nor bounded away from the removed axes. The following generalization [15] sometimes resolves this, provided (4.2) below holds.

Temporarily, we proceed generally and assume \mathbb{S} is a complete, separable metric space and that scalar multiplication is defined. If \mathbb{C} is a cone, $\theta\mathbb{C} = \mathbb{C}$ for $\theta \geq 0$. Suppose further that the metric on \mathbb{S} satisfies

$$d(\theta x, \theta y) = \theta d(x, y), \quad \theta \geq 0, (x, y) \in \mathbb{S} \times \mathbb{S}. \quad (4.2)$$

Note (4.2) holds for a Banach space where distance is defined by a norm. (It does not hold for \mathbb{R}_+^∞ .) If we intend to remove the closed cone \mathbb{C} , set

$$\mathfrak{N}_{\mathbb{C}} = \{s \in \mathbb{S} \setminus \mathbb{C} : d(s, \mathbb{C}) = 1\},$$

which plays the role of the unit sphere and $\mathbb{C}' = \{0\} \times \mathfrak{N}_{\mathbb{C}}$ is closed. Define the generalized polar coordinate transformation

$$\text{GPOLAR} : \mathbb{S} \setminus \mathbb{C} \mapsto (0, \infty) \times \mathfrak{N}_{\mathbb{C}} = [0, \infty) \times \mathfrak{N}_{\mathbb{C}} \setminus (\{0\} \times \mathfrak{N}_{\mathbb{C}}) = \mathbb{S}' \setminus \mathbb{C}'$$

by

$$\text{GPOLAR}(s) = (d(s, \mathbb{C}), s/d(s, \mathbb{C})), \quad s \in \mathbb{S} \setminus \mathbb{C}. \quad (4.3)$$

Since \mathbb{C} is a cone and $d(\cdot, \cdot)$ has property (4.2), we have for any $s \in \mathbb{S} \setminus \mathbb{C}$ that

$$d\left(\frac{s}{d(s, \mathbb{C})}, \mathbb{C}\right) = d\left(\frac{s}{d(s, \mathbb{C})}, \frac{1}{d(s, \mathbb{C})}\mathbb{C}\right) = \frac{1}{d(s, \mathbb{C})}d(s, \mathbb{C}) = 1,$$

so the second coordinate of GPOLAR belongs to $\aleph_{\mathbb{C}}$. For example, if $\mathbb{S} = \mathbb{R}_+^2$ and we remove the cone consisting of the axes through 0_2 , that is, $\mathbb{C} = \{0\} \times [0, \infty) \cup [0, \infty) \times \{0\}$, then $\aleph_{\mathbb{C}} = \{x \in \mathbb{R}_+^2 : x_1 \wedge x_2 = 1\}$. The inverse map $\text{GPOLAR}^{-1} : (0, \infty) \times \aleph_{\mathbb{C}} \mapsto \mathbb{S} \setminus \mathbb{C}$ is

$$\text{GPOLAR}^{-1}(r, a) = ra, \quad r \in (0, \infty), a \in \aleph_{\mathbb{C}}.$$

It is relatively easy to check that if $A' \subset (0, \infty) \times \aleph_{\mathbb{C}}$ is bounded away from $\mathbb{C}' = \{0\} \times \aleph_{\mathbb{C}}$, then $\text{GPOLAR}^{-1}(A')$ is bounded away from \mathbb{C} . On $(0, \infty) \times \aleph_{\mathbb{C}}$ adopt the metric

$$d'((r_1, a_1), (r_2, a_2)) = |r_1 - r_2| \vee d_{\aleph_{\mathbb{C}}}(a_1, a_2),$$

where $d_{\aleph_{\mathbb{C}}}(a_1, a_2)$ is an appropriate metric on $\aleph_{\mathbb{C}}$. Suppose $d'(A', \{0\} \times \aleph_{\mathbb{C}}) = \epsilon > 0$. This means

$$\epsilon = \inf_{\substack{(r_1, a_1) \in A' \\ a_2 \in \aleph_{\mathbb{C}}}} d'((r_1, a_1), (0, a_2))$$

and setting $a_2 = a_1$ this is $\inf_{(r_1, a_1) \in \mathbb{D}'} r_1$. We conclude that $(r, a) \in A'$ implies $r \geq \epsilon$. Since $\text{GPOLAR}^{-1}(A') = \{ra : (r, a) \in A'\}$, we have in $\mathbb{S} \setminus \mathbb{C}$, remembering that \mathbb{C} is assumed to be a cone,

$$\begin{aligned} d(\{ra : (r, a) \in A'\}, \mathbb{C}) &= \inf_{(r, a) \in A'} d(ra, \mathbb{C}) = \inf_{(r, a) \in A'} d(ra, r\mathbb{C}) \\ &= \inf_{(r, a) \in A'} rd(a, \mathbb{C}) \geq \epsilon \cdot 1. \end{aligned}$$

The last line uses (4.2) and the definition of $\aleph_{\mathbb{C}}$.

The hypotheses of Theorem 2.3 are verified so we get the following conclusion about GPOLAR.

Corollary 4.4. *Suppose \mathbb{S} is a complete, separable metric space such that (4.2) holds, scalar multiplication is defined and supposed \mathbb{C} is a closed cone. Then $\mu_n \rightarrow \mu_0$ in $\mathbb{M}(\mathbb{S} \setminus \mathbb{C})$ implies*

$$\mu_n \circ \text{GPOLAR}^{-1} \rightarrow \mu_0 \circ \text{GPOLAR}^{-1}$$

in $\mathbb{M}((0, \infty) \times \aleph_{\mathbb{C}})$. The converse holds as well.

The converse is proven in a similar way.

Remark on (4.2). As mentioned, (4.2) holds if the metric is defined by a norm on the space \mathbb{S} . Thus, (4.2) holds for a Banach space and on $\mathbb{R}_+^p, \mathbb{C}[0, 1]$ with sup-norm and $\mathbb{D}([0, 1], \mathbb{R})$ with Skorohod metric. It fails in \mathbb{R}_+^∞ .

4.3. Reducing \mathbb{R}_+^∞ convergence to finite-dimensional convergence

Corollary 4.2 shows when convergence in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$ implies convergence in $\mathbb{M}(\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}))$. Here is a circumstance where the converse is true and reduces the problem of convergence in infinite-dimensional space to finite dimensions.

Theorem 4.1. *Supppse for every $p \geq 1$, that the closed set $\mathbb{C} \subset \mathbb{R}_+^\infty$ satisfies $\text{PROJ}_p(\mathbb{C})$ is closed in \mathbb{R}_+^p and*

$$(z_1, \dots, z_p) \in \text{PROJ}_p(\mathbb{C}) \quad \text{implies} \quad (z_1, \dots, z_p, 0_\infty) \in \mathbb{C}. \tag{4.4}$$

Then $\mu_n \rightarrow \mu_0$ in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C})$ if and only if for all $p \geq 1$ such that $\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}) \neq \emptyset$ we have

$$\mu_n \circ \text{PROJ}_p^{-1} \rightarrow \mu_0 \circ \text{PROJ}_p^{-1} \tag{4.5}$$

in $\mathbb{M}(\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}))$.

Remark on condition (4.4). The condition says take an infinite sequence z in \mathbb{C} , truncate it to $z|_p \in \mathbb{R}_+^p$, and then make it infinite again by filling in zeros for all the components beyond the p th. The result must still be in \mathbb{C} . Examples:

1. $\mathbb{C} = \{0_\infty\}$.
2. Pick an integer $j \geq 1$ and define

$$\mathbb{C}_{\leq j} = \{x \in \mathbb{R}_+^\infty : \sum_{i=1}^{\infty} \epsilon_{x_i}(0, \infty) \leq j\}, \tag{4.6}$$

where recall $\epsilon_x(A) = 1$, if $x \in A$, and 0, if $x \in A^c$. So $\mathbb{C}_{\leq j}$ is the set of sequences with at most j positive components. Truncation and then insertion of zeros does not increase the number of positive components so $\mathbb{C}_{\leq j}$ is invariant under the operation implied by (4.4).

Proof. Suppose \mathbb{C} satisfies (4.4) and (4.5) holds. Suppose $f \in \mathcal{C}(\mathbb{R}_+^\infty \setminus \mathbb{C})$ and without loss of generality suppose f is uniformly continuous with modulus of continuity

$$\omega_f(\eta) = \sup_{\substack{(x,y) \in \mathbb{R}_+^\infty \setminus \mathbb{C} \\ d_\infty(x,y) < \eta}} |f(x) - f(y)|.$$

There exists $1 > \delta > 0$ such that $d_\infty(x, \mathbb{C}) < \delta$ implies $f(x) = 0$. Observe,

$$d_\infty((x|_p, 0_\infty), x) \leq \sum_{j=p+1}^{\infty} 2^{-j} = 2^{-p}. \tag{4.7}$$

Pick any p so large that $2^{-p} < \delta/2$ and define

$$g(x_1, \dots, x_p) = f(x_1, \dots, x_p, 0_\infty).$$

Then we have

(a) From (4.7),

$$|f(x) - g(x_{|p})| = |f(x) - f(x_{|p}, 0_\infty)| \leq \omega_f(2^{-p}).$$

(b) $g \in \mathcal{C}(\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}))$ and g is uniformly continuous.

To verify that the support of g is positive distance away from $\text{PROJ}_p(\mathbb{C})$, suppose d_p is the L_1 metric on \mathbb{R}_+^p and $d_p((x_1, \dots, x_p), \text{PROJ}_p(\mathbb{C})) < \delta/2$. Then there is $(z_1, \dots, z_p) \in \text{PROJ}_p(\mathbb{C})$ such that $d_p((x_1, \dots, x_p), (z_1, \dots, z_p)) < \delta$. But then if $z \in \mathbb{C}$ with $z_{|p} = (z_1, \dots, z_p)$, we have, since $(z_1, \dots, z_p, 0_\infty) \in \mathbb{C}$ by (4.4),

$$\begin{aligned} d_\infty((x_1, \dots, x_p, 0_\infty), (z_1, \dots, z_p, 0_\infty)) &= \sum_{i=1}^p \frac{|x_i - z_i| \wedge 1}{2^i} \\ &\leq \sum_{i=1}^p |x_i - z_i| \wedge 1 \leq \sum_{i=1}^p |x_i - z_i| \\ &= d_p((x_1, \dots, x_p), (z_1, \dots, z_p)) < \delta, \end{aligned}$$

and therefore $d_p((x_1, \dots, x_p), \text{PROJ}_p(\mathbb{C})) < \delta/2$ implies

$$g(x_1, \dots, x_p) = f(x_1, \dots, x_p, 0_\infty) = 0.$$

So the support of g is bounded away from $\text{PROJ}_p(\mathbb{C})$ as claimed.

Now write

$$\begin{aligned} \mu_n(f) - \mu_0(f) &= [\mu_n(f) - \mu_n(g \circ \text{PROJ}_p)] + [\mu_n(g \circ \text{PROJ}_p) - \mu_0(g \circ \text{PROJ}_p)] \\ &\quad + [\mu_0(g \circ \text{PROJ}_p) - \mu_0(f)] = A + B + C. \end{aligned} \tag{4.8}$$

From (4.5), since $g \circ \text{PROJ}_p \in \mathcal{C}(\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}))$, we have

$$B = \mu_n(g \circ \text{PROJ}_p) - \mu_0(g \circ \text{PROJ}_p) \rightarrow 0$$

as $n \rightarrow \infty$.

How to control A? For $x \in \mathbb{R}_+^\infty \setminus \mathbb{C}$, if $d_\infty((x_1, \dots, x_p, 0_\infty), \mathbb{C}) < \delta/2$, then $f((x_1, \dots, x_p, 0_\infty)) = 0$ and also $d_\infty(x, \mathbb{C}) \leq d_\infty(x, (x_{|p}, 0_\infty)) + d_\infty((x_{|p}, 0_\infty), \mathbb{C}) < 2^{-p} + \delta/2 < \delta$ so $f(x) = 0$. Therefore, on

$$\Lambda = \{x \in \mathbb{R}_+^\infty \setminus \mathbb{C} : d_\infty((x_{|p}, 0), \mathbb{C}) < \delta/2\}$$

both f and $g \circ \text{PROJ}_p$ are zero. Set

$$\Lambda^c = (\mathbb{R}_+^\infty \setminus \mathbb{C}) \setminus \Lambda = \{x \in \mathbb{R}_+^\infty \setminus \mathbb{C} : d_\infty((x_{|p}, 0), \mathbb{C}) \geq \delta/2\}.$$

Then we have

$$|\mu_n(f) - \mu_n(g \circ \text{PROJ}_p)| \leq \int |f - g \circ \text{PROJ}_p| d\mu_n$$

$$\begin{aligned}
 &= \int_{\Lambda^c} |f - g \circ \text{PROJ}_p| d\mu_n \\
 &\leq \mu_n(\Lambda^c) \omega_f(2^{-p}),
 \end{aligned}$$

and similarly for dealing with term C, we would have $|\mu_0(f) - \mu_0(g \circ \text{PROJ}_p)| \leq \mu_0(\Lambda^c) \omega_f(2^{-p})$.

Owing to finite-dimensional convergence (4.5) and (4.8), we have

$$\limsup_{n \rightarrow \infty} |\mu_n(f) - \mu_0(f)| \leq 2\omega_f(2^{-p})\mu_0(\Lambda^c) + 0.$$

Since Λ^c is bounded away from \mathbb{C} , $\mu_0(\Lambda^c) < \infty$ and since the inequality holds for any p sufficiently large such that $2^{-p} < \delta$, we may let $p \rightarrow \infty$ to get $\mu_n(f) \rightarrow \mu_0(f)$. \square

Remark. The proof shows (4.5) only needs to hold for all $p \geq p_0$. For example, if

$$\mathbb{C} = \mathbb{C}_j = \{x \in \mathbb{R}_+^\infty : \sum_{i=1}^\infty \epsilon_{x_i}(0, \infty) \leq j\},$$

then

$$\text{PROJ}_p(\mathbb{C}_j) = \{(x_1, \dots, x_p) \in \mathbb{R}_+^p : \sum_{i=1}^p \epsilon_{x_i}(0, \infty) \leq j\},$$

and for $p < j$, $\text{PROJ}_p(\mathbb{C}_j) = \mathbb{R}_+^p$ and $\mathbb{R}_+^p \setminus \text{PROJ}_p(\mathbb{C}_j) = \emptyset$. However, it suffices for the result to hold for all $p \geq j$.

4.4. Comparing M-convergence on $\mathbb{S} \setminus \mathbb{C}$ with vague convergence when \mathbb{S} is compactified

This continues the discussion of Section 2.3. Conventionally [40] regular variation on $[0, \infty)^p$ has been defined on the punctured compactified space $[0, \infty]^p \setminus \{0_p\}$. This solves the problem of how to make tail regions relatively compact. However, as discussed in [15], when deleting more than $\{0_p\}$, this approach causes problems with the convergence to types lemma and also because certain natural regions are no longer relatively compact. The issue arises when there is mass on the lines through ∞_p , something that is impossible for regular variation on $[0, \infty]^p \setminus \{0_p\}$. The following discussion amplifies what is in [15].

Suppose \mathbb{C} is closed in $[0, \infty]^p$ and set

$$\mathbb{C}_0 = \mathbb{C} \cap [0, \infty)^p, \quad \Omega = [0, \infty]^p \setminus \mathbb{C}, \quad \Omega_0 = [0, \infty)^p \setminus \mathbb{C}_0.$$

Examining the definitions we see that,

- $\Omega_0 \subset \Omega$.
- $\Omega \setminus \Omega_0 = \mathbb{C}^c \cap ([0, \infty]^p \setminus [0, \infty)^p)$

Proposition 4.4. *Suppose for every $n \geq 0$ that $\mu_n \in \mathbb{M}_+(\Omega)$ and μ_n places no mass on the lines through ∞_p :*

$$\mu_n([0, \infty]^p \setminus [0, \infty)^p \cap \mathbb{C}^c) = 0. \tag{4.9}$$

Then

$$\mu_n \xrightarrow{v} \mu_0 \quad \text{in } \mathbb{M}_+(\Omega), \tag{4.10}$$

if and only if the restrictions to the space without the lines through ∞_p converge:

$$\mu_{n0} := \mu_n(\cdot \cap \Omega_0) \rightarrow \mu_0(\cdot \cap \Omega_0) =: \mu_{00} \quad \text{in } \mathbb{M}(\Omega_0). \tag{4.11}$$

Proof. Given (4.11), let $f \in \mathcal{C}_K^+(\Omega)$. Then the restriction to Ω_0 satisfies $f|_{\Omega_0} \in \mathcal{C}(\Omega_0)$ so

$$\mu_n(f) = \mu_{n0}(f|_{\Omega_0}) \rightarrow \mu_{00}(f|_{\Omega_0}) = \mu_0(f),$$

so $\mu_n \xrightarrow{v} \mu_0$ in $M_+(\Omega)$.

Conversely, assume (4.10). Suppose $B \in \mathcal{S}(\Omega_0)$ and $\mu_{00}(\partial_{\Omega_0} B) = 0$, where $\partial_{\Omega_0} B$ is the set of boundary points of B in Ω_0 . This implies $\mu_0(\partial_{\Omega} B) = 0$ since

$$\partial_{\Omega}(B) \subset \partial_{\Omega_0} B \cup \left((\Omega \setminus \Omega_0) \cap \mathbb{C} \right).$$

Therefore $\mu_n(B) \rightarrow \mu_0(B)$ and because of (4.9), $\mu_{n0}(B) \rightarrow \mu_{00}(B)$ which proves (4.11). □

4.5. Regular variation on \mathbb{R}_+^p and \mathbb{R}_+^∞

For this section, either \mathbb{S} is \mathbb{R}_+^p or \mathbb{R}_+^∞ and \mathbb{C} is a closed cone; then $\mathbb{S} \setminus \mathbb{C}$ is still a cone. Applying Definition 3.2, a random element X of $\mathbb{S} \setminus \mathbb{C}$ has a regularly varying distribution if for some regularly varying function $b(t) \rightarrow \infty$, as $t \rightarrow \infty$,

$$tP[X/b(t) \in \cdot] \rightarrow \nu(\cdot) \quad \text{in } \mathbb{M}(\mathbb{S} \setminus \mathbb{C}),$$

for some limit measure $\nu \in \mathbb{M}(\mathbb{S} \setminus \mathbb{C})$. In \mathbb{R}_+^p , if $\mathbb{C} = \{0_p\}$ or if (4.9) holds, this definition is the same as the one using vague convergence on the compactified space.

4.5.1. The iid case: Remove $\{0_\infty\}$

Suppose $X = (X_1, X_2, \dots)$ is iid with non-negative components, each of which has a regularly varying distribution on $(0, \infty)$ satisfying

$$P[X_1 > tx] / P[X_1 > t] \rightarrow x^{-\alpha}, \quad \text{as } t \rightarrow \infty, x > 0, \alpha > 0.$$

Equivalently, as $t \rightarrow \infty$,

$$tF(b(t)\cdot) := tP[X_1/b(t) \in \cdot] \rightarrow \nu_\alpha(\cdot) \quad \text{in } \mathbb{M}((0, \infty)), \tag{4.12}$$

where $\nu_\alpha(x, \infty) = x^{-\alpha}$, $\alpha > 0$. Then in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \{0_\infty\})$, we have

$$\begin{aligned} \mu_t((dx_1, dx_2, \dots)) &:= tP[X/b(t) \in (dx_1, dx_2, \dots)] \\ &\rightarrow \sum_{l=1}^{\infty} \prod_{i \neq l} \epsilon_0(dx_i) \nu_\alpha(dx_l) =: \mu^{(0)}(dx_1, dx_2, \dots), \end{aligned} \quad (4.13)$$

and the limit measure concentrates on

$$\mathbb{C}_{=1} = \{x \in \mathbb{R}_+^\infty : \sum_{i=1}^{\infty} \epsilon_{x_i}((0, \infty)) = 1\},$$

the set of sequences with exactly one component positive. Note $\{0_\infty\} \cup \mathbb{C}_{=1} =: \mathbb{C}_{\leq 1}$, the set of sequences with at most one component positive, is closed.

To verify (4.13), note from Theorem 4.1, it suffices to verify finite-dimensional convergence since $\{0_\infty\}$ satisfies (4.4), so it suffices to prove as $t \rightarrow \infty$, for $p \geq 1$,

$$\begin{aligned} \mu_t \circ \text{PROJ}_p^{-1}((dx_1, \dots, dx_p)) &:= tP[(X_1, \dots, X_p)/b(t) \in (dx_1, \dots, dx_p)] \\ &\rightarrow \mu^{(0)} \circ \text{PROJ}_p^{-1}((dx_1, \dots, dx_p)) = \sum_{l=1}^p \prod_{i \neq l} \epsilon_0(dx_i) \nu_\alpha(dx_l), \end{aligned} \quad (4.14)$$

in $\mathbb{M}(\mathbb{R}_+^p \setminus \{0_p\})$. Since neither $\mu_t \circ \text{PROJ}_p^{-1}$ nor $\mu^{(0)} \circ \text{PROJ}_p^{-1}$ place mass on the lines through ∞_p , \mathbb{M} -convergence and vague convergence are the same and then (4.14) follows from the binding lemma in [40, p. 228, 210].

Applying the operator CUMSUM and Corollary 4.1 to (4.13) gives

$$\begin{aligned} tP[\text{CUMSUM}(X)/b(t) \in (dx_1, dx_2, \dots)] &\rightarrow \mu^{(0)} \circ \text{CUMSUM}^{-1}((dx_1, dx_2, \dots)) \\ &= \sum_{l=1}^{\infty} \prod_{i=1}^{l-1} \epsilon_0(dx_i) \nu_\alpha(dx_l) \prod_{i=l+1}^{\infty} \epsilon_{x_i}(dx_i) \quad \text{in } \mathbb{M}(\mathbb{R}_+^\infty \setminus \{0_\infty\}), \end{aligned} \quad (4.15)$$

where the limit concentrates on non-decreasing sequences with one jump and the size of the jump is governed by ν_α . Then applying the operator PROJ_p we get, by Corollary 4.2,

$$\begin{aligned} tP[(X_1, X_1 + X_2, \dots, \sum_{i=1}^p X_k)/b(t) \in (dx_1, dx_2, \dots, dx_p)] \\ &\rightarrow \mu^{(0)} \circ \text{CUMSUM}^{-1} \circ \text{PROJ}_p^{-1}((dx_1, dx_2, \dots, dx_p)) \\ &= \sum_{l=1}^p \prod_{i=1}^{l-1} \epsilon_0(dx_i) \nu_\alpha(dx_l) \prod_{i=l+1}^p \epsilon_{x_i}(dx_i) \quad \text{in } \mathbb{M}(\mathbb{R}_+^p \setminus \{0_p\}), \end{aligned} \quad (4.16)$$

giving an elaboration of the one-big-jump heuristic saying that summing independent risks which have the same heavy tail results in a tail risk which is the number of summands times the individual tail risk; for example, see [40,

p. 230]. In particular, applying the projection from $\mathbb{R}_+^p \setminus \{0_p\} \mapsto (0, \infty)$ defined by $T : (x_1, \dots, x_p) \mapsto x_p$ gives by Corollary 2.1 that

$$\begin{aligned} tP\left[\sum_{i=1}^p X_i > b(t)x\right] &\rightarrow \mu^{(0)} \circ \text{CUMSUM}^{-1} \circ \text{PROJ}_p^{-1} \circ T^{-1}(x, \infty) \\ &= p\nu_\alpha(x, \infty) = px^{-\alpha}. \end{aligned} \quad (4.17)$$

The projection T is uniformly continuous but also Theorem 2.3 applies to T since for $y > 0$, $T^{-1}(y, \infty) = \{(x_1, \dots, x_p) : x_p > y\}$ is at positive distance from $\{0_p\}$.

The above discussion could have been carried out with minor modifications without the iid assumption by assuming (4.12) and

$$P[X_j > x]/P[X_1 > x] \rightarrow c_j > 0, \quad j \geq 2.$$

4.5.2. The iid case: Remove more; Hidden regular variation

We now investigate how to get past the one-big-jump heuristic by using hidden regular variation. For $j \geq 1$, set

$$\begin{aligned} \mathbb{C}_{=j} &= \{x \in \mathbb{R}_+^\infty : \sum_{i=1}^\infty \epsilon_{x_i}((0, \infty)) = j\}, \\ \mathbb{C}_{\leq j} &= \{x \in \mathbb{R}_+^\infty : \sum_{i=1}^\infty \epsilon_{x_i}((0, \infty)) \leq j\} = \mathbb{C}_{\leq(j-1)} \cup \mathbb{C}_{=j}, \end{aligned} \quad (4.18)$$

so that $\mathbb{C}_{\leq j}$ is closed. We imagine an infinite sequence of reductions of the state space with scaling adjusted at each step. This is suggested by the previous discussion. On $\mathbb{M}(\mathbb{R}_+^\infty \setminus \{0_\infty\})$, the limit measure $\mu^{(0)}$ concentrated on $\mathbb{C}_{=1}$, a small part of the potential state space. Remove $\{0_\infty\} \cup \mathbb{C}_{=1} = \mathbb{C}_{\leq 1}$ and on $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C}_{\leq 1})$ seek a new convergence using adjusted scaling $b(\sqrt{t})$. We get in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C}_{\leq 1})$ as $t \rightarrow \infty$,

$$\begin{aligned} \mu_t^{(1)}(dx_1, dx_2, \dots) &= tP[X/b(\sqrt{t}) \in (dx_1, dx_2, \dots)] \rightarrow \mu^{(1)}((dx_1, dx_2, \dots)) \\ &:= \sum_{l < k} \left(\prod_{j \notin \{l, k\}} \epsilon_0(dx_j) \right) \nu_\alpha(dx_l) \nu_\alpha(dx_k), \end{aligned} \quad (4.19)$$

which concentrates on $\mathbb{C}_{=2}$. In general, we find that in $\mathbb{M}(\mathbb{R}_+^\infty \setminus \mathbb{C}_{\leq j})$ as $t \rightarrow \infty$,

$$\begin{aligned} \mu_t^{(j)}(dx_1, dx_2, \dots) &= tP[X/b(t^{1/(j+1)}) \in (dx_1, dx_2, \dots)] \rightarrow \mu^{(j)}((dx_1, dx_2, \dots)) \\ &:= \sum_{i_1 < i_2 < \dots < i_{j+1}} \left(\prod_{j \notin \{i_1, \dots, i_{j+1}\}} \epsilon_0(dx_j) \right) \nu_\alpha(dx_{i_1}) \nu_\alpha(dx_{i_2}) \dots \nu_\alpha(dx_{i_{j+1}}), \end{aligned} \quad (4.20)$$

which concentrates on $\mathbb{C}_{=(j+1)}$. This is an elaboration of results in [30, 33, 34]. The result in \mathbb{R}_+^∞ can be proven by reducing to \mathbb{R}_+^p by means of Theorem 4.1

TABLE 1
An infinite number of coexisting regular-variation properties

j	remove	scaling	$\mu^{(j)}$	support
1	$\{0\}$	$b(t)$	$\sum_{i=1}^{\infty} \nu_{\alpha}(dx_i) [\prod_{i \neq l} \epsilon_0(dx_i)]$	axes
2	axes	$b(\sqrt{t})$	$\sum_{l,m} \nu_{\alpha}(dx_l) \nu_{\alpha}(dx_m) \left[\prod_{i \notin \{l,m\}} \epsilon_0(dx_i) \right]$	2-dim faces
\vdots	\vdots	\vdots	\vdots	\vdots
m	$\mathbb{C}_{\leq(m-1)}$	$b(t^{\frac{1}{m}})$	$\sum_{(l_1, \dots, l_m)_{p=1}^m} \prod_{p=1}^m \nu_{\alpha}(dx_{l_p}) \left[\prod_{i \notin \{l_1, \dots, l_m\}} \epsilon_0(dx_i) \right]$	$\mathbb{C}_{=m}$

noting that $\mathbb{C}_{\leq j}$ satisfies (4.4) and then observing that neither $\mu_t^{(j)}$ nor $\mu^{(j)}$ puts mass on lines through ∞_p . It is enough to show convergences of the following form: Assume $p \geq j$ and $i_1 < i_2 < \dots < i_{j+1}$ and $y_l > 0$, $l = 1, \dots, j+1$ and

$$\begin{aligned} tP[X_{i_l} > b(t^{1/(j+1)})y_l, l = 1, \dots, j+1] &= \prod_{l=1}^{j+1} t^{1/(j+1)} P[X_{i_l} > b(t^{1/(j+1)})y_l] \\ &\rightarrow \prod_{l=1}^{j+1} \nu_{\alpha}(y_l, \infty) = \prod_{l=1}^{j+1} y_l^{-\alpha}. \end{aligned}$$

A formal statement of the result and a proof relying on a convergence-determining class is given in the next Section 4.5.3. Table 1 gives a summary of the results in tabular form.

Proposition 4.2 and Corollary 2.1 allow application of CUMSUM to get

$$\begin{aligned} \mu_t^{(j)} \circ \text{CUMSUM}^{-1}(dx_1, dx_2, \dots) \\ &= tP[\text{CUMSUM}(X)/b(t^{1/(j+1)}) \in (dx_1, dx_2, \dots)] \\ &\rightarrow \mu^{(j)} \circ \text{CUMSUM}^{-1}((dx_1, dx_2, \dots)) \end{aligned} \quad (4.21)$$

in $\mathbb{M}(\text{CUMSUM}(\mathbb{R}_+^{\infty}) \setminus \text{CUMSUM}(\mathbb{C}_{\leq j}))$. Note $\text{CUMSUM}(\mathbb{R}_+^{\infty}) =: \mathbb{R}_+^{\infty \uparrow}$ is the set of non-decreasing sequences and $\text{CUMSUM}(\mathbb{C}_{\leq j}) =: \mathbb{S}_{\leq j}$ is the set of non-decreasing sequences with at most j positive jumps. Now apply the map PROJ_p to (4.21) to get a p -dimensional result for $(X_1, X_1 + X_2, \dots, X_1 + \dots + X_p)$ and the analogue of (4.16) is

$$\begin{aligned} tP[(X_1, X_1 + X_2, \dots, \sum_{i=1}^p X_k)/b(t^{1/(j+1)}) \in (dx_1, dx_2, \dots, dx_p)] \\ \rightarrow \mu^{(j)} \circ \text{CUMSUM}^{-1} \circ \text{PROJ}_p^{-1}((dx_1, dx_2, \dots, dx_p)) \end{aligned} \quad (4.22)$$

in $\mathbb{M}(\text{CUMSUM}(\mathbb{R}_+^p) \setminus \text{PROJ}_p(\mathbb{S}_{\leq j})) = \mathbb{M}(\mathbb{R}_+^{p \uparrow}) \setminus \text{PROJ}_p(\mathbb{S}_{\leq j})$.

When $j > 1$, unlike the step leading to (4.17), we cannot apply the map $T : (x_1, \dots, x_p) \mapsto x_p$ to (4.22) to get a marginal result for $X_1 + \dots + X_p$.

Although T is uniformly continuous, Corollary 2.1 is not applicable since

$$T(\mathbb{R}_+^{\uparrow}) \setminus T(\text{PROJ}_p(\mathbb{S}_{\leq j})) = [0, \infty) \setminus [0, \infty) = \emptyset.$$

4.5.3. *The iid case: HRV; Formal statement and proof*

Recall $X = (X_l, l \geq 1)$ has iid components each of which has a distribution with a regularly varying tail of index $\alpha > 0$. Define $\mathbb{C}_{\leq j}$ as in (4.18) and set $\mathbb{O}_j = \mathbb{R}_+^\infty \setminus \mathbb{C}_{\leq j}$. The definition of μ_t and $\mu^{(j)}$ are given in (4.20).

Theorem 4.2. *For every $j \geq 1$ there is a nonzero measure $\mu^{(j)} \in \mathbb{M}_{\mathbb{O}_j}$ with support in $\mathbb{C}_{=(j+1)}$ such that $tP[X/b(t^{1/j+1}) \in \cdot] \rightarrow \mu^{(j)}(\cdot)$ in $\mathbb{M}_{\mathbb{O}_j}$ as $t \rightarrow \infty$. The measure $\mu^{(j)}$ is given in (4.20), or more formally,*

$$\mu^{(j)}(A) = \sum_{(i_1, \dots, i_{j+1})} \int I\left\{\sum_{k=1}^{j+1} z_k e_{i_k} \in A\right\} \nu_\alpha(dz_1) \dots \nu_\alpha(dz_{j+1}),$$

where the components of e_{i_k} are all zero except component i_k whose value is 1 and the indices (i_1, \dots, i_{j+1}) run through the ordered subsets of size $j + 1$ of $\{1, 2, \dots\}$.

The proof of Theorem 4.2 uses a particular convergence-determining class $\mathbb{A}_{\geq j}$ of subsets of \mathbb{O}_j . Let $\mathbb{A}_{\geq j}$ denote the set of sets $A_{m,i,a}$ for $m \geq j$, where

$$A_{m,i,a} = \{x \in \mathbb{R}_+^\infty : x_{i_k} > a_k \text{ for } k = 1, \dots, m\}, \quad i_1 < \dots < i_m, a_1, \dots, a_m > 0.$$

Lemma 4.1. *If $\mu_t, \mu \in \mathbb{M}_{\mathbb{O}_j}$ and $\lim_{t \rightarrow \infty} \mu_t(A) = \mu(A)$ for all $A \in \mathbb{A}_{\geq j}$ bounded away from \mathbb{C}_j with $\mu(\partial A) = 0$, then $\mu_t \rightarrow \mu$ in $\mathbb{M}_{\mathbb{O}_j}$ as $t \rightarrow \infty$.*

Proof. Consider the set of finite differences of sets in $\mathbb{A}_{\geq j}$ and note that this set is a π -system. Take $x \in \mathbb{O}_j$ and $\epsilon > 0$. Since $x \in \mathbb{O}_j$ there are $i_1 < \dots < i_j$ such that $x_{i_k} > 0$ for each k . If $2^{-i_j} < \epsilon/2$ choose $m = i_j$. Otherwise, choose $m > i_j$ such that $2^{-m} < \epsilon/2$. Take $\delta < \min\{\epsilon/2, \min\{x_k : x_k > 0 \text{ and } k \leq m\}\}$ and set

$$B = \{y \in \mathbb{R}_+^\infty : y_k \geq 0 \text{ if } x_k = 0 \text{ and } y_k > x_k - \delta \text{ otherwise for } k \leq m\},$$

$$B' = \{y \in \mathbb{R}_+^\infty : y_k > \delta \text{ if } x_k = 0 \text{ and } y_k > x_k + \delta \text{ otherwise for } k \leq m\}.$$

Then $B, B' \in \mathbb{A}_{\geq j}$, B' is a proper subset of B , and $z \in B \setminus B'$ implies that $d(z, x) < \delta \sum_{k=1}^m 2^{-k} + \epsilon/2 < \epsilon$, i.e. that $z \in B_{x,\epsilon}$. Moreover,

$$(B \setminus B')^\circ = \{y \in \mathbb{R}_+^\infty : y_k \in J(x_k) \text{ for } k \leq m\},$$

where $J(x_k) = [0, \delta)$ if $x_k = 0$ and $J(x_k) = (x_k - \delta, x_k + \delta)$ if $x_k \neq 0$. Finally, $\partial(B \setminus B')$ is the set of $y \in \mathbb{R}_+^\infty$ such that $y_k \in [\max\{0, x_k - \delta\}, x_k + \delta]$ for all $k \leq m$ and $y_k = \delta$ or $y_k = x_k \pm \delta$ for some $k \leq m$. In particular, there is an uncountable set of δ -values, for which the boundaries $\partial(B \setminus B')$ are disjoint, satisfying the requirements. Therefore δ can without loss of generality be chosen

so that $\mu(\partial(B \setminus B')) = 0$. The separability of \mathbb{R}_+^∞ implies (cf. the proof of Theorem 2.3 in [6]) that each open set is a countable union of μ -continuity sets of the form $(B \setminus B')^\circ$. The same argument as in the proof of Theorem 2.2 in [6] therefore shows that $\liminf_{t \rightarrow \infty} \mu_t(G) \geq \mu(G)$ for all open $G \subset \mathbb{O}_j$ bounded away from \mathbb{C}_j . Any closed set $F \subset \mathbb{O}_j$ bounded away from \mathbb{C}_j is a subset of some $A \in \mathbb{A}_{\geq j}$. By the same argument as above, we may without loss of generality take A such that $\mu(\partial A) = 0$. The set $A \setminus F$ is open and therefore

$$\mu(A) - \limsup_{t \rightarrow \infty} \mu_t(F) = \liminf_{t \rightarrow \infty} \mu_t(A \setminus F) \geq \mu(A \setminus F) = \mu(A) - \mu(F),$$

i.e. $\limsup_{t \rightarrow \infty} \mu_t(F) \leq \mu(F)$. The conclusion follows from Theorem 2.1(iii). \square

Proof of Theorem 4.2. For any $m > j$ and $a_1, \dots, a_m > 0$, the following two limits hold,

$$\lim_{t \rightarrow \infty} c(t)^j P[X \in tA_{j,i,a}] = \prod_{k=1}^j a_k^{-\alpha} = \mu_j(A_{j,i,a}), \quad \lim_{t \rightarrow \infty} c(t)^j P[X \in tA_{j+1,i,a}] = 0.$$

Therefore, the support of μ is a subset of $\mathbb{C}_{j+1} \setminus \mathbb{C}_j$. Notice that for $j \geq 1$

$$(\mathbb{C}_{j+1} \setminus \mathbb{C}_j) \cap \mathbb{R}_+^p = \cup_{i_1 < \dots < i_j} \{(\lambda_1 e_{i_1}, \dots, \lambda_j e_{i_j}); \lambda_1, \dots, \lambda_j > 0\},$$

where indices i_1, \dots, i_j run through ordered subsets of size j of $\{1, 2, \dots\}$. \square

4.5.4. Poisson points as random elements of \mathbb{R}_+^∞

Considering Poisson points provides a variant to the iid case and leads naturally to considering regular variation of the distribution of a Lévy process with regularly varying measure.

Suppose $\nu \in \mathbb{M}(0, \infty)$ and $x \mapsto \nu(x, \infty)$ is regularly varying at infinity with index $-\alpha < 0$. If $Q(x) = \nu([x, \infty))$, define $Q^-(y) = \inf\{t > 0 : \nu([t, \infty)) < y\}$. Then the function b given by $b(t) = Q^-(1/t)$ satisfies $\lim_{t \rightarrow \infty} t\nu(b(t)x, \infty) = x^{-\alpha}$. It follows that b is regularly varying at infinity with index $1/\alpha$.

Let $\{E_n, n \geq 1\}$ be iid standard exponentially distributed random variables so that if $\{\Gamma_n, n \geq 1\} := \text{CUMSUM}\{E_n, n \geq 1\}$, we get points of a homogeneous Poisson process of rate 1. Transforming [40, p. 121], we find $\{Q^-(\Gamma_n), n \geq 1\}$ are points of a Poisson process with mean measure ν , written in decreasing order.

Define the following subspaces of \mathbb{R}_+^∞ :

$$\begin{aligned} \mathbb{R}_+^{\infty \downarrow} &= \{x \in \mathbb{R}_+^\infty : x_1 \geq x_2 \geq \dots\}, \\ \mathbb{H}_{=j} &= \{x \in \mathbb{R}_+^{\infty \downarrow} : x_j > 0, x_{j+1} = 0\}, \\ \mathbb{H}_{\leq j} &= \{x \in \mathbb{R}_+^{\infty \downarrow} : x_{j+1} = 0\}, \quad \mathbb{O}_j = \mathbb{R}_+^{\infty \downarrow} \setminus \mathbb{H}_{\leq j}, \end{aligned} \tag{4.23}$$

with the usual meaning of multiplication by a scalar. So $\mathbb{H}_{\leq 0} = \{0_\infty\}$ and $\mathbb{R}_+^{\infty \downarrow}$ are sequences with decreasing, non-negative components and $\mathbb{H}_{\leq j}$ are decreasing

sequences such that components are 0 from the $(j + 1)$ st component onwards. Furthermore, for each $j \geq 1$, $\mathbb{H}_{\leq j}$ is closed. To verify the closed property, suppose $\{x(n), n \geq 1\}$ is a sequence in $\mathbb{H}_{\leq j}$ and $x(n) \rightarrow x(\infty)$ in the \mathbb{R}_+^∞ metric. This means componentwise convergence, so for the m th component convergence, where $m > j$, as $n \rightarrow \infty$, $0 = x_m(n) \rightarrow x_m(\infty)$ and $x(\infty)$ is 0 beyond the j th component. The monotonicity of the components for each $x(n)$ is preserved by taking limits. Hence $\mathbb{H}_{\leq j}$ is closed.

Analogous to (4.13), we claim

$$tP[(Q^{\leftarrow}(\Gamma_l)/b(t), l \geq 1) \in \cdot] \rightarrow \mu^{(1)}(\cdot),$$

in $\mathbb{M}(\mathbb{O}_0)$ as $t \rightarrow \infty$, where

$$\mu^{(1)}(dx_1 \times dx_2 \times \dots) = \nu_\alpha(dx_1)1_{[x_1 > 0]} \prod_{l=2}^\infty \epsilon_0(dx_l).$$

To verify this, it suffices to prove finite dimensional convergence and for the biggest component and $x > 0$,

$$\begin{aligned} tP[Q^{\leftarrow}(\Gamma_1)/b(t) > x] &= tP[\Gamma_1 \leq Q(b(t)x)] = t(1 - e^{-Q(b(t)x)}) \\ &\sim tQ(b(t)x) \rightarrow x^{-\alpha} = \nu_\alpha(x, \infty). \end{aligned}$$

For the first two components, let $\text{PRM}(\nu)$ be a Poisson counting function with mean measure ν and for $x > 0, y > 0$,

$$tP[Q^{\leftarrow}(\Gamma_1)/b(t) > x, Q^{\leftarrow}(\Gamma_2)/b(t) > y] \leq tP[\text{PRM}(\nu)(b(t)(x \wedge y, \infty) \geq 2]$$

and writing $p(t) = \nu(b(t)(x \wedge y, \infty))$, we have

$$\begin{aligned} tP[\text{PRM}(\nu)(b(t)(x \wedge y, \infty) \geq 2] &= t(1 - e^{-p(t)} - p(t)e^{-p(t)}) \\ &\leq t(p(t) - p(t)e^{-p(t)}) \leq tp^2(t) \rightarrow 0. \end{aligned}$$

The conclusion now follows from Lemma 4.1 by observing that we have shown convergence for the sets in a convergence-determining class.

Similarly, we claim

$$tP[(Q^{\leftarrow}(\Gamma_l)/b(t^{1/2}), l \geq 1) \in \cdot] \rightarrow \mu^{(2)}(\cdot)$$

in $\mathbb{M}(\mathbb{O}_1)$ as $t \rightarrow \infty$, where

$$\mu^{(2)}(dx_1 \times dx_2 \times \dots) = \nu_\alpha(dx_1)\nu_\alpha(dx_2)1_{[x_1 \geq x_2 > 0]} \prod_{l=3}^\infty \epsilon_0(dx_l).$$

Simple computations show that the distribution of $(\Gamma_1, \Gamma_2) = (E_1, E_1 + E_2)$ satisfies

$$P[\Gamma_1 \leq z, \Gamma_2 \leq w] = \begin{cases} 1 - e^{-z} - ze^{-w}, & z < w, \\ 1 - e^{-w} - we^{-w}, & z \geq w. \end{cases}$$

Notice that, for $x > y > 0$,

$$\begin{aligned} & P[Q^\leftarrow(\Gamma_1)/b(t^{1/2}) > x, Q^\leftarrow(\Gamma_2)/b(t^{1/2}) > y] \\ &= P[\Gamma_1 \leq Q(b(t^{1/2})x), \Gamma_2 \leq Q(b(t^{1/2})y)] \\ &= 1 - e^{-Q(b(t^{1/2})x)} - Q(b(t^{1/2})x)e^{-Q(b(t^{1/2})y)} \\ &\sim Q(b(t^{1/2})x) - Q(b(t^{1/2})x)^2/2 + O(Q(b(t^{1/2})x)^3) \\ &\quad - Q(b(t^{1/2})x)\left(1 - Q(b(t^{1/2})y) + O(Q(b(t^{1/2})y)^2)\right). \end{aligned}$$

In particular, it is easy to verify that for $x > y > 0$

$$\begin{aligned} & \lim_{t \rightarrow \infty} tP[Q^\leftarrow(\Gamma_1)/b(t^{1/2}) > x, Q^\leftarrow(\Gamma_2)/b(t^{1/2}) > y] \\ &= x^{-\alpha}y^{-\alpha} - x^{-2\alpha}/2 \\ &= \mu^{(2)}(z \in \mathbb{R}^{\infty \downarrow} : z_1 > x, z_2 > y). \end{aligned}$$

Similar computations show that, for $y > x > 0$,

$$\begin{aligned} & \lim_{t \rightarrow \infty} tP[Q^\leftarrow(\Gamma_1)/b(t^{1/2}) > x, Q^\leftarrow(\Gamma_2)/b(t^{1/2}) > y] \\ &= y^{-2\alpha}/2 \\ &= \mu^{(2)}(z \in \mathbb{R}^{\infty \downarrow} : z_1 > x, z_2 > y). \end{aligned}$$

Moreover, for $x > 0, y > 0, z > 0$,

$$\begin{aligned} & tP[Q^\leftarrow(\Gamma_1)/b(t^{1/2}) > x, Q^\leftarrow(\Gamma_2)/b(t^{1/2}) > y, Q^\leftarrow(\Gamma_3)/b(t^{1/2}) > z] \\ &\leq tP[\text{PRM}(\nu)(b(t^{1/2})(x \wedge y \wedge z, \infty)) \geq 3], \end{aligned}$$

and writing $p(t) = \nu(b(t^{1/2})(x \wedge y \wedge z, \infty))$, we have

$$\begin{aligned} & tP[\text{PRM}(\nu)(b(t^{1/2})(x \wedge y \wedge z, \infty)) \geq 3] \\ &= t(1 - e^{-p(t)} - p(t)e^{-p(t)} - p(t)^2e^{-p(t)}/2) \\ &\sim t(p(t)^3/3! + o(p(t)^3)) \end{aligned}$$

as $t \rightarrow \infty$. Hence, $\lim_{t \rightarrow \infty} tP[\text{PRM}(\nu)(b(t^{1/2})(x \wedge y \wedge z, \infty)) \geq 3] = 0$.

As in the iid case described by Theorem 4.2 and (4.20), we have an infinite number of regular-variation properties co-existing.

Theorem 4.3. *For the Poisson points $\{Q^\leftarrow(\Gamma_l), l \geq 1\}$, for every $j \geq 1$, we have*

$$tP[(Q^\leftarrow(\Gamma_l)/b(t^{1/j}), l \geq 1) \in \cdot] \rightarrow \mu^{(j)}(\cdot), \quad (4.24)$$

in $\mathbb{M}_{\mathbb{O}_{j-1}}$ as $t \rightarrow \infty$, where $\mu^{(j)}$ is a measure concentrating on $\mathbb{H}_{=j}$ given by

$$\mu^{(j)}(dx_1, dx_2, \dots) = \prod_{i=1}^j \nu_\alpha(dx_i) 1_{[x_1 \geq x_2 \geq \dots \geq x_j > 0]} \prod_{i=j+1}^{\infty} \epsilon_0(dx_i). \quad (4.25)$$

Proof. The explicit computations above, and similarly for $j \geq 3$, together with an application of Lemma 4.1 yields the conclusion. \square

5. Finding the hidden jumps of a Lévy process

In this section we consider a real-valued Lévy process $X = \{X_t, t \in [0, 1]\}$ as a random element of $\mathbb{D} := \mathbb{D}([0, 1], \mathbb{R})$, the space of real-valued càdlàg functions on $[0, 1]$. We metrize \mathbb{D} with the usual Skorohod metric

$$d_{\text{sk}}(x, y) = \inf_{\lambda \in \Lambda} \|\lambda - e\| \vee \|x \circ \lambda - y\|,$$

where $x, y \in \mathbb{D}$, λ is a non-decreasing homeomorphism of $[0, 1]$ onto itself, Λ is the set of all such homeomorphisms, $e(t) = t$ is the identity, and $\|x\| = \sup_{t \in [0, 1]} |x(t)|$ is the sup-norm. The space \mathbb{D} is not complete under the metric d_{sk} , but there is an equivalent metric under which \mathbb{D} is complete [6, page 125]. Therefore, the space \mathbb{D} fits into the framework presented in Section 2 and we may use the Skorohod metric to check continuity of mappings.

For simplicity we suppose X has only positive jumps and its Lévy measure ν concentrates on $(0, \infty)$. Suppose $x \mapsto \nu(x, \infty)$ is regularly varying at infinity with index $-\alpha < 0$. Let $Q(x) = \nu([x, \infty))$ and define $Q^{\leftarrow}(y) = \inf\{t > 0 : \nu([t, \infty)) < y\}$. Then the function b given by $b(t) = Q^{\leftarrow}(1/t)$ satisfies $\lim_{t \rightarrow \infty} t\nu(b(t)x, \infty) = x^{-\alpha}$ and b is regularly varying at infinity with index $1/\alpha$. It is shown in [23, 25] that with scaling function $b(t)$, the distribution of X is regularly varying on $\mathbb{D} \setminus \{0\}$ with a limit measure concentrating on functions which are constant except for one jump. Where did the other Lévy process jumps go? Using weaker scaling and biting more out of \mathbb{D} than just the zero-function 0 allows recovery of the other jumps.

The standard Itô representation [1, 5, 27] of X is

$$X_t = ta + B_t + \int_{|x| \leq 1} x[N([0, t] \times dx) - t\nu(dx)] + \int_{|x| > 1} xN([0, t] \times dx),$$

where B is standard Brownian motion independent of the Poisson random measure N on $[0, 1] \times (0, \infty)$ with mean measure $\text{Leb} \times \nu$. Referring to the discussion preceding (4.23), $\{Q^{\leftarrow}(\Gamma_n), n \geq 1\}$ are points written in decreasing order of a Poisson random measure on $(0, \infty)$ with mean measure ν and by augmentation [40, p. 122], we can represent

$$N = \sum_{l=1}^{\infty} \epsilon_{(U_l, Q^{\leftarrow}(\Gamma_l))},$$

where $(U_l, l \geq 1)$ are iid standard uniform random variables independent of $\{\Gamma_n\}$.

The Lévy-Itô decomposition allows X to be decomposed into the sum of two independent Lévy processes,

$$X = \tilde{X} + J, \tag{5.1}$$

where J is a compound Poisson process of large jumps bounded from below by 1, and $\tilde{X} = X - J$ is a Lévy process of small jumps that are bounded from above by 1. The compound Poisson process can be represented as the random sum $J = \sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l)1_{[U_l, 1]}$, where $N_1 = N([0, 1] \times [1, \infty))$.

Recall the notation in (4.23) for $\mathbb{R}_+^{\infty\downarrow}$, $\mathbb{H}_{=j}$ and $\mathbb{H}_{\leq j}$ and the result in Theorem 4.3. We seek to convert a statement like (4.24) into a statement about X . The first step is to augment (4.24) with a sequence of iid standard uniform random variables. The uniform random variables will eventually serve as jump times for the Lévy process. The following result is an immediate consequence of Theorem 4.3.

Proposition 5.1. *Under the given assumptions on ν and Q , for $j \geq 1$,*

$$tP\left[\left((Q^{\leftarrow}(\Gamma_l)/b(t^{1/j}), l \geq 1), (U_l, l \geq 1)\right) \in \cdot\right] \rightarrow (\mu^{(j)} \times L)(\cdot) \quad (5.2)$$

in $\mathbb{M}((\mathbb{R}_+^{\infty\downarrow} \setminus \mathbb{H}_{\leq j-1}) \times [0, 1]^\infty)$ as $t \rightarrow \infty$, where L is Lebesgue measure on $[0, 1]^\infty$ and $\mu^{(j)}$ concentrates on $\mathbb{H}_{=j}$ and is given by (4.25).

Think of (5.2) as regular variation on the product space $\mathbb{R}_+^{\infty\downarrow} \times [0, 1]^\infty$ when multiplication by a scalar is defined as $(\lambda, (x, y)) \mapsto (\lambda x, y)$.

Recall ν_α is the Pareto measure on $(0, \infty)$ satisfying $\nu_\alpha(x, \infty) = x^{-\alpha}$, for $x > 0$, and we denote by ν_α^j product measure generated by ν_α with j factors. For $m \geq 0$, let $\mathbb{D}_{\leq m}$ be the subspace of the Skorohod space \mathbb{D} consisting of nondecreasing step functions with at most m jumps and define A_m as

$$A_m = \{(x, u) \in \mathbb{R}_+^{\infty\downarrow} \times [0, 1]^\infty : u_i \in (0, 1) \text{ for } 1 \leq i \leq m; u_i \neq u_j \text{ for } i \neq j, 1 \leq i, j \leq m\}. \quad (5.3)$$

Let T_m be the map

$$T_m : A_m \mapsto \mathbb{D} \text{ defined by } T_m(x, u) = \sum_{i=1}^m x_i 1_{[u_i, 1]}, \quad (5.4)$$

and we think of T_m as mapping a jump-size sequence and a sequence of distinct jump times into a step function in $\mathbb{D}_{\leq m} \subset \mathbb{D}$. Our approach applies T_m to the convergence in (5.2) to get a sequence of regular-variation properties of the distribution of X . Whereas in Section 4.5.2, we could rely on uniform continuity of CUMSUM, T_m is not uniformly continuous and hence the mapping Theorem 2.3 must be used and its hypotheses verified. We will prove the following.

Theorem 5.1. *Under the regular variation assumptions on ν and Q , for $j \geq 1$,*

$$\begin{aligned} tP[X/b(t^{1/j}) \in \cdot] &\rightarrow (\mu^{(j)} \times L) \circ T_j^{-1}(\cdot) \\ &= E\left[\nu_\alpha^j\left\{y \in (0, \infty)^j : \sum_{i=1}^j y_i 1_{[U_i, 1]} \in \cdot\right\}\right] \end{aligned} \quad (5.5)$$

in $\mathbb{M}(\mathbb{D} \setminus \mathbb{D}_{\leq j-1})$ as $t \rightarrow \infty$.

The first expression after taking the limit in (5.5) follows from the mapping Theorem 2.3 and the second from applying T_j to (5.2) and then using Fubini to hold the integration with respect to Lebesgue measure L outside as an expectation.

Proof. Here is the outline; more detail is given in the next section. We prove convergence using Theorem 2.1(iii). Take F and G closed and open sets respec-

tively in \mathbb{D} that are bounded away from $\mathbb{D}_{\leq j-1}$. Take $\delta > 0$ small enough so that also $F_\delta = \{x \in \mathbb{D} : d_{\text{sk}}(x, F) \leq \delta\}$ is bounded away from $\mathbb{D}_{\leq j-1}$. Then

$$\begin{aligned} tP[X/b(t^{1/j}) \in F] &= tP\left[X \in b(t^{1/j})F, \sup_{s \in [0,1]} |\tilde{X}_s| \leq b(t^{1/j})\delta\right] \\ &\quad + tP\left[X \in b(t^{1/j})F, \sup_{s \in [0,1]} |\tilde{X}_s| > b(t^{1/j})\delta\right] \\ &\leq tP[J \in b(t^{1/j})F_\delta] + tP\left[\sup_{s \in [0,1]} |\tilde{X}_s| > b(t^{1/j})\delta\right]. \end{aligned} \quad (5.6)$$

The Lévy process \tilde{X} has all moments finite and does not contribute asymptotically. Application of Lemmas 5.2 and 5.1, and letting $\delta \downarrow 0$ gives

$$\limsup_{t \rightarrow \infty} tP[X/b(t^{1/j}) \in F] \leq (\mu^{(j)} \times L) \circ T_j^{-1}(F).$$

To deal with the lower bound using open G , take $\delta > 0$ small enough so that

$$G^{-\delta} := ((G^c)_\delta)^c = \{x \in G : d_{\text{sk}}(x, y) < \delta \text{ implies } y \in G\}$$

is nonempty and bounded away from $\mathbb{D}_{\leq j-1}$. Then

$$\begin{aligned} tP[X/b(t^{1/j}) \in G] &\geq tP\left[J \in b(t^{1/j})G^{-\delta}, \sup_{s \in [0,1]} |\tilde{X}_s| \leq b(t^{1/j})\delta\right] \\ &= tP\left[J \in b(t^{1/j})G^{-\delta}\right] P\left[\sup_{s \in [0,1]} |\tilde{X}_s| \leq b(t^{1/j})\delta\right]. \end{aligned}$$

Applying Lemmas 5.2 and 5.1 and letting $\delta \downarrow 0$ gives

$$\liminf_{t \rightarrow \infty} tP[X/b(t^{1/j}) \in G] \geq (\mu^{(j)} \times L) \circ T_j^{-1}(G). \quad \square$$

5.1. Details

We now provide more detail for the proof of Theorem 5.1.

In the decomposition (5.1), the process \tilde{X} represents small jumps that should not affect asymptotics. We make this precise with the next Lemma.

Lemma 5.1. *For $j \geq 1$, and any $\delta > 0$,*

$$\limsup_{t \rightarrow \infty} tP\left[\sup_{s \in [0,1]} |\tilde{X}_s| > b(t^{1/j})\delta\right] = 0.$$

Proof. We rely on Skorohod's inequality for Lévy processes [9], [38, Section 7.3]. For $a > 0$,

$$P\left[\sup_{s \in [0,1]} |\tilde{X}_s| > 2a\right] \leq (1-c)^{-1} P[|\tilde{X}_1| > a],$$

where $c = \sup_{s \in [0,1]} P[|\tilde{X}_s| > a]$. Thus, since \tilde{X}_1 has all moments finite, for any $m > 1$,

$$tP\left[\sup_{s \in [0,1]} |\tilde{X}_s| > b(t^{1/j})\delta\right] \leq t(1-c(t))^{-1} P[|\tilde{X}_1| > b(t^{1/j})\delta/2]$$

$$\leq t(1 - c(t))^{-1} \frac{E|\tilde{X}_1|^m}{b^m(t^{1/j})(\delta/2)^m}.$$

For large enough m , $t/b^m(t^{1/j}) \rightarrow 0$ as $t \rightarrow \infty$ and

$$\begin{aligned} c(t) &:= \sup_{s \in [0,1]} P[|\tilde{X}_s| > b(t^{1/j})\delta/2] \leq \sup_{s \in [0,1]} \frac{E|\tilde{X}_s|^m}{b^m(t^{1/j})(\delta/2)^m} \\ &= \sup_{s \in [0,1]} \frac{s^m E|\tilde{X}_1|^m}{b^m(t^{1/j})(\delta/2)^m} \leq \frac{E|\tilde{X}_1|^m}{b^m(t^{1/j})(\delta/2)^m} \rightarrow 0, \end{aligned}$$

as $t \rightarrow \infty$ since $b(t) \rightarrow \infty$. \square

Lemma 5.2. For $j \geq 1$, $tP[J \in b(t^{1/j}) \cdot] \rightarrow (\mu^{(j)} \times L) \circ T_j^{-1}(\cdot)$ in $\mathbb{M}(\mathbb{D} \setminus \mathbb{D}_{\leq j-1})$ as $t \rightarrow \infty$.

Proof. We apply Theorem 2.1(iii).

CONSTRUCTION OF THE LOWER BOUND FOR OPEN SETS: Let $G \subset \mathbb{D}$ be open and bounded away from $\mathbb{D}_{\leq j-1}$. This implies that functions in G have no fewer than j jumps. Recall that $\Gamma_l = E_1 + \dots + E_l$, where the E_{ks} are iid standard exponentials. Take $M \geq j$ and notice that

$$\begin{aligned} &tP\left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l)1_{[U_l,1]} \in b(t^{1/j})G\right] \\ &\geq tP\left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l)1_{[U_l,1]} \in b(t^{1/j})G, N_1 \leq M\right] \\ &= tP\left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l)1_{[U_l,1]} \in b(t^{1/j})G, j \leq N_1 \leq M\right] \\ &\geq tP\left[\sum_{l=1}^j Q^{\leftarrow}(\Gamma_l)1_{[U_l,1]} \in b(t^{1/j})G^\delta, \sum_{l=j+1}^M Q^{\leftarrow}(\Gamma_l) \leq b(t^{1/j})\delta, Q^{\leftarrow}(\Gamma_{M+1}) < 1\right] \\ &\geq tP\left[\sum_{l=1}^j \frac{Q^{\leftarrow}(\Gamma_l)}{b(t^{1/j})} 1_{[U_l,1]} \in G^\delta, M \frac{Q^{\leftarrow}(E_{j+1})}{b(t^{1/j})} \leq \delta, Q^{\leftarrow}(\Gamma_{M+1} - \Gamma_{j+1}) < 1\right] \\ &\geq tP\left[\sum_{l=1}^j \frac{Q^{\leftarrow}(\Gamma_l)}{b(t^{1/j})} 1_{[U_l,1]} \in b(t^{1/j})G^\delta\right] P\left[MQ^{\leftarrow}(E_{j+1}) \leq b(t^{1/j})\delta\right] \\ &\quad \times P\left[Q^{\leftarrow}(\Gamma_{M+1} - \Gamma_{j+1}) < 1\right] \\ &\geq tP\left[\left(\left(\frac{Q^{\leftarrow}(\Gamma_l)}{b(t^{1/j})}, l \geq 1\right), (U_l, l \geq 1)\right) \in T_j^{-1}(G^\delta)\right] P\left[MQ^{\leftarrow}(E_{j+1}) \leq b(t^{1/j})\delta\right] \\ &\quad \times P\left[Q^{\leftarrow}(\Gamma_{M+1} - \Gamma_{j+1}) < 1\right]. \end{aligned}$$

Let $t \rightarrow \infty$ and apply Theorem 2.1(iii) to (5.2), so the lim inf of the first factor above has a lower bound. As $t \rightarrow \infty$, the second factor approaches 1. Let $M \rightarrow \infty$

and the third factor also approaches 1. Let $\delta \downarrow 0$ and we obtain

$$\liminf_{t \rightarrow \infty} tP \left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})G \right] \geq (\mu^{(j)} \times L) \circ T_j^{-1}(G).$$

CONSTRUCTION OF THE UPPER BOUND FOR CLOSED SETS: Let $F \subset \mathbb{D}$ be closed and bounded away from $\mathbb{D}_{\leq j-1}$. Take $\beta \in (0, 1)$ close to 1 and let

$$M_t = \sum_{l=1}^{N_1} 1_{(b(t^{1/j})^\beta, \infty)}(Q^{\leftarrow}(\Gamma_l)).$$

Choose $\delta > 0$ small enough so that $F_\delta := \{x \in \mathbb{D} : d(x, F) \leq \delta\}$ is bounded away from $\mathbb{D}_{\leq j-1}$. Then

$$\begin{aligned} & tP \left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})F \right] \\ &= tP \left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})F, \sum_{l=M_t+1}^{N_1} Q^{\leftarrow}(\Gamma_l) \leq b(t^{1/j})\delta \right] \\ &\quad + tP \left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})F, \sum_{l=M_t+1}^{N_1} Q^{\leftarrow}(\Gamma_l) > b(t^{1/j})\delta \right] \\ &\leq tP \left[\sum_{l=1}^{M_t} Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})F_\delta \right] + tP \left[\sum_{l=M_t+1}^{N_1} Q^{\leftarrow}(\Gamma_l) > b(t^{1/j})\delta \right]. \end{aligned}$$

Decompose the first summand according to whether $M_t \leq j$ or $M_t \geq j + 1$. Notice $M_t < j$ is incompatible with $\sum_{l=1}^{M_t} Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})F_\delta$ since F_δ is bounded away from $\mathbb{D}_{\leq j-1}$. Thus we get the upper bound

$$\begin{aligned} &\leq tP \left[\sum_{l=1}^j Q^{\leftarrow}(\Gamma_l) 1_{[U_l, 1]} \in b(t^{1/j})F_\delta \right] + tP[M_t \geq j + 1] \\ &\quad + tP \left[\sum_{l=M_t+1}^{N_1} Q^{\leftarrow}(\Gamma_l) > b(t^{1/j})\delta \right]. \end{aligned}$$

We now show that the second and third of the three terms above vanish as $t \rightarrow \infty$. Firstly, the definition of M_t implies that $Q^{\leftarrow}(\Gamma_l) \leq b(t^{1/j})^\beta$ for $M_t + 1 \leq l \leq N_1$. Thus,

$$\begin{aligned} tP \left[\sum_{l=M_t+1}^{N_1} Q^{\leftarrow}(\Gamma_l) > b(t^{1/j})\delta \right] &\leq tP[(N_1 - M_t)b(t^{1/j})^\beta > b(t^{1/j})\delta] \\ &\leq tP[N_1 > b(t^{1/j})^{1-\beta}\delta]. \end{aligned}$$

The right-hand side converges to 0 as $t \rightarrow \infty$ since the tail probability has a Markov bound of $tE(N_1^p)/[b(t^{1/j})^{1-\beta}\delta]^p$ for any p . Secondly,

$$\begin{aligned} P[M_t \geq j+1] &\leq P[Q^{\leftarrow}(\Gamma_{j+1}) > b(t^{1/j})^\beta] \\ &\leq P[\Gamma_{j+1} \leq \nu([b(t^{1/j})^\beta, \infty))] \\ &\leq P[\max(E_1, \dots, E_{j+1}) \leq \nu([b(t^{1/j})^\beta, \infty))] \\ &= P[E_1 \leq \nu([b(t^{1/j})^\beta, \infty))]^{j+1}. \end{aligned}$$

Since $P[E_1 \leq y] \sim y$ as $y \downarrow 0$, and since $\nu([x, \infty))$ is regularly varying at infinity with index $-\alpha$ and b is regularly varying at infinity with index $1/\alpha$, we find that

$$\limsup_{t \rightarrow \infty} t\nu([b(t^{1/j})^\beta, \infty))^{j+1} = \limsup_{t \rightarrow \infty} L(t)t^{1-\beta(j+1)/j}$$

for some slowly varying function L . In particular, choosing $\beta \in (\frac{j}{j+1}, 1)$ ensures that $\lim_{t \rightarrow \infty} tP[M_t \geq j+1] = 0$.

We now deal with the remaining term. Since

$$\begin{aligned} tP\left[\sum_{l=1}^j Q^{\leftarrow}(\Gamma_l)1_{[U_l, 1]} \in b(t^{1/j})F_\delta\right] \\ = tP\left[\left((Q^{\leftarrow}(\Gamma_l), l \geq 1), (U_l, l \geq 1)\right) \in b(t^{1/j}) \circ T_j^{-1}(F_\delta)\right] \\ + tP\left[\left((Q^{\leftarrow}(\Gamma_l), l \geq 1), (U_l, l \geq 1)\right) \in A_m^c\right] \end{aligned}$$

and, by Lemmas 5.3 and 5.4, $T_j^{-1}(F_\delta)$ is, if nonempty, closed and bounded away from $\mathbb{H}_{\leq j-1} \times [0, 1]^\infty$, Proposition 5.1 and the fact that A_m^c is a P -null set yield that

$$\limsup_{t \rightarrow \infty} tP\left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l)1_{[U_l, 1]} \in b(t^{1/j})F\right] \leq (\mu^{(j)} \times L) \circ T_j^{-1}(F_\delta).$$

Letting $\delta \downarrow 0$ shows that

$$\limsup_{t \rightarrow \infty} tP\left[\sum_{l=1}^{N_1} Q^{\leftarrow}(\Gamma_l)1_{[U_l, 1]} \in b(t^{1/j})F\right] \leq (\mu^{(j)} \times L) \circ T_j^{-1}(F).$$

We have thus shown that $\liminf_{t \rightarrow \infty} tP[J \in b(t^{1/j})G] \geq (\mu^{(j)} \times L_j) \circ T_j^{-1}(G)$ and $\limsup_{t \rightarrow \infty} tP[J \in b(t^{1/j})F] \leq (\mu^{(j)} \times L_j) \circ T_j^{-1}(F)$ for all open G and closed F bounded away from $\mathbb{D}_{\leq j-1}$. The conclusion follows from Theorem 2.1. \square

Recall the definitions of A_m and T_m in (5.3) and (5.4).

Lemma 5.3. For $m \geq 1$, $T_m : A_m \mapsto \mathbb{D}$ is continuous.

Proof. The projection

$$A_m \ni (x, u) \mapsto ((x_1, \dots, x_m), (u_1, \dots, u_m)) \in \mathbb{R}_+^{m\downarrow} \times (0, 1)^{m, \neq}$$

where $(0, 1)^{m, \neq} = \{(u_1, \dots, u_m) \in (0, 1)^m : u_i \neq u_j \text{ for } i \neq j\}$, is continuous. Since compositions of continuous functions are continuous, it remains to check that

$$\mathbb{R}_+^{m\downarrow} \times (0, 1)^{m, \neq} \ni ((x_1, \dots, x_m), (u_1, \dots, u_m)) \mapsto \sum_{i=1}^m x_i 1_{[u_i, 1]} \in \mathbb{D}$$

is continuous. Take $(x, u) \in \mathbb{R}_+^{m\downarrow} \times (0, 1)^{m, \neq}$. Then there exists some $\delta > 0$ such that, for $(\tilde{x}, \tilde{u}) \in \mathbb{R}_+^{m\downarrow} \times (0, 1)^{m, \neq}$, $d_{2m}((x, u), (\tilde{x}, \tilde{u})) < \delta$, where d_{2m} is the usual metric in \mathbb{R}^{2m} , implies that the components of \tilde{u} appear in the same order as do the components of u . If $0 = u_{(0)} < u_{(1)} < \dots < u_{(m)} < u_{(m+1)} = 1$, with corresponding notation for the ordered \tilde{u} 's, make sure $3 \cdot \delta < \bigvee_{i=1}^{m+1} |u_{(i)} - u_{(i-1)}| \vee |\tilde{u}_{(i)} - \tilde{u}_{(i-1)}|$. Consider the piece-wise linear function λ_l for which $\lambda_l(0) = 0$, $\lambda_l(1) = 1$, and $\lambda_l(u_i) = \tilde{u}_i$ for each i . Notice that λ_l is strictly increasing and satisfies $\|\lambda_l - e\| < \delta$. Therefore,

$$\sup_{t \in [0, 1]} \left| \sum_{i=1}^m x_i 1_{[\lambda_l(u_i), 1]}(t) - \sum_{i=1}^m \tilde{x}_i 1_{[\tilde{u}_i, 1]}(t) \right| < \sum_{i=1}^m |x_i - \tilde{x}_i| < m\delta.$$

In particular,

$$d_{\text{sk}} \left(\sum_{i=1}^m x_i 1_{[u_i, 1]}, \sum_{i=1}^m \tilde{x}_i 1_{[\tilde{u}_i, 1]} \right) < m\delta,$$

which shows the continuity. \square

Lemma 5.4. *Suppose $A \subset \mathbb{D}$ is bounded away from $\mathbb{D}_{\leq j-1}$. For $m \geq j$, if $T_m^{-1}(A)$ is nonempty, then it is bounded away from $\mathbb{H}_{\leq j-1} \times [0, 1]^\infty$.*

Proof. If $A \cap \mathbb{D}_{\leq m} = \emptyset$, then $T_m^{-1}(A) = \emptyset$. Therefore, without loss of generality we may take $A \subset \mathbb{D}_{\leq m}$. Assume $d_{\text{sk}}(A, \mathbb{D}_{\leq j-1}) > \delta > 0$ and notice that $x \in \mathbb{D}_{\leq m}$ if and only if

$$x = \sum_{i=1}^m y_i 1_{[u_i, 1]} \text{ for } y_1 \geq \dots \geq y_m \geq 0, u_i \in [0, 1].$$

If $x \in A$, $\sum_{i=j}^m y_i > \delta$ as a consequence of $d_{\text{sk}}(A, \mathbb{D}_{\leq j-1}) > \delta$ and because the y 's are non-increasing, $y_j > \delta / (m - j + 1)$. Consequently,

$$T_m^{-1}(A) \subset \left\{ (x_i, i \geq 1) \in \mathbb{R}_+^{\infty\downarrow} : x_j > \delta / (m - j + 1) \right\} \times [0, 1]^\infty,$$

and the latter set is bounded away from $\mathbb{H}_{\leq j-1} \times [0, 1]^\infty$. \square

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