# Multifractal analysis of Lévy fields 

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Received: 17 May 2010 / Revised: 10 December 2010 / Published online: 1 February 2011
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#### Abstract

We study the pointwise regularity properties of the Lévy fields introduced by T. Mori; these fields are the most natural generalization of Lévy processes to the multivariate setting. We determine their spectrum of singularities, and we show that their Hölder singularity sets satisfy a large intersection property in the sense of K. Falconer.


Keywords Lévy random fields • Multifractal analysis • Hausdorff measures and dimension • Sets with large intersection • Diophantine approximation • Ubiquity

Mathematics Subject Classification (2000) Primary: 60G60 • 60G51 • Secondary: 60G17 • 60D05 • 28A78 • 28A80

## 1 Introduction

The determination of the uniform and pointwise regularity of stochastic processes has been a longstanding issue, starting with the discovery of the uniform modulus of continuity of Brownian motion by Lévy and Wiener, the law of the iterated logarithm by Khintchine and Kolmogorov, and the study of the irregularity of the Brownian paths by Paley, Wiener and Zygmund, as well as Dvoretzky.

[^0]Remarkably, a connexion between stochastic processes and fractals was also first established in the case of Brownian motion: slow points where the modulus of continuity is smaller than almost everywhere (the $\sqrt{\log \log \left(1 /\left|t-t_{0}\right|\right)}$ term is replaced by a constant $a>0$ ), and fast points (where this term is replaced by $a \sqrt{\log \left(1 /\left|t-t_{0}\right|\right)}$ ) were shown to occur on random fractal sets; the dimensions of these collections of sets (indexed by the parameter $a$ ) were determined by E. Perkins and by S. Orey and S.J. Taylor, respectively. This connexion did not remain confined to the specific case of Brownian motion: in the mid 1980s, multifractal analysis was introduced in order to propose a general framework for the study of the local variations of the regularity of stochastic processes, using the mathematical tools supplied by fractal analysis. Let us be more specific.
Definition 1 Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a locally bounded function, $t_{0} \in \mathbb{R}^{d}$ and $\alpha>0$. The function $f$ belongs to $C^{\alpha}\left(t_{0}\right)$ if there exist $C>0$ and a polynomial $P_{t_{0}}$ of degree less than $\alpha$ such that, for all $t$ in a neighborhood of $t_{0}$,

$$
\left|f(t)-P_{t_{0}}(t)\right| \leq C\left\|t-t_{0}\right\|^{\alpha}
$$

The Hölder exponent of $f$ at $t_{0}$ is

$$
\alpha_{f}\left(t_{0}\right)=\sup \left\{\alpha \geq 0 \mid f \in C^{\alpha}\left(t_{0}\right)\right\}
$$

Note that $\alpha_{f}$ takes values in $[0, \infty]$. Multifractal analysis is concerned with the determination of the Hausdorff dimension of the set of points where the Hölder exponent takes a given value. Two collections of sets play a particular role: the iso-Hölder sets of $f$, defined by

$$
E_{f}(h)=\left\{t \in \mathbb{R}^{d} \mid \alpha_{f}(t)=h\right\}
$$

and the singularity sets $E_{f}^{\prime}(h)$ which consist of the points $t$ where $f$ is continuous and satisfies $\alpha_{f}(t) \leq h$. Note that, though these sets can be studied for any random field, they are really pertinent only when the Hölder exponent is nonconstant (which excludes for instance Brownian motion, where it takes everywhere the value $1 / 2$ ). If the Hölder exponent changes from point to point, one is interested in determining the Hausdorff dimension of the iso-Hölder sets $E_{f}(h)$. The corresponding notion is supplied by the local spectrum of singularities defined as follows.

Definition 2 Let $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a locally bounded function and let $W$ be a nonempty open subset of $\mathbb{R}^{d}$. The local spectrum of singularities of $f$ is

$$
d_{f}(h, W)=\operatorname{dim}_{\mathrm{H}}\left(E_{f}(h) \cap W\right),
$$

where $\operatorname{dim}_{H}$ denotes Hausdorff dimension (with the convention that $\operatorname{dim}_{H} \emptyset=-\infty$ ).
For many examples of random fields $Y$, the spectrum of singularities does not depend on the particular region $W$ that is considered, and is actually a deterministic function. If such is the case, i.e. if there is a deterministic function $d_{Y}(h)$ such that

$$
\text { a.s. } \forall h \in[0, \infty] \quad \forall W \neq \emptyset \text { open } \quad d_{Y}(h, W)=d_{Y}(h),
$$

the random field $Y$ is called homogeneous (following the terminology of [25]). In that case, the support of the spectrum is the set of $h \in[0, \infty]$ such that $d_{Y}(h) \geq 0$ or, equivalently, such that $\alpha_{f}(t)=h$ for some $t \in \mathbb{R}^{d}$. A spectrum is degenerate if its support is reduced to a single point. This means that the Hölder exponent takes the same value everywhere (as is the case for Brownian motion). In that situation, the field is called a mono-Hölder field. Otherwise, it is called multifractal.

As shown in [18,24], Lévy processes supply examples of multifractal processes which are homogeneous. Note that this is not the case for all classical random processes, see e.g. Barral et al. [7] for examples of Markov processes whose spectrum is random and depends on the region, or Durand [15] for a model of wavelet series based on a tree-indexed Markov chain whose spectrum is random too. Moreover, though many processes with a fairly general spectrum have been constructed (such as the random processes built by Ayache et al. [6]), the only natural large class of homogeneous random processes which has been studied is that of Lévy processes.

The introduction of multifractal analysis was motivated by classification and model selection issues in signal processing (the analysis of one-dimensional turbulence data), and its fast development was boosted by its pertinence for an unexpectedly large number of applications. Recent developments in image processing have stimulated similar needs in 2D, and classification tools based on multifractal parameters have recently proven to be promising directions of research [43]. These results motivate the investigation of random fields that could be used in image modeling, and the study of their multifractal properties (and in particular for porous media, medical imaging, etc.). Since Lévy processes were the first class of processes proven to be multifractal, and as they play an important role in both theoretical probability and modeling, it is very natural to wonder if their multivariate extensions display similar properties.

We shall prove that it is the case for the fields that we study in this paper, which are multivariate analogs of Lévy processes that jump along random hyperplanes. Due to that very anisotropic situation, new ubiquity techniques will have to be employed in order to perform their multifractal analysis. This strongly differs from other situations where geometry does not play any specific role and the results in the one-dimensional case may straightforwardly be extended to higher dimensions. Such is the case for several multifractal models of random wavelet series [3,4,18].

Two main extensions of Lévy processes to the multivariate setting have been proposed. The first one is that of Adler et al. [1], which finds its origins in the work of Straf [41], and whose regularity was studied by Vares [42] and Lagaize [28] in the two-dimensional case. The second extension is due to Mori [33], and covers the following important particular cases. In the Gaussian case, the Lévy Brownian motion [30] is the first instance of Lévy field (in the sense of Mori) which has been considered. This field is isotropic, but its geometric construction, due to Chentsov [14], may easily be extended to the anisotropic setting. Subsequently, Chentsov's construction motivated the introduction of a class of isotropic and stable Lévy fields, which also satisfy a remarkable selfsimilarity property. These fields, usually termed as Lévy-Chentsov fields, are discussed in Samorodnitsky and Taqqu's book [37], and were studied by Shieh [39], who established their local boundedness and the existence of local times.

In order to compare the relevance of these models, we propose two natural criteria for selecting a multivariate extension of Lévy processes:

- Stability under linear transforms of coordinates: If $M$ is an invertible (deterministic) linear transform, and $Y(t)$ is a Lévy field, then $Y(M t)$ should also be a Lévy field. In particular, the coordinate axes do not play any specific role in the construction or the analysis of these fields.
- Stability under trace: The restriction of a $d$-dimensional Lévy field to any $d^{\prime}$-dimensional linear subspace is a $d^{\prime}$-dimensional Lévy field. Note that the first requirement implies that we do not need to specify a particular system of coordinates on that subspace. Furthermore, restrictions to arbitrary half lines are Lévy processes.

Among the aforementioned extensions, Mori's is the only one that fulfills these two criteria. It turns out that Mori's definition of a Lévy field is also a natural multivariate analog of the definition of a Lévy process. Indeed, recall that a stochastic process indexed by $[0, \infty)$ is a Lévy process if it satisfies the following properties: it has stationary and independent increments, it is stochastically continuous, and it vanishes at zero almost surely, see e.g. [12,38]. This implies in particular that the one-dimensional marginals are infinitely divisible.

Definition 3 A random field $Y=\left\{Y(t), t \in \mathbb{R}^{d}\right\}$ is a Lévy field (in the sense of Mori) if the following conditions are satisfied:

1. It is stochastically continuous and vanishes at zero almost surely.
2. It has stationary increments, i.e. $Y(a+\cdot)-Y(a) \stackrel{d}{=} Y$ for any $a \in \mathbb{R}^{d}$.
3. Its finite-dimensional marginals are infinitely divisible.
4. For any $a, b \in \mathbb{R}^{d}$, the increments of $\{Y(a+\lambda b), \lambda \in \mathbb{R}\}$ are independent.

In the previous definition, $\stackrel{d}{=}$ stands for equality of the finite dimensional distributions. It is easy to see that for any $a, b \in \mathbb{R}^{d}$, the process $\{Y(a+\lambda b)-Y(a), \lambda \geq 0\}$ has stationary and independent increments, is stochastically continuous, and vanishes at zero almost surely, thereby being a Lévy process. More generally, it is clear that the two stability requirements listed above are satisfied.

## 2 Main results

### 2.1 Representation of Lévy fields

Mori established a remarkable decomposition of Lévy fields into three parts which is similar to the Lévy-Itō decomposition of a Lévy process into a linear drift, a Brownian motion, and a jump component controlled by the Lévy measure, see e.g. [38, Chap. 4]. In the $d$-dimensional case, the three components are the following:

1. A linear drift, that is, a mapping of the form $t \mapsto\langle a, t\rangle$, where $a$ is a deterministic vector and $\langle\cdot, \cdot\rangle$ denotes the standard inner product in $\mathbb{R}^{d}$.
2. A Gaussian field $B_{\mu}=\left\{B_{\mu}(t), t \in \mathbb{R}^{d}\right\}$ that depends on a finite nonnegative symmetric Borel measure $\mu$ defined on the unit sphere $\mathbb{S}^{d-1}$ of $\mathbb{R}^{d}$. Here, symmetric means invariant under the mapping $s \mapsto-s$. The construction of $B_{\mu}$ is detailed in Sect. 3. Note that $B_{\mu}$ is almost surely constant equal to zero when $\mu=0$.
3. A jump field $L_{v}=\left\{L_{v}(t), t \in \mathbb{R}^{d}\right\}$. Recall that the jump component of a Lévy process is a sum of independent compensated compound Poisson processes; the structure is the same in the multivariate setting, except that Poisson processes, which jump at points, are replaced by random fields which jump along Poisson distributed hyperplanes. Their distribution is described by a $d$-dimensional analog of the Lévy measure, namely, a nonnegative Borel measure $\nu$ defined on $\mathbb{S}^{d-1} \times \mathbb{R}^{*}$, where $\mathbb{R}^{*}$ means $\mathbb{R} \backslash\{0\}$. The measure $v$ is symmetric, i.e. invariant under $(s, x) \mapsto(-s,-x)$, and

$$
\begin{equation*}
\int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in \mathbb{R}^{*}}}\left(1 \wedge x^{2}\right) v(\mathrm{~d} s, \mathrm{~d} x)<\infty \tag{1}
\end{equation*}
$$

where $\wedge$ stands for minimum. Intuitively, $v(\mathrm{~d} s, \mathrm{~d} x)$ describes the amount of hyperplanes orthogonal to $s$ where a jump of size $x$ occurs. In particular, the jump field $L_{v}$ is almost surely zero everywhere when $v=0$. The construction of $L_{v}$ is detailed in Sect. 4.

Theorem 1 (Mori) Every Lévy field $Y=\left\{Y(t), t \in \mathbb{R}^{d}\right\}$ may be represented in the following manner:

$$
Y \stackrel{d}{=}\langle a, \cdot\rangle+B_{\mu}+L_{v}
$$

where $a \in \mathbb{R}^{d}, \mu$ is a finite symmetric measure on $\mathbb{S}^{d-1}$, and $\nu$ is a symmetric measure on $\mathbb{S}^{d-1} \times \mathbb{R}^{*}$ satisfying (1). Moreover, $B_{\mu}$ and $L_{v}$ are independent and the triple $(a, \mu, \nu)$ is uniquely determined by the field $Y$.

Due to Theorem 1, studying the regularity of a Lévy field reduces to analyzing each of these three components, and then to understanding what happens when combining them. To begin with, note that linear drifts are $C^{\infty}$ and thus play no role in our analysis.

### 2.2 Regularity of the Gaussian component

Our main result on the regularity of the Gaussian field $B_{\mu}$ is the next statement.
Theorem 2 Let $\mu$ be a finite symmetric measure on $\mathbb{S}^{d-1}$ such that $\mu \neq 0$. Then, $B_{\mu}$ is a homogeneous mono-Hölder field and

$$
\text { a.s. } \forall t \in \mathbb{R}^{d} \quad \alpha_{B_{\mu}}(t)=1 / 2 .
$$

This theorem is in fact a consequence of two slightly more precise results on the pointwise modulus of continuity of $B_{\mu}$, namely, Propositions 2 and 3 below. Their
proofs are given in Sect. 5 and rely on standard techniques for studying the oscillations of Gaussian random fields.

### 2.3 Regularity of the jump component

Precise results on the size of the iso-Hölder sets of $L_{v}$ are detailed in Sect. 6. In order to give a first insight into them, let us single out a representative consequence: Theorem 3 below, which yields the spectrum of singularities of $L_{\nu}$. Its statement involves an index $\beta_{\nu}$ associated with $v$ and holds under an admissibility condition bearing on that measure.

Definition 4 Let $v$ be a symmetric measure on $\mathbb{S}^{d-1} \times \mathbb{R}^{*}$ satisfying (1). The index of $v$ is

$$
\begin{equation*}
\beta_{\nu}=\inf \left\{\gamma \geq 0 \mid \int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in(0,1]}} x^{\gamma} \nu(\mathrm{d} s, \mathrm{~d} x)<\infty\right\} . \tag{2}
\end{equation*}
$$

For any integer $j \geq 1$, let $v_{j}=v\left(\mathbb{S}^{d-1} \times\left(2^{-j}, 2^{-j+1}\right]\right)$; $v$ is admissible if

$$
\begin{equation*}
\chi_{\nu}=\sum_{j=1}^{\infty} 2^{-j}\left(j v_{j}\right)^{1 / 2}<\infty \tag{3}
\end{equation*}
$$

The index $\beta_{\nu}$ is the $d$-dimensional analog of the index that Blumenthal and Getoor associated to a Lévy process [13, Definition 2.1]; note that (1) implies that $\beta_{\nu} \in[0,2]$. We shall almost always assume below that the measure $\nu$ is admissible. A similar assumption had been made in [24]. This condition is slightly stronger than (1), which ensures the existence of $L_{v}$ and amounts to the finiteness of $\sum 2^{-2 j} v_{j}$. Actually, assuming (3) is a mild restriction since, for instance, every measure $v$ with index less than two is admissible.

Theorem 3 Let $v$ be an admissible measure with $\beta_{v}>0$. Then, $L_{v}$ is a homogeneous multifractal field and with probability one,

$$
\forall h \in[0, \infty] \quad \forall W \neq \emptyset \text { open } \quad d_{L_{v}}(h, W)= \begin{cases}d-1+\beta_{v} h & \text { if } h \leq 1 / \beta_{v} \\ -\infty & \text { if } h>1 / \beta_{v}\end{cases}
$$

This theorem is a straightforward consequence of Proposition 4 and Corollary 2 below. As well as the results of Sect. 6, it covers the isotropic and stable case of the Lévy-Chentsov random fields [37], for which $v(\mathrm{~d} s, \mathrm{~d} x)=\sigma(\mathrm{d} s) \mathrm{d} x /|x|^{\alpha+1}$ where $\alpha \in(0,2)$ and $\sigma$ denotes the uniform measure on $\mathbb{S}^{d-1}$. In that situation, $\beta_{\nu}=\alpha$.

Contrary to what the hypothesis of Theorem 3 may suggest, the analysis developed in Sect. 6 also includes the case in which $\beta_{\nu}=0$. In particular, we shall discuss
the regularity of compound Poisson fields (multivariate analogs of compound Poisson processes), for which $v$ has finite total mass. In that situation, the field is not homogeneous. Indeed, locally, the field jumps on at most finitely many hyperplanes, thereby being $C^{\infty}$ except on a set of dimension $d-1$ where its Hölder exponent vanishes, see Proposition 5 for a precise statement.

### 2.4 Spectrum of singularities of a general Lévy field

We call a Lévy field canonical if it is of the form $Y_{a, \mu, \nu}=\langle a, \cdot\rangle+B_{\mu}+L_{\nu}$ for some $a \in \mathbb{R}^{d}$, some finite symmetric measure $\mu$ on $\mathbb{S}^{d-1}$, and some symmetric measure $v$ on $\mathbb{S}^{d-1} \times \mathbb{R}^{*}$ satisfying (1). By virtue of Theorem 1 , every Lévy field has the same finite-dimensional distributions as a canonical one. In view of Definition 4, we call $\beta_{\nu}$ the index of $Y_{a, \mu, \nu}$, and we call that field admissible if $v$ is admissible. Moreover, $Y_{a, \mu, \nu}$ is said to have a Gaussian component if and only if $\mu \neq 0$.

We show in Sect. 7 that the Hölder exponent of $Y_{a, \mu, \nu}$ is everywhere the minimum of that of its components $B_{\mu}$ and $L_{\nu}$. In view of Theorems 2 and 3, this leads immediately to the next statement. This approach also easily enables one to deduce comparable results for the cases that are treated in Sect. 6 but not covered by Theorem 3. For instance, one could easily infer the spectrum of singularities of a canonical Lévy field whose jump component is a compound Poisson field.

Corollary 1 Let $Y$ be an admissible canonical Lévy field with index $\beta>0$. Then, $Y$ is homogeneous. Moreover,

- if Y has no Gaussian component, then with probability one,

$$
\forall h \in[0, \infty] \quad \forall W \neq \emptyset \text { open } \quad d_{Y}(h, W)= \begin{cases}d-1+\beta h & \text { if } h \leq 1 / \beta \\ -\infty & \text { if } h>1 / \beta\end{cases}
$$

- if Y has a Gaussian component, then with probability one,

$$
\forall h \in[0, \infty] \quad \forall W \neq \emptyset \text { open } \quad d_{Y}(h, W)= \begin{cases}d-1+\beta h & \text { if } h<1 / 2 \\ d & \text { if } h=1 / 2 \\ -\infty & \text { if } h>1 / 2\end{cases}
$$

As expected, when $d=1$, the previous result boils down to [24, Theorem 1] which gives the spectrum of singularities of a Lévy process.

Remark 1 The notion of pointwise regularity given in Definition 1 does not take into account directional regularity but yields the worst possible regularity in all directions. Therefore, the results obtained in this paper do not take into account possible directional irregularity phenomena. However, such phenomena are to be expected in the case of Lévy fields. Indeed, given that they display jumps along hyperplanes, they have, by construction, a very anisotropic nature. Therefore, it would be of great interest to perform a multifractal analysis of these fields using a more flexible notion of pointwise smoothness, which can take directionality into account (see e.g. [10,26] for appropriate definitions).

### 2.5 Large intersection properties and ubiquity

We will not restrict our study to the determination of the Hausdorff dimension of the random sets of points related with the definition of Hölder regularity, but we will also investigate some of their geometric properties. It turns out that, rather than the iso-Hölder sets, the singularity sets mentioned in Sect. 1 display the most striking features. We shall show in Sect. 6.3 that the singularity sets of the jump component $L_{v}$ satisfy a remarkable counterintuitive property introduced by Falconer [21]: they fall in the category of sets with large intersection. Recall that the intersection of two subsets of $\mathbb{R}^{d}$ with dimension $d_{1}$ and $d_{2}$, respectively, is usually expected to be equal to $d_{1}+d_{2}-d$ (the codimensions add up), as in the case of affine subspaces, see [22, Chap. 8] for precise statements. Sets with large intersection disprove this heuristic in a striking way: their size properties are not altered by taking countable intersection. Indeed, the Hausdorff dimension of the intersection of countably many sets with large intersection is equal to the infimum of their Hausdorff dimensions. Sets with large intersection play a prominent role in metric number theory (Diophantine approximation) and dynamical systems, see e.g. $[16,17,21]$ and references therein. In the context of probability, they are relevant to the multifractal analysis of random wavelet series [4, 15], as well as the study of random coverings of the circle [19]. Durand [18] also proved that the singularity sets of Lévy processes are sets with large intersection. In all these situations, large intersection properties arise because the sets under study are derived from an underlying ubiquitous system. It is also the case here, and our proofs will require new extensions of ubiquity (see Theorem 7 below and its proof). Besides the aforementioned papers, we also refer to [8,9,11] for important results on ubiquity and its applications.

### 2.6 Behavior of traces

An important and difficult subject of investigation is to understand how the multifractal properties of a field and its traces on linear subspaces are related, see e.g. [5]. This question initially came up in the context of the analysis of turbulence: the only high precision experimental data available are one-dimensional cuts and the challenge is to infer from these cuts information about the multifractal properties of the whole field. In general, it is expected that the spectrum of the trace is the initial spectrum lowered by the codimension of the subspace (see (4) below) and the parts which become negative are set to $-\infty$. Lévy fields provide a case study of that effect, because the characteristic parameters of their traces are easy to obtain.

To be more specific, let $e=\left(e_{1}, \ldots, e_{d^{\prime}}\right)$ be an arbitrary orthonormal set of vectors of $\mathbb{R}^{d}$, with $1 \leq d^{\prime} \leq d$, and let $Y_{a, v, \mu}$ denote a canonical Lévy field. Then, the random field $Y_{a, v, \mu}^{e}$ defined by

$$
\forall t_{1}, \ldots, t_{d^{\prime}} \in \mathbb{R}^{d^{\prime}} \quad Y_{a, v, \mu}^{e}\left(t_{1}, \ldots, t_{d^{\prime}}\right)=Y_{a, v, \mu}\left(t_{1} e_{1}+\cdots+t_{d^{\prime}} e_{d^{\prime}}\right)
$$

is a canonical Lévy field indexed by $\mathbb{R}^{d^{\prime}}$ whose characteristic triple ( $a_{e}, \mu_{e}, v_{e}$ ) may be deduced from $(a, \mu, v)$ with the help of the mapping $\mathrm{p}_{e}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d^{\prime}}$ defined by
$\mathrm{p}_{e}(t)=\left(\left\langle t, e_{1}\right\rangle, \ldots,\left\langle t, e_{d^{\prime}}\right\rangle\right)$ for any $t \in \mathbb{R}^{d}$. To be specific, $a_{e}=\mathrm{p}_{e}(a)$, the measure $\mu_{e}$ is the image under $s \mapsto \mathrm{p}_{e}(s) /\left\|\mathrm{p}_{e}(s)\right\|$ of $\left\|\mathrm{p}_{e}(s)\right\| \mu(\mathrm{d} s)$, and the measure $v_{e}$ is the image under $(s, x) \mapsto\left(\mathrm{p}_{e}(s) /\left\|\mathrm{p}_{e}(s)\right\|, x\right)$ of $\left\|\mathrm{p}_{e}(s)\right\| \nu(\mathrm{d} s, \mathrm{~d} x)$. In particular, the field $Y_{a, v, \mu}^{e}$ has no Gaussian component if and only if $\mu$ is supported in the orthogonal complement of the linear span of $e$. Moreover, $v_{e}$ is admissible whenever $v$ is, and the index of $v_{e}$ is at most the index of $v$.

In the isotropic case, for which $v$ is the product of the uniform measure on $\mathbb{S}^{d-1}$ and a given measure on $\mathbb{R}^{*}$, the indices of $v_{e}$ and $v$ coincide, regardless of the choice of $e$. It follows from Corollary 1 that, as expected,

$$
\begin{equation*}
d_{Y_{a, v, \mu}^{e}}=d_{Y_{a, v, \mu}}-\left(d-d^{\prime}\right) \tag{4}
\end{equation*}
$$

However, for appropriate anisotropic choices of the measure $v$, one may obtain a whole range of values for the index $\beta_{\nu_{e}}$. Our results then lead to a whole variety of spectra for $Y_{a, \nu, \mu}^{e}$, depending on the choice of the directions for the trace.

### 2.7 Roadmap

The paper is organized as follows. In Sects. 3 and 4, we detail the construction of the two main components appearing in Mori's decompositon of Lévy fields: the Gaussian part and the jump part. We will also derive some basic properties which will be useful for their multifractal analysis.

Section 5 is devoted to the proof of Theorem 2, according to which the Gaussian part has everywhere the Hölder exponent $1 / 2$. Precise results on the size (in terms of Hausdorff measures, Hausdorff dimension and packing dimension) and large intersection properties of the iso-Hölder and the singularity sets of the jump part are stated in Sect. 6. These results lead to Theorem 3 above. In Sect. 7, we explain how the results on the Gaussian and the jump components may be combined to obtain the pointwise regularity of a general canonical Lévy field.

The rest of the paper is devoted to establishing the results of Sect. 6. The structure of the proof is described in Sect. 8; we present there the required tools, specifically, a precise knowledge of the location of the singularities of the jump part, and a description of the size and large intersection properties of the set of points approximated at a certain rate by random hyperplanes that are distributed in a Poissonian way. The first tool is detailed in Sect. 9. The second one is presented in Sects. 10 and 11, and relies heavily on ubiquity. Last, Sect. 12 details the proofs of the results of Sect. 6, and the paper ends with the proof of a lemma which is called upon by the first ingredient, see Sect. 13.

## 3 The Gaussian component

The Gaussian part of a Lévy field is a Gaussian field $B_{\mu}=\left\{B_{\mu}(t), t \in \mathbb{R}^{d}\right\}$ depending on a finite nonnegative symmetric Borel measure $\mu$ defined on $\mathbb{S}^{d-1}$. We now recall the construction of such random fields $B_{\mu}$, which is essentially due to Chentsov [14].

Note that, in the isotropic case (where $\mu$ is the uniform measure on the sphere), one basically ends up with a geometric representation of the Lévy Brownian motion [30].

### 3.1 Definition of $B_{\mu}$

Let us consider the collection $\mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$ of all relatively compact Borel subsets of the product space $\mathbb{H}_{d}=(0, \infty) \times \mathbb{S}^{d-1}$, and the centered Gaussian process $\mathfrak{B}_{\mu}=$ $\left\{\mathfrak{B}_{\mu}(V), V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)\right\}$ with covariance function given by

$$
\forall V, V^{\prime} \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right) \quad \mathbb{E}\left[\mathfrak{B}_{\mu}(V) \mathfrak{B}_{\mu}\left(V^{\prime}\right)\right]=\int_{(\rho, s) \in \mathbb{H}_{d}} \mathbb{1}_{\left\{(\rho, s) \in V \cap V^{\prime}\right\}} \mathrm{d} \rho \mu(\mathrm{~d} s)
$$

Such a process $\mathfrak{B}_{\mu}$ is often referred to as a white noise, and may roughly be regarded as a random signed measure on $\mathbb{H}_{d}$ (although strictly speaking it is not); the reason is that for any disjoint sets $V, V^{\prime} \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$, the two random variables $\mathfrak{B}_{\mu}\left(V \cup V^{\prime}\right)$ and $\mathfrak{B}_{\mu}(V)+\mathfrak{B}_{\mu}\left(V^{\prime}\right)$ coincide almost surely.

At this point, it is useful to mention that $\mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$ is a $\delta$-ring, i.e. is closed under symmetric difference and countable intersections. Moreover, there is a one-to-one correspondence between the set $\mathbb{H}_{d}$ and the collection of all $(d-1)$-dimensional hyperplanes in $\mathbb{R}^{d}$ that do not contain the origin. Indeed, any such hyperplane $h$ may be represented in a unique manner by a pair $(\rho, s) \in \mathbb{H}_{d}$ since it coincides with the set of all $t \in \mathbb{R}^{d}$ satisfying $\rho=\langle s, t\rangle$. In that correspondence, the hyperplanes that separate a given point $t \in \mathbb{R}^{d}$ and the origin are those which are represented by a pair that belongs to the set $V_{t} \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$ given by

$$
\begin{equation*}
V_{t}=\left\{(\rho, s) \in \mathbb{H}_{d} \mid \rho<\langle s, t\rangle\right\} . \tag{5}
\end{equation*}
$$

The geometric intuition behind the definition of the random field $B_{\mu}$ is that its value at a given point $t$ is determined by the mass that $\mathfrak{B}_{\mu}$ assigns to the hyperplanes separating $t$ and the origin. Specifically,

$$
\forall t \in \mathbb{R}^{d} \quad B_{\mu}(t)=\mathfrak{B}_{\mu}\left(V_{t}\right)
$$

More generally, the increment of the field $B_{\mu}$ between two points in $\mathbb{R}^{d}$ is determined by the mass of the hyperplanes that separate them.

### 3.2 Basic properties

It is quite straightforward to establish that $B_{\mu}$ falls in the category of Lévy fields. In particular, checking the stationary increments property for $B_{\mu}$ calls upon the fact that for any $t, t^{\prime} \in \mathbb{R}^{d}$, the random variable $B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)$ is normally distributed with
mean zero and variance

$$
\begin{equation*}
\mathfrak{d}_{\mu}\left(t, t^{\prime}\right)^{2}=\mathbb{E}\left[\left(B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)\right)^{2}\right]=\frac{1}{2} \int_{s \in \mathbb{S}^{d-1}}\left|\left\langle s, t^{\prime}-t\right\rangle\right| \mu(\mathrm{d} s), \tag{6}
\end{equation*}
$$

which depends on $t$ and $t^{\prime}$ only through $t^{\prime}-t$. We refer to [33, Eq. (5.11)] for details on how to derive the last equality in (6). In addition, it is customary to observe that $\mathfrak{d}_{\mu}$ defines a pseudometric on $\mathbb{R}^{d}$ which satisfies

$$
\begin{equation*}
\forall t, t^{\prime} \in \mathbb{R}^{d} \quad \mathfrak{d}_{\mu}\left(t, t^{\prime}\right) \leq c_{\mu}\left\|t-t^{\prime}\right\|^{1 / 2} \tag{7}
\end{equation*}
$$

with $c_{\mu}=\left(\mu\left(\mathbb{S}^{d-1}\right) / 2\right)^{1 / 2}$, in view of (6) and the Cauchy-Schwarz inequality. Here and below, $\|\cdot\|$ denotes the Euclidean norm. This observation implies that the field $B_{\mu}$ is stochastically continuous, and also admits a separable modification, see e.g. [31, Sect. 4]. For technical reasons that come into play in Sect. 5 (see the proof of Proposition 2), we need to work with such a modification. Therefore, we assume throughout that the field $B_{\mu}$ is separable.

## 4 The jump component

As seen in Sect. 2.1, the jump component of a Lévy field is a random field $L_{v}=$ $\left\{L_{v}(t), t \in \mathbb{R}^{d}\right\}$ that depends on a nonnegative Borel measure $v$ on $\mathbb{S}^{d-1} \times \mathbb{R}^{*}$. Recall that $\nu$ is symmetric and satisfies (1). Moreover, let $\mathcal{L}_{+}^{1}$ be the Lebesgue measure on $(0, \infty)$ and let N denote a Poisson random measure on $(0, \infty) \times \mathbb{S}^{d-1} \times \mathbb{R}^{*}$ with intensity $\mathcal{L}_{+}^{1} \otimes v$, see e.g. [34].

### 4.1 Multivariate compound Poisson processes

For any set $V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$, let

$$
\mathfrak{L}_{v, 0}(V)=\int_{\substack{(\rho, s) \in V \\|s|>1}} x \mathrm{~N}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x)
$$

In our situation, what plays the role of a compound Poisson process with jumps of magnitude larger than one is the random field $L_{v, 0}=\left\{L_{v, 0}(t), t \in \mathbb{R}^{d}\right\}$ defined by $L_{v, 0}(t)=\mathfrak{L}_{\nu, 0}\left(V_{t}\right)$ for any $t \in \mathbb{R}^{d}$, where $V_{t}$ is given by (5). In fact, letting ( $\mathrm{P}_{n}, S_{n}, X_{n}$ ), for $n \geq 1$, denote the atoms of the Poisson measure N , we have

$$
\begin{equation*}
L_{v, 0}(t)=\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\mathrm{P}_{n}<\left\langle S_{n}, t\right\rangle,\left|X_{n}\right|>1\right\}} \tag{8}
\end{equation*}
$$

for all $t \in \mathbb{R}^{d}$. For any integer $A \geq 1$, when $t$ ranges over the closed ball with radius $A$ centered at the origin, the sum in (8) may actually be restricted to the integers
$n \geq 1$ satisfying both $\mathrm{P}_{n}<A$ and $\left|X_{n}\right|>1$, which are almost surely finitely many. Therefore, with probability one, $L_{v, 0}$ is piecewise constant, with jumps of magnitude $\left|X_{n}\right|>1$ located on the hyperplanes $H_{n}$ parametrized by $\left(\mathrm{P}_{n}, S_{n}\right) \in \mathbb{H}_{d}$, which are given by

$$
\begin{equation*}
H_{n}=\left\{t \in \mathbb{R}^{d} \mid \mathrm{P}_{n}=\left\langle S_{n}, t\right\rangle\right\} \tag{9}
\end{equation*}
$$

### 4.2 Multivariate compensated sums of jumps

For any integer $j \geq 1$, the compensated sum corresponding to the jumps of magnitude in $\mathcal{I}_{j}=\left(2^{-j}, 2^{-j+1}\right]$ is the field $L_{v, j}=\left\{L_{v, j}(t), t \in \mathbb{R}^{d}\right\}$ given by $L_{v, j}(t)=$ $\mathfrak{L}_{v, j}\left(V_{t}\right)$ for any $t \in \mathbb{R}^{d}$, where

$$
\begin{equation*}
\mathfrak{L}_{v, j}(V)=\int_{\substack{(\rho, s) \in V \\|x| \in \mathcal{I}_{j}}} x \mathrm{~N}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x)-\int_{\substack{(\rho, s) \in V \\|x| \in \mathcal{I}_{j}}} x \mathrm{~d} \rho v(\mathrm{~d} s, \mathrm{~d} x) \tag{10}
\end{equation*}
$$

for any Borel set $V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$. Note that $\mathfrak{L}_{v, j}(V)$ may be regarded as an integral with respect to the compensated Poisson measure $\mathrm{N}^{*}=\mathrm{N}-\mathcal{L}_{+}^{1} \otimes v$ associated with N . In addition, due to the symmetry of $v$, we have

$$
\begin{equation*}
L_{v, j}(t)=\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\mathrm{P}_{n}<\left\langle S_{n}, t\right\rangle,\left|X_{n}\right| \in \mathcal{I}_{j}\right\}}-\int_{\substack{s \in \mathbb{S}_{d-1}^{d-1} \\ x \in \mathcal{I}_{j}}} x\langle s, t\rangle \nu(\mathrm{d} s, \mathrm{~d} x) \tag{11}
\end{equation*}
$$

for all $t \in \mathbb{R}^{d}$. The sum in (11) is almost surely piecewise constant with jumps of magnitude $\left|X_{n}\right| \in \mathcal{I}_{j}$ located on the hyperplanes $H_{n}$ given by (9), while the compensating integral depends linearly on $t$.

### 4.3 Definition and basic properties of $L_{v}$

The series formed by the compensated sums (10) for $j \geq 1$ converges and yields a Lévy field with jumps of magnitude at most one. As a matter of fact, for any $V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$ and $j \geq 1$, Campbell's theorem [27] ensures that the random variable $\mathfrak{L}_{v, j}(V)$ is centered with variance

$$
\int_{\substack{(\rho, s) \in V \\|x| \in \mathcal{I _ { j }}}} x^{2} \mathrm{~d} \rho v(\mathrm{~d} s, \mathrm{~d} x)
$$

In view of (1), it follows that the series $\sum_{j \geq 1} \mathfrak{L}_{v, j}(V)$ converges in $L^{2}$. Thus, it converges almost surely in view of the Lévy-Itō-Nisio theorem for sums of independent random variables, see [29, p. 151]. This enables us to define

$$
\mathfrak{L}_{v}(V)=\mathfrak{L}_{v, 0}(V)+\sum_{j=1}^{\infty} \mathfrak{L}_{v, j}(V)
$$

Then, $\mathfrak{L}_{v}=\left\{\mathfrak{L}_{v}(V), V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)\right\}$ is an infinitely divisible random measure on $\mathbb{H}_{d}$ with control measure $\nu$, in the sense that it satisfies the next properties:

- it is independently scattered, i.e. for any disjoint sets $V_{1}, \ldots, V_{n} \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$, the random variables $\mathfrak{L}_{v}\left(V_{1}\right), \ldots, \mathfrak{L}_{v}\left(V_{n}\right)$ are independent;
- it is $\sigma$-additive, i.e. for any sequence $\left(V_{n}\right)_{n \geq 1}$ of disjoint sets in $\mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$ whose union belongs to $\mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$, the series $\sum_{n \geq 1} \mathfrak{L}_{v}\left(V_{n}\right)$ converges almost surely and its sum is equal to $\mathfrak{L}_{v}\left(\bigcup_{n \geq 1} V_{n}\right)$;
- for every $V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$, the characteristic function of $\mathfrak{L}_{v}(V)$ is given by

$$
\begin{equation*}
\mathbb{E}\left[\mathrm{e}^{i \theta \mathfrak{L}_{v}(V)}\right]=\exp \int_{\substack{(\rho, s) \in V \\ x \in \mathbb{R}^{*}}}\left(\mathrm{e}^{i \theta x}-1-i \theta x \mathbb{1}_{\{|x| \leq 1\}}\right) \mathrm{d} \rho \nu(\mathrm{~d} s, \mathrm{~d} x) \tag{12}
\end{equation*}
$$

The first two properties directly follow from standard results on Poisson random measures, while the third one is a consequence of Campbell's theorem, see [27,34].

Making use of these properties, it is straightforward to check that the random field $L_{v}=\left\{L_{v}(t), t \in \mathbb{R}^{d}\right\}$ defined by

$$
\forall t \in \mathbb{R}^{d} \quad L_{v}(t)=\mathfrak{L}_{v}\left(V_{t}\right)
$$

is a Lévy field, see [33] for details. Replacing the measure $\nu$ by appropriate restrictions, this implies that the fields $L_{v, j}$ defined above are of Lévy type as well.

### 4.4 Comments

Mori opted for $1 /\left(1+x^{2}\right)$ instead of the cut-off function $\mathbb{1}_{\{|x| \leq 1\}}$ in (12). This leads to a slightly different expression of the field $\mathfrak{L}_{v}$ and of the drift coefficient $a$. Yet, our choice clearly does not compromise the validity of Theorem 1 and does alter the value of neither $\mu$ nor $v$, given a Lévy field $Y$. Furthermore, Mori did not detail the construction of the infinitely divisible random measure $\mathfrak{L}_{v}$ on which the field $L_{v}$ is based. However, the proofs below call upon a precise knowledge of the jump structure of $L_{v}$, and this explains why we chose above to present the construction of $\mathfrak{L}_{v}$ in terms of compensated Poisson integrals.

Furthermore, in order to study the regularity of $L_{v}$, we first need to make sure that its Hölder exponent is a well-defined quantity. This boils down to verifying that the sample functions of that field exist almost surely. In fact, as yet, $L_{v}$ has been defined in a pointwise manner only: we merely proved the almost sure convergence of the series $\sum_{j \geq 0} L_{\nu, j}(t)$ defining $L_{v}(t)$, for every fixed $t \in \mathbb{R}^{d}$. This enables us to consider the finite-dimensional marginals of the field $L_{v}$, which is sufficient to state Theorem 1. However, in order to determine the value of the Hölder exponent $\alpha_{L_{v}}(t)$ at a given point $t \in \mathbb{R}^{d}$, we need to consider $L_{v}$ everywhere near $t$. The next result indicates that this is possible.

Proposition 1 Let v be an admissible measure. Then, with probability one,

$$
\forall t \in \mathbb{R}^{d} \quad L_{v}(t)=\sum_{j=0}^{\infty} L_{v, j}(t) \text { exists. }
$$

The proof of Proposition 1 is postponed to Sect. 12.1. It relies on precise estimates of the increments of the fields $L_{v, j}$ that are given by Lemma 2 below. It can also easily be adapted to show that, when the measure $v$ is admissible, the sample paths of the field $L_{v}$ are almost surely locally bounded; this may be seen as an extension of [39, Theorem 2.3], which concerns the Lévy-Chentsov fields only.

In what follows, whenever $v$ is supposed to be admissible, we implicitly work on the almost sure event where the sample paths of the field $L_{v}$ exist.

## 5 Regularity of the Gaussian component

This section is devoted to the proof of Theorem 2; we shall in fact establish two slightly more precise results. First, the Hölder exponent of the random field $B_{\mu}$ is almost surely at least $1 / 2$ everywhere, as a direct consequence of the following result on its modulus of continuity. Recall that $c_{\mu}=\left(\mu\left(\mathbb{S}^{d-1}\right) / 2\right)^{1 / 2}$.

Proposition 2 There exists a universal constant $\mathrm{K}>0$ such that for any finite symmetric measure $\mu$ on $\mathbb{S}^{d-1}$ satisfying $c_{\mu}>0$ and for any integer $A \geq 1$,

$$
\begin{equation*}
\text { a.s. } \quad \limsup _{\delta \rightarrow 0} \frac{1}{(\delta \log (1 / \delta))^{1 / 2}} \sup _{\substack{t, t^{\prime} \in[-A, A]^{d} \\\left\|t-t^{\prime}\right\| \leq \delta}}\left|B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)\right| \leq \mathrm{K} c_{\mu} d^{1 / 2} \tag{13}
\end{equation*}
$$

Proof As mentioned in Sect. 3, the field $B_{\mu}$ is assumed to be separable. We may therefore apply [2, Theorem 1.3.5]. Accordingly, there exists a universal constant $\mathrm{K}>0$ such that with probability one, for $\eta>0$ small enough,

$$
\begin{equation*}
\sup _{\substack{t, t^{\prime} \in[-A, A]^{d} \\ \mathfrak{o} \mu\left(t, t^{\prime}\right) \leq \eta}}\left|B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)\right| \leq \mathrm{K} \int_{0}^{\eta}\left(\log N\left([-A, A]^{d}, \mathfrak{o}_{\mu}, \varepsilon\right)\right)^{1 / 2} \mathrm{~d} \varepsilon . \tag{14}
\end{equation*}
$$

Here, $N\left([-A, A]^{d}, \mathfrak{d}_{\mu}, \varepsilon\right)$ denotes the minimal number of balls with radius $\varepsilon$ that cover the cube $[-A, A]^{d}$, the balls being closed, centered in that cube and taken in the sense of the pseudometric $\mathfrak{d}_{\mu}$ defined by (6). Letting $\lfloor\cdot\rfloor$ stand for the floor function, it is easy to check that the aforementioned cube is covered by $\left(1+2\left\lfloor A\left(c_{\mu} / \varepsilon\right)^{2} d^{1 / 2}\right\rfloor\right)^{d}$ closed Euclidean balls centered in it with radius $\left(\varepsilon / c_{\mu}\right)^{2}$, and (7) implies that each of these balls is included in a closed ball with radius $\varepsilon$ for the pseudometric $\mathfrak{d}_{\mu}$. Hence,
the right-hand side of (14) is bounded above by

$$
\mathrm{K} \int_{0}^{\eta}\left(d \log \left(1+\frac{2 A c_{\mu}^{2} d^{1 / 2}}{\varepsilon^{2}}\right)\right)^{1 / 2} \mathrm{~d} \varepsilon \sim \mathrm{~K} \eta\left(2 d \log \frac{c_{\mu}}{\eta}\right)^{1 / 2} \quad \text { as } \eta \rightarrow 0
$$

To conclude, it now suffices to let $\eta=c_{\mu} \sqrt{\delta}$ and observe that the supremum in left-hand side of (13) is bounded above by the left-hand side of (14), due to (7).

Second, the fact that the Hölder exponent of $B_{\mu}$ is almost surely at most $1 / 2$ everywhere follows directly from the next proposition; $\overline{\mathrm{B}}_{t}(\delta)$ will denote the closed Euclidean ball centered at $t$ with radius $\delta$.

Proposition 3 For any finite symmetric measure $\mu$ on $\mathbb{S}^{d-1}$ satisfying $c_{\mu}>0$, there exists $\kappa_{d, \mu}>0$ such that

$$
\text { a.s. } \forall t \in \mathbb{R}^{d} \quad \forall \delta>0 \quad \exists t^{\prime} \in \overline{\mathrm{B}}_{t}(\delta) \quad\left|B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)\right|>\kappa_{d, \mu}\left\|t^{\prime}-t\right\|^{1 / 2}
$$

Proof We shall adapt some ideas that Dvoretzky [20] employed in the case of Brownian motion. To begin with, let $\left(e_{1}, \ldots, e_{d}\right)$ denote the canonical basis of $\mathbb{R}^{d}$. We necessarily have $\mathfrak{d}_{\mu}\left(0, e_{i}\right)>0$ for some $i \in\{1, \ldots, d\}$, because

$$
\sum_{i=1}^{d} \mathfrak{d}_{\mu}\left(0, e_{i}\right)^{2} \geq \frac{1}{2} \sum_{i=1}^{d} \int_{s \in \mathbb{S}^{d-1}}\left|\left\langle s, e_{i}\right\rangle\right|^{2} \mu(\mathrm{~d} s)=\frac{1}{2} \int_{s \in \mathbb{S}^{d-1}}\|s\|^{2} \mu(\mathrm{~d} s)=c_{\mu}^{2}>0
$$

due to (6). Then, let us consider an integer $A \geq 1$, a real number $\kappa>0$ and let us assume that for any $\delta>0$, there exists a point $t \in[-A, A]^{d}$ such that

$$
\begin{equation*}
\forall t^{\prime} \in \overline{\mathrm{B}}_{t}(\delta) \quad\left|B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)\right| \leq \kappa\left\|t^{\prime}-t\right\|^{1 / 2} \tag{15}
\end{equation*}
$$

Hence, for every integer $n \geq 1$, there exists $k \in\{-A n, \ldots, A n-1\}^{d}$ such that (15) holds for some $t \in \mathbb{R}^{d}$ with $t-k / n \in[0,1 / n]^{d}$. Let $\log _{2}$ stand for base two logarithm and $j_{n}=\left\lfloor\log _{2}\left(n \delta / d^{1 / 2}\right)\right\rfloor-1$, and assume that $n$ is large enough to ensure that $j_{n} \geq 1$. For any $j \in\left\{0, \ldots, j_{n}\right\}$, the point $\left(k+2^{j} e_{i}\right) / n$ is in the ball $\overline{\mathrm{B}}_{t}(\delta)$, so that

$$
\left|B_{\mu}\left(\frac{k+2^{j} e_{i}}{n}\right)-B_{\mu}(t)\right| \leq \kappa\left\|\frac{k+2^{j} e_{i}}{n}-t\right\|^{1 / 2} \leq \kappa\left(\frac{2^{j+1} d^{1 / 2}}{n}\right)^{1 / 2}
$$

thanks to (15). Due to the triangle inequality, this implies that

$$
\begin{equation*}
\forall j \in\left\{1, \ldots, j_{n}\right\} \quad\left|B_{\mu}\left(\frac{k+2^{j} e_{i}}{n}\right)-B_{\mu}\left(\frac{k+2^{j-1} e_{i}}{n}\right)\right| \leq \kappa \frac{d^{1 / 4}}{\sqrt{n}} 2^{(j+3) / 2} \tag{16}
\end{equation*}
$$

The increments in the left-hand side, for $j \in\left\{1, \ldots, j_{n}\right\}$, are independent and normally distributed with mean zero and variance $2^{j-1} \mathfrak{d}_{\mu}\left(0, e_{i}\right)^{2} / n$ by virtue of (6). Thus, (16) occurs with probability at most

$$
\prod_{j=1}^{j_{n}} \mathbb{P}\left(\left|\frac{2^{j / 2}}{\sqrt{2 n}} \mathfrak{d}_{\mu}\left(0, e_{i}\right) \zeta\right| \leq \kappa \frac{d^{1 / 4}}{\sqrt{n}} 2^{(j+3) / 2}\right)=q(\kappa)^{j_{n}},
$$

where $q(\kappa)=\mathbb{P}\left(|\zeta| \leq 4 \kappa d^{1 / 4} / \mathfrak{o}_{\mu}\left(0, e_{i}\right)\right)$ and $\zeta$ denotes a standard normal random variable. As a result, the probability that (15) holds for some point $t \in[-A, A]^{d}$ is at most $(2 A n)^{d} q(\kappa)^{j_{n}}$ for every integer $n$ such that $j_{n} \geq 1$. Clearly, $q\left(\kappa_{d, \mu}\right)<2^{-d}$ for some $\kappa_{d, \mu}>0$, which implies that $(2 A n)^{d} q\left(\kappa_{d, \mu}\right)^{j_{n}} \rightarrow 0$ as $n \rightarrow \infty$. We deduce that for any integer $A \geq 1$ and any $\delta>0$,

$$
\text { a.s. } \forall t \in[-A, A]^{d} \quad \exists t^{\prime} \in \overline{\mathrm{B}}_{t}(\delta) \quad\left|B_{\mu}\left(t^{\prime}\right)-B_{\mu}(t)\right|>\kappa_{d, \mu}\left\|t^{\prime}-t\right\|^{1 / 2}
$$

The desired result clearly follows.

## 6 Regularity of the jump component

Let us detail our results on the regularity of the random field $L_{v}$ defined in Sect. 4. The corresponding proofs are given in Sect. 12, and their architecture is presented in Sect. 8. Throughout, we assume that $v$ is admissible.

In Sect. 6.2, we describe the size properties of the iso-Hölder sets $E_{L_{v}}(h)$ in a precise manner: we give the value of their Hausdorff $g$-measure in every open subset of $\mathbb{R}^{d}$, for every gauge function $g$. In what follows, the iso-Hölder sets are denoted by $E_{v}(h)$ instead of $E_{L_{v}}(h)$, for the sake of brevity. Specifically,

$$
\forall h \in[0, \infty] \quad E_{v}(h)=\left\{t \in \mathbb{R}^{d} \mid \alpha_{L_{v}}(t)=h\right\}
$$

We get similar results for the singularity sets raised in Sect. 1 and defined by

$$
\begin{equation*}
E_{v}^{\prime}(h)=\left\{t \in \mathbb{R}^{d} \backslash J_{v} \mid \alpha_{L_{v}}(t) \leq h\right\} . \tag{17}
\end{equation*}
$$

We also establish that the sets $E_{v}^{\prime}(h)$ fall in the category of sets with large intersection, see Sect. 6.3. As explained in Sect. 6.4, this property has straightforward consequences on the value of their packing dimension.

The set $J_{v}$ in (17) consists of the points at which $L_{v}$ jumps. Specifically, $J_{v}$ is the union over $n \geq 1$ of the hyperplanes $H_{n}$ defined by (9), in terms of the atoms of the Poisson measure N that arises in the construction of $L_{v}$. Equivalently,

$$
\begin{equation*}
J_{v}=\left\{t \in \mathbb{R}^{d} \mid \int_{\substack{(\rho, s) \in \mathbb{H}_{d} d \\ x \in \mathbb{R}^{*}}} \mathbb{1}_{\{\rho=\langle s, t\rangle\}} \mathrm{N}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x) \geq 1\right\} . \tag{18}
\end{equation*}
$$

### 6.1 Preliminary remark

We adopt the convention $1 / \beta_{v}=\infty$ if the index $\beta_{v}$ defined by (2) vanishes. The size and large intersection properties of $E_{\nu}(h)$ and $E_{\nu}^{\prime}(h)$ are nontrivial only when $h \in\left[0,1 / \beta_{v}\right)$. In fact, when $h \geq 1 / \beta_{v}$, this description follows essentially from the next result, which is proven in Sect. 12.5.

Proposition 4 Let v be an admissible measure. Then,

$$
\text { a.s. } \forall h \in\left[1 / \beta_{v}, \infty\right] \quad E_{v}^{\prime}(h)=\mathbb{R}^{d} \backslash J_{v} .
$$

Moreover, $E_{v}\left(1 / \beta_{v}\right)$ has full Lebesgue measure in $\mathbb{R}^{d}$ with probability one, and

$$
\text { a.s. } \forall h \in\left(1 / \beta_{v}, \infty\right] \quad E_{v}(h)=\emptyset .
$$

Given that $J_{\nu}$ has Lebesgue measure zero, Proposition 4 ensures that $E_{\nu}\left(1 / \beta_{v}\right)$ and $E_{\nu}^{\prime}\left(1 / \beta_{\nu}\right)$ have full Lebesgue measure in $\mathbb{R}^{d}$ almost surely. This result has direct implications in terms of Hausdorff measures and large intersection properties, which we shall detail below (see the comments following Theorems 4 and 6).
6.2 Size properties of the sets $E_{v}(h)$ and $E_{v}^{\prime}(h)$ : Hausdorff measures and dimension

We call a gauge function any nondecreasing function $g$ defined on $[0, \infty)$ such that $\lim _{0+} g=g(0)=0$ and $r \mapsto g(r) / r^{d}$ is positive and nonincreasing near zero (this last assumption is not particularly restrictive and may be removed using [16, Proposition 2]); $\mathfrak{D}_{d}$ will denote the set of gauge functions. For any $g \in \mathfrak{D}_{d}$, the Hausdorff $g$-measure of a subset $F$ of $\mathbb{R}^{d}$ is defined by

$$
\mathcal{H}^{g}(F)=\lim _{\delta \downarrow 0} \uparrow \mathcal{H}_{\delta}^{g}(F) \quad \text { with } \quad \mathcal{H}_{\delta}^{g}(F)=\inf _{\substack{F \\\left|\cup_{n}\right|<\delta}} \sum_{n=1}^{\infty} g\left(\left|U_{n}\right|\right) .
$$

Here, the infimum is taken over all sequences $\left(U_{n}\right)_{n \geq 1}$ of subsets of $\mathbb{R}^{d}$ satisfying $F \subseteq \bigcup_{n} U_{n}$ and $\left|U_{n}\right|<\delta$ for all $n \geq 1$, where $|\cdot|$ denotes diameter. As shown e.g. in [36], $\mathcal{H}^{g}$ is a Borel measure on $\mathbb{R}^{d}$. For simplicity, the Hausdorff measure corresponding to a gauge function of the form $r \mapsto r^{d-1} g(r)$ with $g \in \mathfrak{D}_{1}$ (obtained by letting $d=1$ in the definition of $\mathfrak{D}_{d}$ ) is denoted by $\mathcal{H}^{d-1, g}$.

When $g$ is of the form $r \mapsto r^{s}$ for some $s \in(0, d]$, the Hausdorff $g$-measure is referred to as the $s$-dimensional Hausdorff measure and denoted by $\mathcal{H}^{s}$. These particular measures lead to the notion of Hausdorff dimension. Specifically, the Hausdorff dimension of a nonempty set $F \subseteq \mathbb{R}^{d}$ is defined by

$$
\operatorname{dim}_{\mathrm{H}} F=\sup \left\{s \in(0, d) \mid \mathcal{H}^{s}(F)=\infty\right\}=\inf \left\{s \in(0, d) \mid \mathcal{H}^{s}(F)=0\right\}
$$

see [22]. Here, we adopt the convention that $\sup \emptyset=0$ and $\inf \emptyset=d$, and also that $\operatorname{dim}_{H} \emptyset=-\infty$.

### 6.2.1 General results

When $h \in\left[0,1 / \beta_{\nu}\right)$, the size properties of $E_{\nu}(h)$ and $E_{\nu}^{\prime}(h)$ are described by the next result. In its statement, $h_{\nu}(g)$ is given by

$$
\begin{equation*}
h_{v}(g)=\inf \left\{h>0 \mid \int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in(0,1]}} g\left(x^{1 / h}\right) v(\mathrm{~d} s, \mathrm{~d} x)=\infty\right\}, \tag{19}
\end{equation*}
$$

for $g \in \mathfrak{D}_{1}$, with the usual convention that $\inf \emptyset=\infty$. Clearly, $h_{\nu}(g) \leq 1 / \beta_{\nu}$.
Theorem 4 Let $v$ be an admissible measure. Then, with probability one, for any $h \in\left[0,1 / \beta_{v}\right)$, any gauge function $g \in \mathfrak{D}_{1}$ and any nonempty open $W \subseteq \mathbb{R}^{d}$,

$$
\mathcal{H}^{d-1, g}\left(E_{\nu}(h) \cap W\right)= \begin{cases}0 & \text { if } h<h_{\nu}(g) \\ \infty & \text { if } h=h_{\nu}(g)\end{cases}
$$

and

$$
\mathcal{H}^{d-1, g}\left(E_{v}^{\prime}(h) \cap W\right)= \begin{cases}0 & \text { if } h<h_{v}(g) \\ \infty & \text { if } h \geq h_{v}(g) .\end{cases}
$$

We refer to Sect. 12.2 for a proof of Theorem 4. In this result, it is remarkable that the almost sure event on which the statement holds does not depend on $h$, the gauge function $g$ or the open set $W$. In other words, the previous description of the size properties of the sets $E_{v}(h)$ and $E_{v}^{\prime}(h)$ holds for almost every sample function of the random field $L_{v}$. On top of that, the description is both precise and local, in the sense that we do not restrict our attention to the mere gauge functions of the form $r \mapsto r^{s}$ that lead to the Hausdorff dimension, or to the case in which the open set $W$ is equal to the whole space $\mathbb{R}^{d}$.

With the help of Proposition 4 above, it is easy to obtain an analog of Theorem 4 for the case where $h \geq 1 / \beta_{\nu}$. Specifically, with probability one, the sets $E_{\nu}\left(1 / \beta_{\nu}\right)$ and $E_{\nu}^{\prime}\left(1 / \beta_{\nu}\right)$ both have full Lebesgue measure in $\mathbb{R}^{d}$, so that

$$
\left.\forall g \in \mathfrak{D}_{d} \quad \forall W \text { open } \begin{array}{l}
\mathcal{H}^{g}\left(E_{\nu}\left(1 / \beta_{\nu}\right) \cap W\right) \\
\mathcal{H}^{g}\left(E_{\nu}^{\prime}\left(1 / \beta_{\nu}\right) \cap W\right)
\end{array}\right\}=\mathcal{H}^{g}(W) .
$$

The same result holds for the sets $E_{v}^{\prime}(h)$, with $h \in\left(1 / \beta_{v}, \infty\right]$, because they all coincide with $E_{v}^{\prime}\left(1 / \beta_{v}\right)$. For those values of $h$, the study of the size of the sets $E_{v}(h)$ is pointless because they are empty almost surely.

In the remainder of this section, we make additional assumptions on the positivity of $\beta_{v}$ or the finiteness of $v$ with a view to obtaining further results. Due to (2), it is clear that $\beta_{v}=0$ when $v$ has finite mass. This corresponds to the case of a compound Poisson field.

### 6.2.2 Case where $\beta_{v}>0$

Here, Theorem 4 leads to the following more compact result, which is established in Sect. 12.3. In its statement, $\mathfrak{D}_{1}^{*}$ is the set of gauge functions $g \in \mathfrak{D}_{1}$ for which $\lim _{r \rightarrow 0}(\log g(r)) / \log r$ exists.

Corollary 2 Let $v$ be an admissible measure satisfying $\beta_{v}>0$. Then, with probability one, for any $h \in\left[0,1 / \beta_{v}\right)$, any gauge function $g \in \mathfrak{D}_{1}^{*}$ and any nonempty open subset $W$ of $\mathbb{R}^{d}$,

$$
\mathcal{H}^{d-1, g}\left(E_{\nu}(h) \cap W\right)=\mathcal{H}^{d-1, g}\left(E_{v}^{\prime}(h) \cap W\right)= \begin{cases}0 & \text { if } h<h_{v}(g) \\ \infty & \text { if } h \geq h_{v}(g),\end{cases}
$$

so that

$$
\operatorname{dim}_{\mathrm{H}}\left(E_{v}(h) \cap W\right)=\operatorname{dim}_{\mathrm{H}}\left(E_{v}^{\prime}(h) \cap W\right)=d-1+\beta_{v} h
$$

In particular, the sets $E_{v}(h) \cap W$ have Hausdorff dimension $d-1+\beta_{v} h$, with an infinite $\left(d-1+\beta_{\nu} h\right)$-dimensional Hausdorff measure.

### 6.2.3 Case where $\beta_{v}=0$ and $v$ has infinite total mass

In that situation, we have $h_{v}\left(r \mapsto r^{s}\right)=\infty$ for all $s \in(0,1]$ and, borrowing ideas from the proof of [18, Proposition 5], we may build a gauge function $g \in \mathfrak{D}_{1}$ such that $h_{\nu}(g)=0$. Hence, Theorem 4 leads to the following statement.

Corollary 3 Let $v$ be a measure with $\beta_{v}=0$ and infinite total mass. Then, with probability one, for any $h \in[0, \infty)$ and any nonempty open subset $W$ of $\mathbb{R}^{d}$,

$$
\operatorname{dim}_{H}\left(E_{\nu}^{\prime}(h) \cap W\right)=d-1
$$

and the $(d-1)$-dimensional Hausdorff measure of $E_{v}^{\prime}(h) \cap W$ is infinite.
Let us now focus on the iso-Hölder sets $E_{v}(h)$. Theorem 4, coupled with the preceding remarks, also implies that with probability one, for any $h \in[0, \infty)$ and any nonempty open subset $W$ of $\mathbb{R}^{d}$,

$$
\begin{equation*}
\operatorname{dim}_{\mathrm{H}}\left(E_{\nu}(h) \cap W\right) \leq d-1 \tag{20}
\end{equation*}
$$

However, we cannot conclude that (20) is an equality, except when $h=0$. According to our approach, inferring that equality holds in (20) requires that we build, for any given $h \in[0, \infty)$, a gauge function $g \in \mathfrak{D}_{1}$ with $h_{v}(g)=h$, and apply Theorem 4 with this gauge function. As mentioned above, such a construction is feasible when $h=0$, but is not always possible otherwise. Indeed, some measures $v$ yield $h_{v}(g) \in\{0, \infty\}$ for all $g \in \mathfrak{D}_{1}$. This is the case of the product of the uniform measure on $\mathbb{S}^{d-1}$ and the sum over $j \in \mathbb{Z} \backslash\{0\}$ of the atoms $\delta_{1 / j} /|j|$ on $\mathbb{R}^{*}$.

### 6.2.4 Case where v has finite total mass

Making use of (8) and (11), we get

$$
\begin{equation*}
\text { a.s. } \forall t \in \mathbb{R}^{d} \quad L_{\nu}(t)=\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\mathrm{P}_{n}<\left\langle S_{n}, t\right\rangle\right\}}-\int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in(0,1]}} x\langle s, t\rangle \nu(\mathrm{d} s, \mathrm{~d} x) . \tag{21}
\end{equation*}
$$

With probability one, the above sum is piecewise constant, with jumps located on the set $J_{v}$ given by (18). Moreover, the integral in (21) depends linearly on $t$. Thus, it is natural that the previous results lead to the next statement.

Proposition 5 Let v be a measure with finite total mass. Then, with probability one, for any $h \in[0, \infty)$,

$$
E_{v}^{\prime}(h)=\emptyset \text { and } E_{\nu}(h)= \begin{cases}J_{v} & \text { if } h=0 \\ \emptyset & \text { if } h>0\end{cases}
$$

Although Proposition 5 is elementary, a formal proof is given in Sect. 12.4 for the sake of completeness. By virtue of Proposition 4, we also have $E_{v}(\infty)=E_{v}^{\prime}(\infty)=$ $\mathbb{R}^{d} \backslash J_{v}$ almost surely. Hence,

$$
\text { a.s. } \quad \forall h \in[0, \infty] \quad \operatorname{dim}_{\mathrm{H}} E_{v}(h)= \begin{cases}d-1 & \text { if } h=0 \\ -\infty & \text { if } 0<h<\infty \\ d & \text { if } h=\infty\end{cases}
$$

In the present case, the sample paths of $L_{v}$ are not homogeneous. Indeed, with probability one, for any $A \geq 1$, the ball $\overline{\mathrm{B}}_{0}(A)$ intersects only finitely many hyperplanes $H_{n}$. Thus, with probability one, there exist nonempty open sets $W$ such that $J_{v} \cap W=\emptyset$, which implies that $d_{L_{v}}(0, W)=-\infty \neq d-1=d_{L_{v}}\left(0, \mathbb{R}^{d}\right)$.

### 6.3 Large intersection properties of the sets $E_{v}^{\prime}(h)$

As shown below, the sets $E_{v}^{\prime}(h)$ defined by (17) are sets with large intersection. Specifically, they belong to certain classes $\mathrm{G}^{g}(W)$ of subsets of $\mathbb{R}^{d}$, which were introduced in [16] in order to generalize the original classes of sets with large intersection defined by Falconer [21]. Given a gauge function $g \in \mathfrak{D}_{d}$ and a nonempty open set $W \subseteq \mathbb{R}^{d}$, the class $\mathrm{G}^{g}(W)$ may be defined in the following manner. To begin with, recall that the function $r \mapsto g(r) / r^{d}$ is positive and nonincreasing in a neighborhood of zero. Let $\varepsilon_{g}$ denote the supremum of all $\varepsilon \in(0,1]$ such that this property holds on the interval $(0, \varepsilon]$. Moreover, let $\Lambda_{g}$ be the set of all dyadic cubes of diameter less than $\varepsilon_{g}$, that is, sets of the form $\lambda=2^{-j}\left(k+[0,1)^{d}\right)$, where $j$ is an integer larger than $\log _{2}\left(\sqrt{d} / \varepsilon_{g}\right)$ and $k \in \mathbb{Z}^{d}$. The outer net measure associated with $g \in \mathfrak{D}_{d}$ is defined by

$$
\forall F \subseteq \mathbb{R}^{d} \quad \mathcal{M}_{\infty}^{g}(F)=\inf _{\left(\lambda_{n}\right)_{n \geq 1}} \sum_{n=1}^{\infty} g\left(\left|\lambda_{n}\right|\right)
$$

where the infimum is taken over all sequences $\left(\lambda_{n}\right)_{n \geq 1}$ in $\Lambda_{g} \cup\{\emptyset\}$ satisfying $F \subseteq$ $\bigcup_{n} \lambda_{n}$. The outer measure $\mathcal{M}_{\infty}^{g}$ is intimately related with the corresponding Hausdorff measure $\mathcal{H}^{g}$, so that in particular $\mathcal{H}^{g}(F)>0$ for any set $F \subseteq \mathbb{R}^{d}$ with $\mathcal{M}_{\infty}^{g}(F)>0$, see [36, Theorem 49]. In addition, for $\bar{g}, g \in \mathfrak{D}_{d}$, let us write $\bar{g} \prec g$ if $\bar{g} / g$ tends monotonically to infinity at zero. This means essentially that $\bar{g}$ grows faster than $g$ near the origin. We can now define the class $\mathrm{G}^{g}(W)$. Recall that a $G_{\delta}$-set is one that may be expressed as a countable intersection of open sets.

Definition 5 Let $g \in \mathfrak{D}_{d}$ and let $W$ be a nonempty open subset of $\mathbb{R}^{d}$. The class $\mathrm{G}^{g}(W)$ of subsets of $\mathbb{R}^{d}$ with large intersection in $W$ with respect to $g$ is the collection of all $G_{\delta}$-subsets $F$ of $\mathbb{R}^{d}$ such that

$$
\forall U \subseteq W \text { open } \quad \forall \bar{g} \in \mathfrak{D}_{d} \quad \bar{g} \prec g \Longrightarrow \mathcal{M}_{\infty}^{\bar{g}}(F \cap U)=\mathcal{M}_{\infty}^{\bar{g}}(U)
$$

The class $\mathrm{G}^{g}(W)$ satisfies several remarkable properties which are detailed in [16]. We collect the most significant ones in the following statement.

Theorem 5 Let $g \in \mathfrak{D}_{d}$ and let $W$ be a nonempty open subset of $\mathbb{R}^{d}$. Then,

1. the class $\mathrm{G}^{g}(W)$ is closed under countable intersections;
2. the set $f^{-1}(F)$ belongs to $\mathrm{G}^{g}(W)$ for any bi-Lipschitz mapping $f: W \rightarrow \mathbb{R}^{d}$ and any set $F \in \mathrm{G}^{g}(f(W))$;
3. any set $F \in \mathrm{G}^{g}(W)$ satisfies $\mathcal{H}^{\bar{g}}(F \cap W)=\infty$ for any $\bar{g} \in \mathfrak{D}_{d}$ with $\bar{g} \prec g$;
4. any $G_{\delta}$-subset of $\mathbb{R}^{d}$ with full Lebesgue measure in $W$ belongs to $\mathrm{G}^{g}(W)$.

Making use of Theorem 5(3), it is easy to check that any set that belongs to the class $\mathrm{G}^{g}(W)$ is of Hausdorff dimension at least

$$
\begin{equation*}
\sigma_{g}=\sup \left\{s \in(0, d) \mid\left(r \mapsto r^{s}\right) \prec g\right\} \tag{22}
\end{equation*}
$$

(with the convention that $\sup \emptyset=0$ ). In other words, the fact that a set satisfies a large intersection property leads to a lower bound on its Hausdorff dimension, an information which is usually difficult to derive.

More generally, Theorem 5 (3) embodies the connection between size and large intersection properties, thereby suggesting a link between the following result and Theorem 4. We shall make this link more apparent when proving these two theorems concurrently in Sect. 12.2. In the next statement, $\mathrm{G}^{d-1, g}(W)$ denotes the class of sets with large intersection in $W$ with respect to $r \mapsto r^{d-1} g(r)$, where $g \in \mathfrak{D}_{1}$.

Theorem 6 Let v be an admissible measure. Then, with probability one, for any $h \in\left[0,1 / \beta_{\nu}\right)$, any gauge function $g \in \mathfrak{D}_{1}$ and any nonempty open subset $W$ of $\mathbb{R}^{d}$,

$$
E_{\nu}^{\prime}(h) \in \mathrm{G}^{d-1, g}(W) \Longleftrightarrow h \geq h_{\nu}(g)
$$

The previous result only concerns the case where $h<1 / \beta_{\nu}$. As regards the opposite case, recall that $E_{\nu}^{\prime}\left(1 / \beta_{\nu}\right)$ has full Lebesgue measure in $\mathbb{R}^{d}$ with probability one, by Proposition 4. Hence, owing to Theorem 5(4),

$$
\text { a.s. } \forall g \in \mathfrak{D}_{d} \quad \forall W \neq \emptyset \text { open } \quad E_{v}^{\prime}\left(1 / \beta_{v}\right) \in \mathrm{G}^{g}(W) .
$$

The sets $E_{v}^{\prime}(h)$, for $h \in\left(1 / \beta_{v}, \infty\right]$, satisfy the same property because they are all identical to $E_{v}^{\prime}\left(1 / \beta_{\nu}\right)$.

### 6.4 Size properties of the sets $E_{v}^{\prime}(h)$ : packing dimension

We do not recall the definition of the packing dimension $\operatorname{dim}_{\mathrm{P}} F$ of a subset $F$ of $\mathbb{R}^{d}$ here, and refer to [22, Chap. 3] or [32, Chap. 5] for a full exposition. The fact that the sets $E_{v}^{\prime}(h)$ satisfy a large intersection property has a straightforward consequence on the value of their packing dimension, as we now explain.

First, in view of Proposition 4, it is clear that with probability one, the sets $E_{\nu}^{\prime}(h)$, for $h \geq 1 / \beta_{\nu}$, all have packing dimension equal to $d$. Therefore, in what follows, we may restrict our attention to the case where $h<1 / \beta_{v}$.

Assume that $v$ has infinite total mass and that $d>1$. As mentioned above, there exists a gauge function $g \in \mathfrak{D}_{1}$ such that $h_{\nu}(g)=0$. Thanks to Theorem 6, almost surely, for every $h \in\left[0,1 / \beta_{\nu}\right)$, the set $E_{\nu}^{\prime}(h)$ belongs to the class $G^{\widetilde{g}}\left(\mathbb{R}^{d}\right)$, where $\widetilde{g}: r \mapsto r^{d-1} g(r)$. It is shown in [16] that if the parameter $\sigma_{\widetilde{g}}$ defined by (22) is positive, then this class is included in the class $\mathcal{G}^{\sigma_{\sigma}}$ introduced by Falconer [21]. This is clearly the case here, since $\sigma_{\tilde{g}} \geq d-1>0$. By virtue of [21, Theorem D], every set of the latter class has packing dimension $d$ in every nonempty open subset of $\mathbb{R}^{d}$. Hence, we end up with the next statement.

Corollary 4 If $d>1$, then for any admissible measure $v$ with infinite total mass,

$$
\text { a.s. } \forall h \in\left[0,1 / \beta_{v}\right) \quad \forall W \neq \emptyset \text { open } \quad \operatorname{dim}_{\mathrm{P}}\left(E_{v}^{\prime}(h) \cap W\right)=d .
$$

The preceding result remains valid for $d=1$ under the additional assumption that $\beta_{\nu}>0$, but for $h \in\left(0,1 / \beta_{\nu}\right)$ only. This follows from applying the previous method with the gauge function $\widetilde{g}: r \mapsto r^{\beta_{v} h}$, which satisfies $\sigma_{\tilde{g}}=\beta_{v} h>0$.

## 7 Superposition of the Gaussian and the jump components

The next lemma gives an expression of the Hölder exponent of an admissible canonical Lévy field $Y_{a, \mu, \nu}=\langle a, \cdot\rangle+B_{\mu}+L_{\nu}$, in terms of the Hölder exponents of its Gaussian component $B_{\mu}$ and its jump component $L_{\nu}$. In particular, it allows one to easily infer the spectrum of singularities of $Y_{a, \mu, \nu}$ from those of $B_{\mu}$ and $L_{\nu}$, that is, to deduce Corollary 1 from Theorems 2 and 3.

Lemma 1 Let $a \in \mathbb{R}^{d}$, let $\mu$ be a finite symmetric measure on $\mathbb{S}^{d-1}$, and let $v$ be an admissible measure. Then,

$$
\begin{equation*}
\text { a.s. } \forall t \in \mathbb{R}^{d} \quad \alpha_{Y_{a, \mu, v}}(t)=\alpha_{B_{\mu}}(t) \wedge \alpha_{L_{v}}(t) \tag{23}
\end{equation*}
$$

Proof As linear drifts are $C^{\infty}$ everywhere, we may clearly assume that $a$ is zero. Moreover, note that the Hölder exponent of the sum of two functions is the minimum of the two exponents, except maybe when these exponents coincide, in which case the exponent of the sum may be larger. In view of Theorem 2, it remains to show that with probability one, $\alpha_{Y_{0, \mu, v}}(t) \leq 1 / 2$ for any $t \in \mathbb{R}^{d}$ with $\alpha_{L_{v}}(t)=1 / 2$.

To this end, it suffices to observe that Proposition 9 below (which corresponds to the upper bound in Proposition 6) still holds when replacing $L_{\nu}$ by $B_{\mu}+L_{\nu}$, because $B_{\mu}$ is continuous everywhere. Therefore, with probability one, $\alpha_{Y_{0, \mu, v}}(t) \leq \mathrm{A}_{\nu}(t)$ for any $t \in \mathbb{R}^{d}$ with $\alpha_{L_{v}}(t)=1 / 2$. We conclude by Proposition 6 .

## 8 Architecture of the proofs concerning the jump component

Let us now present the key ideas involved in the proofs of the main results of Sect. 6, which describe the size and large intersection properties of the iso-Hölder sets $E_{v}(h)$ and the singularity sets $E_{v}^{\prime}(h)$.

### 8.1 Location of the singularities

The first ingredient in our proofs is a precise determination of the location of the singularities of $L_{\nu}$. This follows from a characterization of the Hölder exponent of its sample paths in terms of the atoms of the Poisson measure N arising in the construction described in Sect. 4. In fact, the value of the exponent $\alpha_{L_{v}}(t)$ at a point $t \in \mathbb{R}^{d}$ depends on how well $t$ is approximated by the hyperplanes on which $L_{v}$ has a jump of size at most one.

To be more specific, for any real number $\alpha>0$, let us consider

$$
K_{\nu}(\alpha)=\left\{t \in \mathbb{R}^{d}\left|\mathrm{~d}\left(t, H_{n}\right)<\left|X_{n}\right|^{1 / \alpha} \text { for i.m. } n \geq 1 \text { with }\right| X_{n} \mid \leq 1\right\}
$$

where i.m. stands for "infinitely many" and $\mathrm{d}\left(t, H_{n}\right)$ denotes the distance, equal to $\left|\mathrm{P}_{n}-\left\langle S_{n}, t\right\rangle\right|$, between the point $t$ and the hyperplane $H_{n}$ defined by (9). In the terminology of Diophantine approximation, $K_{\nu}(\alpha)$ is the set of points that are approximated by the hyperplanes $H_{n}$ with a precision given by $\left|X_{n}\right|^{1 / \alpha}$. Equivalently,

$$
K_{\nu}(\alpha)=\left\{t \in \mathbb{R}^{d} \mid \int_{\substack{(\rho, s) \in \mathbb{H}_{d} \\|x| \in(0,1]}} \mathbb{1}_{\left\{|\rho-\langle s, t\rangle|<|x|^{1 / \alpha}\right\}} \mathrm{N}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x)=\infty\right\}
$$

Given that the mapping $\alpha \mapsto K_{\nu}(\alpha)$ is nondecreasing, it is possible to define

$$
\mathrm{A}_{v}(t)=\inf \left\{\alpha>0 \mid t \in K_{v}(\alpha)\right\}
$$

for each $t \in \mathbb{R}^{d}$. Proposition 7 below implies that with probability one, $K_{\nu}(\alpha)$ covers the whole space $\mathbb{R}^{d}$ when $\alpha>1 / \beta_{v}$. As a consequence,

$$
\begin{equation*}
\text { a.s. } \forall t \in \mathbb{R}^{d} \quad 0 \leq \mathrm{A}_{v}(t) \leq 1 / \beta_{v} \tag{24}
\end{equation*}
$$

The next result, which is proven in Sect. 9, gives a simple connection between the value of the Hölder exponent of $L_{v}$ at a given point $t$ and that of $\mathrm{A}_{v}(t)$.

Proposition 6 If $v$ is admissible, then

$$
\text { a.s. } \forall t \in \mathbb{R}^{d} \quad \alpha_{L_{v}}(t)= \begin{cases}0 & \text { if } t \in J_{v} \\ \mathrm{~A}_{v}(t) & \text { else. }\end{cases}
$$

Recall that $J_{v}$ is the set given by (18) and composed of the points at which the field $L_{v}$ jumps. Therefore the Hölder exponent of this field vanishes everywhere in $J_{v}$. Furthermore, it follows from Proposition 6 that with probability one, the Hölder exponent of $L_{\nu}$ is at most $1 / \beta_{\nu}$ everywhere. So, for any $h>1 / \beta_{\nu}$,

$$
\begin{equation*}
E_{v}^{\prime}(h)=\mathbb{R}^{d} \backslash J_{v} \quad \text { and } \quad E_{v}(h)=\emptyset, \tag{25}
\end{equation*}
$$

thus confirming some of the results announced in Proposition 4. When $h \leq 1 / \beta_{\nu}$, Proposition 6 enables us to relate $E_{v}(h)$ and $E_{v}^{\prime}(h)$ with the sets $K_{v}(\alpha)$ and $J_{v}$ as follows.

Corollary 5 If $v$ is admissible, then with probability one,

1. for any $h \in\left[0,1 / \beta_{\nu}\right]$,

$$
E_{v}^{\prime}(h)=\left(\bigcap_{h<\alpha \leq 1 / \beta_{v}} K_{\nu}(\alpha)\right) \backslash J_{v} ;
$$

2. $E_{\nu}(0)=J_{v} \cup E_{v}^{\prime}(0)$ and for any $h \in\left(0,1 / \beta_{v}\right]$,

$$
E_{\nu}(h)=E_{v}^{\prime}(h) \backslash \bigcup_{0<\alpha<h} K_{v}(\alpha) .
$$

We adopt here the standard convention that an intersection and a union indexed by the empty set are equal, respectively, to $\mathbb{R}^{d}$ and to the empty set. In particular, $E_{v}^{\prime}\left(1 / \beta_{v}\right)=\mathbb{R}^{d} \backslash J_{v}$ almost surely. The proof of Corollary 5, assuming that Proposition 6 holds, is straightforward and therefore omitted.

### 8.2 Diophantine approximation by Poisson hyperplanes

Corollary 5 suggests that the proofs of the results of Sect. 6 follow from a precise understanding of the size and large intersection properties of the sets $K_{v}(\alpha)$. In fact, these proofs make a crucial use of the next two results, see Sect. 12.2. The first one shows that $K_{v}(\alpha)=\mathbb{R}^{d}$ almost surely whenever $\alpha>1 / \beta_{v}$, thereby leading to (24). We refer to Sect. 10 for its proof.

Proposition 7 For any $\alpha>1 / \beta_{v}$, with probability one, $K_{v}(\alpha)=\mathbb{R}^{d}$.
The second one deals with the size and large intersection properties of $K_{v}(\alpha)$ and is proven in Sect. 11. Recall that $h_{v}(g)$ is defined by (19).

Theorem 7 For any $\alpha>0$, with probability one, for any gauge function $g \in \mathfrak{D}_{1}$,

$$
\left\{\begin{array}{l}
h_{v}(g)>\alpha \Longrightarrow \mathcal{H}^{d-1, g}\left(K_{v}(\alpha)\right)=0 \\
h_{v}(g)<\alpha \Longrightarrow K_{v}(\alpha) \in \mathrm{G}^{d-1, g}\left(\mathbb{R}^{d}\right) .
\end{array}\right.
$$

The fact that we make use of that result in the proof of Theorems 4 and 6 hints at why they describe the size and large intersection properties of the sets $E_{\nu}(h)$ and $E_{v}^{\prime}(h)$ only in terms of the gauge functions of the form $r \mapsto r^{d-1} g(r)$ with $g \in \mathfrak{D}_{1}$. This follows essentially from the ubiquity techniques that we use in Sect. 11 (see Lemma 9 in particular), where we regard $K_{\nu}(\alpha)$ as an enlargement, measured through $g$, of the random hyperplanes $H_{n}$, which are of dimension $d-1$.

## 9 Location of the singularities of the jump component

### 9.1 Preliminaries

The purpose of this section is to establish Proposition 6. The proof relies on suitable estimates of the increments of the random fields $L_{\nu, j}$ that come into play in the construction of $L_{v}$, as detailed in Sect. 4. To be specific, for all integers $A, j, k \geq 1$, let

$$
\begin{equation*}
\zeta_{v}(A, j, k)=\sup _{\substack{\|t\| \leq A \\\|\tau\| \leq 2^{-k}}}\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right| . \tag{26}
\end{equation*}
$$

Even though the supremum is taken over an uncountable set of parameters, there is no measurability issue here, in the sense that $\zeta_{v}(A, j, k)$ is a random variable. In fact, it is easy to check that the field $L_{v, j}$ is separable and that $\mathbb{Q}^{d}$ may be taken as the separant dense countable subset of $\mathbb{R}^{d}$ involved in this property. Accordingly, in addition to (26), we have

$$
\begin{equation*}
\zeta_{\nu}(A, j, k)=\sup _{\substack{t \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}(A) \\ \tau \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}\left(2^{-k}\right)}}\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right| . \tag{27}
\end{equation*}
$$

The next lemma yields an upper bound on $\zeta_{\nu}(A, j, k)$, that is, a control of the increments of $L_{\nu, j}$. It is proven in Sect. 13, and also comes into play in the proof of Proposition 1, see Sect. 12.1. Recall that $\left(v_{j}\right)_{j \geq 1}$ is the sequence arising in (3).

Lemma 2 For every integer $A \geq 1$, with probability one,

$$
Z_{v}(A)=\sup _{(j, k) \in \mathbb{N}^{2}} \frac{\zeta_{v}(A, j, k)}{2^{-j} k\left(j+2^{-k / 2}\left(j v_{j}\right)^{1 / 2}\right)}<\infty
$$

In view of Lemma 2, we shall work below on the almost sure event consisting in the finiteness of $\mathrm{Z}_{v}(A)$, for all $A \geq 1$. Moreover, recall that ( $\mathrm{P}_{n}, S_{n}, X_{n}$ ), for $n \geq 1$, are the atoms of the Poisson random measure N arising in the construction of $L_{\nu}$. We may assume that the real numbers $\mathrm{P}_{n}$ are distinct. Indeed, for any $j \geq 1$, the image under $(\rho, s, x) \mapsto \rho$ of the restriction of N to $(0, \infty) \times \mathbb{S}^{d-1} \times\left(\mathbb{R} \backslash\left[-2^{-j}, 2^{-j}\right]\right)$ is a Poisson measure on $(0, \infty)$ with intensity proportional to $\mathcal{L}_{+}^{1}$, thereby being almost surely simple, see [34, p. 299]. In addition, we define $L_{\nu}^{j_{1}, j_{2}}=\sum_{j=j_{1}}^{j_{2}} L_{\nu, j}$ for $0 \leq j_{1} \leq j_{2} \leq \infty$.

We now split the proof of Proposition 6 into three parts. First, we show that $\alpha_{L_{v}}(t)$ vanishes at every jump point $t \in J_{v}$. Second, we show that $\alpha_{L_{v}}(t) \leq \mathrm{A}_{v}(t)$ at every $t \notin J_{v}$. Third, we show that $\alpha_{L_{v}}(t) \geq \mathrm{A}_{v}(t)$ at any such $t$. Throughout, we assume that the measure $v$ is admissible.

### 9.2 Value of the Hölder exponent at the jump points

With regard to the next statement, recall that the hyperplanes $H_{n}$ are defined by (9).
Lemma 3 Almost surely, for any $t \in \mathbb{R}^{d}$ such that $t \in H_{n_{0}}$ for a unique $n_{0} \geq 1$,

$$
\lim _{\ell \rightarrow \infty} L_{v}\left(t+\frac{S_{n_{0}}}{\ell}\right)=L_{v}(t)+X_{n_{0}}
$$

Proof Given an integer $A \geq 1$, let $t$ be a point in $\overline{\mathrm{B}}_{0}(A)$ such that $t \in H_{n_{0}}$ for a unique $n_{0} \geq 1$. Then, with probability one, for all integers $j_{0}>-\log _{2}\left|X_{n_{0}}\right|$ and $\ell$ larger than some $\ell_{0} \geq 1$, the only $n \geq 1$ satisfying $\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t+S_{n_{0}} / \ell} \Delta V_{t}$ and $\left|X_{n}\right|>2^{-j_{0}}$ simultaneously is $n_{0}$, and in fact $\left(\mathrm{P}_{n_{0}}, S_{n_{0}}\right) \in V_{t+S_{n_{0}} / \ell} \backslash V_{t}$. Here, $\Delta$ stands for symmetric difference of sets, and both $V_{t}$ and $V_{t+S_{n_{0}} / \ell}$ are given by (5). Hence, for all such $j_{0}$ and $\ell$, we have

$$
L_{v}^{0, j_{0}}\left(t+\frac{S_{n_{0}}}{\ell}\right)-L_{v}^{0, j_{0}}(t)=X_{n_{0}}-\frac{1}{\ell} \int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in\left(2^{-j}, 1\right]}} x\left\langle s, S_{n_{0}}\right\rangle v(\mathrm{~d} s, \mathrm{~d} x)
$$

thanks to (8) and (11). In addition, given that $\|t\| \leq A$ and $\left\|S_{n_{0}} / \ell\right\| \leq 1 / 2$, we get

$$
\left|L_{v}^{j_{0}+1, \infty}\left(t+\frac{S_{n_{0}}}{\ell}\right)-L_{v}^{j_{0}+1, \infty}(t)\right| \leq \sum_{j=j_{0}+1}^{\infty} \zeta_{v}(A, j, 1)
$$

by virtue of (26). As a consequence, making use of Lemma 2, we deduce that

$$
\left|L_{v}\left(t+\frac{S_{n_{0}}}{\ell}\right)-L_{v}(t)-X_{n_{0}}\right| \leq \frac{1}{\ell} \sum_{j=1}^{j_{0}} v_{j}+\mathrm{Z}_{v}(A) \sum_{j=j_{0}+1}^{\infty} 2^{-j}\left(j+\left(j v_{j}\right)^{1 / 2}\right),
$$

and conclude by letting $\ell \rightarrow \infty$, and then by letting $j_{0} \rightarrow \infty$ while using of the fact that the sum $\chi_{\nu}$ defined by (3) is finite, because $v$ is admissible.

Thanks to Lemma 3, we may now prove that the Hölder exponent of $L_{v}$ vanishes at every jump point.

Proposition 8 With probability one, $\alpha_{L_{v}}(t)=0$ for every $t \in J_{v}$.
Proof Let $t \in J_{v}$ and suppose that $\alpha_{L_{v}}(t)>0$. So, there exist $\varepsilon, \delta, C>0$ such that for any $\tau \in \mathbb{R}^{d}$,

$$
\begin{equation*}
\|\tau\| \leq \delta \Longrightarrow\left|L_{v}(t+\tau)-L_{v}(t)\right| \leq C\|\tau\|^{\varepsilon} . \tag{28}
\end{equation*}
$$

Since $t \in J_{v}$, there is an $n_{0} \geq 1$ such that $t \in H_{n_{0}}$. However, $n_{0}$ need not be unique, and we cannot apply Lemma 3 directly. To cope with that problem, recall that the real numbers $\mathrm{P}_{n}$ are distinct, so that the hyperplanes $H_{n}$ are distinct too. Hence, for any integer $m \geq 1$, the set $\mathrm{B}_{t}(1 / m) \cap\left(H_{n_{0}} \backslash \bigcup_{n \neq n_{0}} H_{n}\right)$ contains a point $t_{m}$. Here, $\mathrm{B}_{t}(1 / m)$ is the open Euclidean ball centered at $t$ with radius $1 / m$. Then, $n_{0}$ is the only integer such that $t_{m} \in H_{n_{0}}$. Applying Lemma 3 with $t_{m}$, we get

$$
\begin{equation*}
\forall m \geq 1 \quad \lim _{\ell \rightarrow \infty} L_{v}\left(t_{m}+\frac{S_{n_{0}}}{\ell}\right)=L_{v}\left(t_{m}\right)+X_{n_{0}} \tag{29}
\end{equation*}
$$

Now, for all integers $m, \ell \geq 1$, we have $\left\|t_{m}-t\right\|<1 / m$ and $\left\|t_{m}+S_{n_{0}} / \ell-t\right\|<$ $1 / m+1 / \ell$. So, assuming that $1 / m+1 / \ell<\delta$ and applying (28), we obtain

$$
\left|L_{v}\left(t_{m}+\frac{S_{n_{0}}}{\ell}\right)-L_{v}\left(t_{m}\right)\right| \leq C\left(\frac{1}{m^{\varepsilon}}+\left(\frac{1}{m}+\frac{1}{\ell}\right)^{\varepsilon}\right) .
$$

Letting $\ell \rightarrow \infty$ and using (29), we infer that $\left|X_{n_{0}}\right| \leq 2 C / m^{\varepsilon}$ for any $m>1 / \delta$. Then, letting $m \rightarrow \infty$, we get $X_{n_{0}}=0$, which contradicts the fact that $X_{n_{0}} \in \mathbb{R}^{*}$.
9.3 Upper bound on the Hölder exponent

We now consider the points at which the field $L_{v}$ does not jump.
Proposition 9 With probability one, $\alpha_{L_{v}}(t) \leq \mathrm{A}_{v}(t)$ for every $t \in \mathbb{R}^{d} \backslash J_{v}$.
Proof Let $t \in \mathbb{R}^{d} \backslash J_{v}$ and $\alpha>\mathrm{A}_{v}(t)$, and suppose that $\alpha_{L_{v}}(t) \geq \alpha+\varepsilon$ for some $\varepsilon>0$. First, there are $\delta, C>0$ and a polynomial $Q_{t}$ such that for any $\tau \in \mathbb{R}^{d}$,

$$
\begin{equation*}
\|\tau\| \leq \delta \Longrightarrow\left|L_{v}(t+\tau)-Q_{t}(\tau)\right| \leq C\|\tau\|^{\alpha+\varepsilon} . \tag{30}
\end{equation*}
$$

Second, $t \in K_{\nu}(\alpha) \backslash J_{\nu}$, so that for any $B \in(0,1]$ with $B^{-\alpha / \varepsilon}>3 C 2^{\alpha+\varepsilon}$ and $2 B^{1 / \alpha}<\delta$, there is an $n_{0} \geq 1$ with $\left|X_{n_{0}}\right| \leq B$ and $0<\mathrm{d}\left(t, H_{n_{0}}\right) \leq\left|X_{n_{0}}\right|^{1 / \alpha}$.

Let us suppose that for any $\tau \in \mathbb{R}^{d}$,

$$
\begin{equation*}
\|\tau\|<2 \mathrm{~d}\left(t, H_{n_{0}}\right) \Longrightarrow 3\left|L_{v}(t+\tau)-Q_{t}(\tau)\right|<\left|X_{n_{0}}\right| . \tag{31}
\end{equation*}
$$

As in the proof of Proposition 8, the set $\mathrm{B}_{t}\left(2 \mathrm{~d}\left(t, H_{n_{0}}\right)\right) \cap\left(H_{n_{0}} \backslash \bigcup_{n \neq n_{0}} H_{n}\right)$ contains a point $t^{\prime}$. For $\ell$ large enough, $\left\|t^{\prime}-t\right\|$ and $\left\|t^{\prime}+S_{n_{0}} / \ell-t\right\|$ are both smaller than $2 \mathrm{~d}\left(t, H_{n_{0}}\right)$, so that (31) leads to

$$
\left|L_{v}\left(t^{\prime}+\frac{S_{n_{0}}}{\ell}\right)-L_{v}\left(t^{\prime}\right)\right| \leq \frac{2}{3}\left|X_{n_{0}}\right|+\left|Q_{t}\left(t^{\prime}+\frac{S_{n_{0}}}{\ell}-t\right)-Q_{t}\left(t^{\prime}-t\right)\right|
$$

Since $n_{0}$ is the only integer such that $t^{\prime} \in H_{n_{0}}$, it follows from Lemma 3 that the left-hand side tends to $\left|X_{n_{0}}\right|$ as $\ell \rightarrow \infty$. This contradicts the fact that the right-hand side goes to $2\left|X_{n_{0}}\right| / 3$. As a result, there is a $\tau \in \mathbb{R}^{d}$ for which (31) does not hold.

Therefore, we have $\|\tau\|<2 \mathrm{~d}\left(t, H_{n_{0}}\right) \leq 2\left|X_{n_{0}}\right|^{1 / \alpha} \leq 2 B^{1 / \alpha}<\delta$ and

$$
\left|X_{n_{0}}\right| \leq 3\left|L_{v}(t+\tau)-Q_{t}(\tau)\right| \leq 3 C\|\tau\|^{\alpha+\varepsilon} \leq 3 C 2^{\alpha+\varepsilon}\left|X_{n_{0}}\right|^{1+\varepsilon / \alpha},
$$

thanks to (30). This implies that $B^{-\alpha / \varepsilon} \leq 3 C 2^{\alpha+\varepsilon}$, which contradicts the choice of $B$. Finally, $\alpha_{L_{v}}(t) \leq \alpha$, and we conclude by letting $\alpha \downarrow \mathrm{A}_{v}(t)$.

### 9.4 Lower bound on the Hölder exponent

The following result remains to be established. Its proof is split into several parts for the sake of clarity.

Proposition 10 With probability one, $\alpha_{L_{v}}(t) \geq \mathrm{A}_{v}(t)$ for every $t \in \mathbb{R}^{d} \backslash J_{v}$.
In view of (1), with probability one, for any integer $A \geq 1$ and any $\varepsilon>0$, there are only finitely many $n \geq 1$ such that $\mathrm{P}_{n}<A$ and $\left|X_{n}\right|>\varepsilon$ simultaneously. We may therefore suppose, in addition to the assumptions made at the beginning of this section, that the corresponding almost sure event occurs. We may also assume that the almost sure event given by (24) occurs too.

Now, given an integer $A \geq 1$, let $t \in \mathrm{~B}_{0}(A) \backslash J_{v}$. To show that $\alpha_{L_{v}}(t) \geq \mathrm{A}_{v}(t)$, we may obviously assume that $\mathrm{A}_{v}(t)>0$. Then, let $\alpha \in\left(0, \mathrm{~A}_{v}(t)\right)$. As $t \notin K_{v}(\alpha)$, there are only finitely many $n \geq 1$ such that $\mathrm{d}\left(t, H_{n}\right)<\left|X_{n}\right|^{1 / \alpha} \leq 1$. Hence, there is an integer $k_{0} \geq 1$ such that $\mathrm{d}\left(t, H_{n}\right) \geq\left|X_{n}\right|^{1 / \alpha}$ for any $n \geq 1$ with $\left|X_{n}\right| \leq 2^{-\left\lfloor\alpha k_{0}\right\rfloor}$. Besides, note that $\alpha<1 / \beta_{v}$, owing to (24).

### 9.4.1 Reduction to the study of the component with small jumps

Thanks to (8) and (11), the value at $t$ of the component of the field $L_{\nu}$ that corresponds to the jumps of size larger than $2^{-\left\lfloor\alpha k_{0}\right\rfloor}$ may be written as

$$
L_{\nu}^{0,\left\lfloor\alpha k_{0}\right\rfloor}(t)=\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t}, \mid X_{n}\right\rfloor>2^{\left.-\left\lfloor\alpha k_{0}\right\rfloor\right\}}}-\int_{\substack{\left.s \in \mathbb{S}^{d-1} \\ x \in\left(2^{-\left\lfloor\alpha k_{0}\right\rfloor}\right\rfloor, 1\right\rfloor}} x\langle s, t\rangle \nu(\mathrm{d} s, \mathrm{~d} x)
$$

where $V_{t}$ is defined by (5). Moreover, for any $\delta \in(0, A-\|t\|)$, let $V_{t, \delta}$ denote the complement of $\bigcap_{\|\tau\| \leq \delta} V_{t+\tau}$ in $\bigcup_{\|\tau\| \leq \delta} V_{t+\tau}$. If $n \geq 1$ satisfies $\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t, \delta}$, then $\mathrm{P}_{n}<A$. Thus, the set $\mathcal{N}_{\delta}$ of all $n \geq 1$ such that $\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t, \delta}$ and $\left|X_{n}\right|>2^{-\left\lfloor\alpha k_{0}\right\rfloor}$ is finite. Moreover, given that $t \notin J_{v}$, it is clear that $\bigcap_{\delta>0} \downarrow \mathcal{N}_{\delta}=\emptyset$, so that $\mathcal{N}_{\delta}=\emptyset$ for $\delta$ small enough. For such a $\delta$ and for $\|\tau\| \leq \delta$, we have $V_{t+\tau} \Delta V_{t} \subseteq V_{t, \delta}$. Hence, no integer $n \geq 1$ can satisfy both $\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t} \Delta V_{t+\tau}$ and $\left|X_{n}\right|>2^{-\left\lfloor\alpha k_{0}\right\rfloor}$, so

$$
L_{v}^{0,\left\lfloor\alpha k_{0}\right\rfloor}(t+\tau)-L_{v}^{0,\left\lfloor\alpha k_{0}\right\rfloor}(t)=-\int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in\left(2^{-\left\lfloor\alpha k_{0}\right\rfloor}, 1\right]}} x\langle s, \tau\rangle v(\mathrm{~d} s, \mathrm{~d} x)
$$

Hence, $L_{v}^{0,\left\lfloor\alpha k_{0}\right\rfloor}$ coincides with an affine form near $t$, and the Hölder exponent at $t$ of $L_{v}$ is equal to that of the component with jumps of size at most $2^{-\left\lfloor\alpha k_{0}\right\rfloor}$.

### 9.4.2 Study of the component with small jumps

In order to study the Hölder exponent of the component with jumps of size at most $2^{-\left\lfloor\alpha k_{0}\right\rfloor}$, let us consider a vector $\tau \in \mathbb{R}^{d}$ such that $2^{-(k+1)}<\|\tau\| \leq 2^{-k}$ for some integer $k \geq k_{0}$. First,

$$
\left|L_{\nu}^{\lfloor\alpha k\rfloor+1, \infty}(t+\tau)-L_{\nu}^{\lfloor\alpha k\rfloor+1, \infty}(t)\right| \leq \mathrm{Z}_{\nu}(A) \sum_{j=\lfloor\alpha k\rfloor+1}^{\infty} 2^{-j} k\left(j+2^{-k / 2}\left(j v_{j}\right)^{1 / 2}\right)
$$

in view of Lemma 2. Furthermore, if $\beta_{\nu}<2$, then there exists $\gamma \in\left(\beta_{v},(1 / \alpha) \wedge 2\right)$ and (2) implies that $c_{v, \gamma}=\sum_{j \geq 1} 2^{-\gamma j} j v_{j}$ is finite. Hence, we have
$k 2^{-k / 2} \sum_{j=\lfloor\alpha k\rfloor+1}^{\infty} 2^{-j}\left(j v_{j}\right)^{1 / 2} \leq k 2^{-k / 2} c_{\nu, \gamma}^{1 / 2} \sum_{j=\lfloor\alpha k\rfloor+1}^{\infty} 2^{(\gamma / 2-1) j} \leq \frac{c_{\nu, \gamma}^{1 / 2}}{1-2^{\gamma / 2-1}} k 2^{-\alpha k}$.

If $\beta_{\nu}=2$, we observe that $\alpha<1 / 2$, so that the left-hand side above is at most $\chi_{v} k 2^{-\alpha k}$, where $\chi_{v}$ is finite and defined by (3). The upshot is that there exists a deterministic real number $D_{v, \alpha}>0$ that depends on $v$ and $\alpha$ only such that

$$
\begin{equation*}
\left|L_{\nu}^{\lfloor\alpha k\rfloor+1, \infty}(t+\tau)-L_{\nu}^{\lfloor\alpha k\rfloor+1, \infty}(t)\right| \leq \mathrm{Z}_{v}(A) D_{v, \alpha} k^{2} 2^{-\alpha k} . \tag{32}
\end{equation*}
$$

Second, no integer $n \geq 1$ can verify at the same time $2^{-\lfloor\alpha k\rfloor}<\left|X_{n}\right| \leq 2^{-\left\lfloor\alpha k_{0}\right\rfloor}$ and $\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t} \Delta V_{t+\tau}$ (otherwise, $\left|\mathrm{P}_{n}-\left\langle S_{n}, t\right\rangle\right|$ would be at most $2^{-k}$ and at least $\left|X_{n}\right|^{1 / \alpha}>2^{-k}$ simultaneously, which is impossible). Along with (11), this yields

$$
\left|L_{\nu}^{\left\lfloor\alpha k_{0}\right\rfloor+1,\lfloor\alpha k\rfloor}(t+\tau)-L_{\nu}^{\left\lfloor\alpha k_{0}\right\rfloor+1,\lfloor\alpha k\rfloor}(t)\right| \leq\|\tau\| \int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in\left(2^{\left.-\lfloor\alpha k\rfloor, 2^{-\left\lfloor\alpha k_{0}\right\rfloor}\right\rfloor}\right.}} x v(\mathrm{~d} s, \mathrm{~d} x) .
$$

### 9.4.3 End of the proof for $\beta_{\nu} \geq 1$

There exists $\gamma \in[1,2]$ such that $\gamma<1 / \alpha$ and the integral $I_{\gamma}$ of $(s, x) \mapsto x^{\gamma}$ over $\mathbb{S}^{d-1} \times(0,1]$ with respect to $v$ is finite. Indeed, one may choose $\gamma=\beta_{v}$ if $\beta_{v}=2$ and $\gamma>\beta_{v}$ sufficiently small otherwise. Hence, using both (32) and the above bound, we infer that

$$
\begin{aligned}
\left|L_{\nu}^{\left\lfloor\alpha k_{0}\right\rfloor+1, \infty}(t+\tau)-L_{\nu}^{\left\lfloor\alpha k_{0}\right\rfloor+1, \infty}(t)\right| & \leq\|\tau\| I_{\gamma} 2^{(\gamma-1)\lfloor\alpha k\rfloor}+\mathrm{Z}_{v}(A) D_{\nu, \alpha} k^{2} 2^{-\alpha k} \\
& \leq\|\tau\|^{\alpha}\left(I_{\gamma}+2^{\alpha} Z_{v}(A) D_{\nu, \alpha}\left(\log _{2}\|\tau\|\right)^{2}\right) .
\end{aligned}
$$

It follows that $\alpha_{L_{v}}(t) \geq \alpha$. To conclude, it remains to let $\alpha \uparrow \mathrm{A}_{v}(t)$.

### 9.4.4 End of the proof for $\beta_{v}<1$

Here, thanks to (11),

$$
L_{\nu}^{\left\lfloor\alpha k_{0}\right\rfloor+1, \infty}(t)=\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t},\left|X_{n}\right| \leq 2^{-\left\lfloor\alpha k_{0}\right\rfloor}\right\}}-\int_{\substack{s \in \mathbb{S}^{d-1} \\ x \in\left(0,2^{-\left\lfloor\alpha k_{0}\right\rfloor_{]}}\right.}} x\langle s, t\rangle v(\mathrm{~d} s, \mathrm{~d} x)
$$

The second term is a linear form, so we just need to study the increments of the first term. To this end, observe that for an arbitrary $\gamma \in\left(\beta_{\nu},(1 / \alpha) \wedge 1\right)$,

$$
\begin{aligned}
& \left|\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t+\tau},\left|X_{n}\right| \leq 2^{\left.-\left\lfloor\alpha k_{0}\right\rfloor\right\}}\right.}-\sum_{n=1}^{\infty} X_{n} \mathbb{1}_{\left\{\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t},\left|X_{n}\right| \leq 2^{\left.-\left\lfloor\alpha k_{0}\right\rfloor\right\}}\right.}\right| \\
& \leq\left|L_{\nu}^{\lfloor\alpha k\rfloor+1, \infty}(t+\tau)-L_{\nu}^{\lfloor\alpha k\rfloor+1, \infty}(t)\right|+\|\tau\| \int_{\substack{s \in \mathbb{S}^{d-1} \\
x \in\left(0,2^{-\lfloor\alpha k\rfloor}\right]}} x v(\mathrm{~d} s, \mathrm{~d} x)
\end{aligned}
$$

$$
\begin{aligned}
& \leq \mathrm{Z}_{v}(A) D_{v, \alpha} k^{2} 2^{-\alpha k}+\|\tau\| I_{\gamma} 2^{-(1-\gamma)\lfloor\alpha k\rfloor} \\
& \leq\|\tau\|^{\alpha}\left(2^{(1-\gamma)(1+\alpha)} I_{\gamma}+2^{\alpha} Z_{v}(A) D_{v, \alpha}\left(\log _{2}\|\tau\|\right)^{2}\right)
\end{aligned}
$$

because of (32) and the fact that no integer $n \geq 1$ can satisfy $\left(\mathrm{P}_{n}, S_{n}\right) \in V_{t} \Delta V_{t+\tau}$ and $2^{-\lfloor\alpha k\rfloor}<\left|X_{n}\right| \leq 2^{-\left\lfloor\alpha k_{0}\right\rfloor}$ simultaneously. We deduce that $\alpha_{L_{v}}(t) \geq \alpha$, and conclude by letting $\alpha \uparrow \mathrm{A}_{\nu}(t)$.

## 10 Approximation by Poisson hyperplanes: covering the whole space

We now prove Proposition 7. To begin with, given $\alpha>1 / \beta_{v}$, let $v_{\alpha}$ denote the image under $(s, x) \mapsto\left(s,|x|^{1 / \alpha}\right)$ of the restriction to $\mathbb{S}^{d-1} \times([-1,1] \backslash\{0\})$ of the measure $\nu$. Then, for an arbitrary orthonormal basis $\left(e_{1}, \ldots, e_{d}\right)$ of $\mathbb{R}^{d}$ and for any $s \in \mathbb{S}^{d-1}$, there necessarily exists an integer $i \in\{1, \ldots, d\}$ such that $\left\langle s, e_{i}\right\rangle \neq 0$. Together with (2), this shows that for some $i$ and some $\varepsilon>0$,

$$
\begin{equation*}
\int_{\substack{s \in \mathbb{S}^{d-1} \\ r \in(0,1]}} \mathbb{1}_{\left\{\left\langle s, e_{i}\right\rangle \neq 0\right\}} r^{1+\varepsilon} v_{\alpha}(\mathrm{d} s, \mathrm{~d} r)=\infty . \tag{33}
\end{equation*}
$$

Now, for any integers $A, j \geq 1$, let $U_{A, j}=\mathrm{B}_{0}(A) \cap\left(2^{-j} / \sqrt{d}\right) \mathbb{Z}^{d}$. Moreover, for $j>j_{0} \geq 1$, let $\mathcal{E}_{A, j_{0}, j}$ denote the event consisting in the existence of a point $u \in U_{A, j}$ satisfying $\mathrm{d}\left(u, H_{n}\right) \geq\left|X_{n}\right|^{1 / \alpha}-2^{-j}$ for any integer $n \geq 1$ with $\left\langle S_{n}, e_{i}\right\rangle \neq 0$ and $2^{-j}<\left|X_{n}\right|^{1 / \alpha} \leq 2^{-j_{0}}$. Then, we have

$$
\begin{equation*}
\left\{K_{\nu}(\alpha) \neq \mathbb{R}^{d}\right\} \subseteq \bigcup_{A=1}^{\infty} \uparrow \bigcup_{j_{0}=1}^{\infty} \uparrow \bigcap_{j=j_{0}+1}^{\infty} \mathcal{E}_{A, j_{0}, j} \tag{34}
\end{equation*}
$$

The event $\mathcal{E}_{A, j_{0}, j}$ occurs with probability at most $\sum_{u \in U_{A, j}} \mathrm{e}^{-I_{j_{0}, j, u}}$, where

$$
I_{j_{0}, j, u}=\int_{\substack{(\rho, s) \in \mathbb{H}_{d} \\ r \in(0,1]^{2}}} f_{j_{0}, j, u}(\rho, s, r) \mathrm{d} \rho v_{\alpha}(\mathrm{d} s, \mathrm{~d} r)
$$

and $f_{j_{0}, j, u}(\rho, s, r)$ is equal to one when $|\rho-\langle s, u\rangle|<r-2^{-j},\left\langle s, e_{i}\right\rangle \neq 0$ and $2^{-j}<r \leq 2^{-j_{0}}$, and is equal to zero otherwise. Using the symmetry of $\nu$, we infer that $I_{j_{0}, j, u}$ is equal to

$$
\int_{\substack{\left\langle s, e_{i}\right\rangle>0 \\ r \in\left(2^{-j}, 2^{-j_{0}}\right]}} \int_{\substack{ \\\rho \in \mathbb{R}}} \mathbb{1}_{\left\{|\rho-\langle s, u\rangle|<r-2^{-j}\right\}} \mathrm{d} \rho v_{\alpha}(\mathrm{d} s, \mathrm{~d} r)=\int_{\substack{\left\langle s, e_{i}\right\rangle \neq 0 \\ r \in\left(2^{-j}, 2^{-j_{0}}\right]}}\left(r-2^{-j}\right) \nu_{\alpha}(\mathrm{d} s, \mathrm{~d} r) .
$$

Let $I_{j_{0}, j}^{\prime}=v_{\alpha}\left(\left\{s \in \mathbb{S}^{d-1} \mid\left\langle s, e_{i}\right\rangle \neq 0\right\} \times\left(2^{-j}, 2^{-j_{0}}\right]\right)$. Then, due to Fubini's theorem,

$$
I_{j_{0}, j, u}=\int_{w \in\left(2^{-j}, 2^{-j_{0}}\right)} \int_{\substack{\left\langle s, e_{i}\right\rangle \neq 0 \\ r \in\left(w, 2^{-j_{0}}\right]}} v_{\alpha}(\mathrm{d} s, \mathrm{~d} r) \mathrm{d} w \geq 2^{-j} I_{j_{0}, j-1}^{\prime}
$$

Therefore, given that $U_{A, j}$ has cardinality at most $\left(2^{j+2} A \sqrt{d}\right)^{d}$, we get

$$
\begin{equation*}
\mathbb{P}\left(\mathcal{E}_{A, j_{0}, j}\right) \leq(4 A \sqrt{d})^{d} \exp \left(j d \log 2-2^{-j} I_{j_{0}, j-1}^{\prime}\right) . \tag{35}
\end{equation*}
$$

Finally, employing Fubini's theorem again, we also have

$$
\begin{aligned}
\int_{\substack{\left\langle s, e_{i}\right\rangle \neq 0 \\
r \in\left(0,2^{\left.-j_{0}\right]}\right.}} r^{1+\varepsilon} v_{\alpha}(\mathrm{d} s, \mathrm{~d} r) & =(1+\varepsilon) \int_{w \in\left(0,2^{-j_{0}}\right)} w^{\varepsilon} \int_{\substack{\left\langle s, e_{i}\right\rangle \neq 0 \\
r \in\left(w, 2^{-j} j_{0}\right.}} v_{\alpha}(\mathrm{d} s, \mathrm{~d} r) \mathrm{d} w \\
& \leq(1+\varepsilon) \sum_{j=j_{0}}^{\infty} I_{j_{0}, j+1}^{\prime} \int_{w \in\left(2^{-(j+1)}, 2^{-j}\right)} w^{\varepsilon} \mathrm{d} w
\end{aligned}
$$

Together with (33), this implies that $I_{j_{0}, j-1}^{\prime}>2^{(1+\varepsilon) j} / j^{2}$ for infinitely many $j>j_{0}$. We conclude with the help of (34) and (35).

## 11 Approximation by Poisson hyperplanes: size and large intersection properties

This section is devoted to the proof of Theorem 7. We begin by establishing a series of preliminary lemmas. Then, we deal with the case where $h_{\nu}(g)>\alpha$. We finally end the proof with the case where $h_{v}(g)<\alpha$. Note that Theorem 7 clearly holds when $v$ has finite total mass. Indeed, in this case, it is easy to check that $h_{v}(g)=\infty$ for any $g \in \mathfrak{D}_{1}$, while $K_{\nu}(\alpha)=\emptyset$ with probability one for every $\alpha>0$ (given $A \geq 1$, there are almost surely finitely many $n \geq 1$ such that $H_{n} \cap \overline{\mathrm{~B}}_{0}(A) \neq \emptyset$ or, equivalently, such that $\mathrm{P}_{n} \leq A$ ). Therefore, we may assume throughout the section that $v$ has infinite total mass.

### 11.1 A Bernstein-type inequality

The first preliminary lemma yields an analog of Bernstein's inequality for integrals with respect to a compensated Poisson random measure, and is a direct consequence of [23, Corollary 5.1], see also [35, Proposition 7]. It comes into play in the proof of Lemma 2 too, see Sect. 13.

Lemma 4 Let $(E, \mathcal{E})$ be a measurable space endowed with a finite nonnegative measure $\mu$, let M be a Poisson random measure on $E$ with intensity $\mu$ and let $\mathrm{M}^{*}=\mathrm{M}-\mu$.

Then, for any real-valued measurable function $f$ defined on $E$ such that $S=\sup _{E}|f|$ and $V=\int_{E} f^{2} \mathrm{~d} \mu$ are both positive and finite, we have

$$
\forall \xi>0 \quad \mathbb{P}\left(\left|\int_{E} f \mathrm{dM}^{*}\right| \geq \xi\right) \leq 2 \exp \left(-\frac{3 \xi^{2}}{2 S \xi+6 V}\right)
$$

### 11.2 A law of large numbers for Poisson measures

The second preliminary lemma concerns the behavior at zero of a Poisson random measure on $(0,1]$ and is reminiscent of the strong law of large numbers for the homogeneous Poisson process on $(0, \infty)$, see [27]. Let $\mathcal{P}$ be the set of all nonnegative Borel measures $\pi$ on $(0,1]$ such that $\pi((0,1])=\infty$ and $\pi([\varepsilon, 1])<\infty$ for any $\varepsilon>0$.

Lemma 5 For any Poisson random measure $\Pi$ on $(0,1]$ with intensity $\pi \in \mathcal{P}$, with probability one,

$$
\Pi([w, 1]) \sim \pi([w, 1]) \text { as } w \rightarrow 0
$$

Proof Let $\mathcal{A}_{0}$ be the (countable) set of all $r \in(0,1]$ such that $\pi(\{r\}) \geq 1$ and let $\mathcal{A}_{1}=(0,1] \backslash \mathcal{A}_{0}$. Then, for all $\ell \in\{0,1\}$ and $w \in(0,1]$, let $\Phi_{\ell}(w)=\Pi\left(\mathcal{A}_{\ell} \cap[w, 1]\right)$ and $\varphi_{\ell}(w)=\pi\left(\mathcal{A}_{\ell} \cap[w, 1]\right)$. It is easy to see that the proof reduces to showing that for any $\ell \in\{0,1\}$ such that $\pi\left(\mathcal{A}_{\ell}\right)=\infty$, with probability one,

$$
\begin{equation*}
\Phi_{\ell}(w) \sim \varphi_{\ell}(w) \text { as } \quad w \rightarrow 0 \tag{36}
\end{equation*}
$$

To this purpose, let us begin by observing that for any $\xi>0$ and any $w>0$ small enough to ensure that $\varphi_{\ell}(w)>0$,

$$
\begin{equation*}
\mathbb{P}\left(\left|\frac{\Phi_{\ell}(w)}{\varphi_{\ell}(w)}-1\right| \geq \xi\right) \leq 2 \exp \left(-\frac{3 \xi^{2}}{2 \xi+6} \varphi_{\ell}(w)\right) \tag{37}
\end{equation*}
$$

as a consequence of Lemma 4 . Now, if $\pi\left(\mathcal{A}_{0}\right)=\infty$, there exists a decreasing sequence $\left(a_{n}\right)_{n \geq 1}$ of positive real numbers that converges to zero and whose terms form the set $\mathcal{A}_{0}$. The previous inequality then implies that for all integers $m, n \geq 1$,

$$
\begin{equation*}
\mathbb{P}\left(\left|\frac{\Phi_{0}\left(a_{n}\right)}{\varphi_{0}\left(a_{n}\right)}-1\right| \geq \frac{1}{m}\right) \leq 2 \exp \left(-\frac{3 n}{(6 m+2) m}\right) \tag{38}
\end{equation*}
$$

because $\varphi_{0}\left(a_{n}\right)=\pi\left(\left\{a_{1}, \ldots, a_{n}\right\}\right) \geq n$. Summing over $n \geq 1$ and making use of the Borel-Cantelli lemma, we infer that for any integer $m \geq 1$, with probability one, $\Phi_{0}\left(a_{n}\right) / \varphi_{0}\left(a_{n}\right)$ is between $1-1 / m$ and $1+1 / m$, for $n$ large enough. The same property holds for $\Phi_{0}(w) / \varphi_{0}(w)$ with $w>0$ small enough, due to the fact that $\Phi_{0}(w)=\Phi_{0}\left(a_{n(w)}\right)$ and $\varphi_{0}(w)=\varphi_{0}\left(a_{n(w)}\right)$, where $n(w)$ is the number of integers
$n \geq 1$ such that $a_{n} \geq w$. As a consequence,

$$
\text { a.s. } \forall m \geq 1 \quad \limsup _{w \rightarrow 0}\left|\frac{\Phi_{0}(w)}{\varphi_{0}(w)}-1\right| \leq \frac{1}{m} \text {, }
$$

and we get (36) for $\ell=0$ by letting $m \rightarrow \infty$.
If $\pi\left(\mathcal{A}_{1}\right)=\infty$, it is possible to consider, for each integer $n \geq 1$,

$$
w_{n}=\sup \left\{w>0 \mid \varphi_{1}(w) \geq n\right\} .
$$

The reals $w_{n}$ are positive, satisfy $\varphi_{1}\left(w_{n}\right) \geq n$ and form a nonincreasing sequence that converges to zero. Moreover, (37) ensures that for any $m, n \geq 1$, the bound (38) holds with $a_{n}$ replaced by $w_{n}$, and $\Phi_{0}$ and $\varphi_{0}$ replaced by $\Phi_{1}$ and $\varphi_{1}$, respectively. Summing over $n \geq 1$ and using the Borel-Cantelli lemma again, it follows that with probability one, for $n$ large enough, $\Phi_{1}\left(w_{n}\right) / \varphi_{1}\left(w_{n}\right)$ is between $1-1 / m$ and $1+1 / m$. In addition, we have

$$
n \leq \varphi_{1}\left(w_{n}\right)=\pi\left(\mathcal{A}_{1} \cap\left\{w_{n}\right\}\right)+\lim _{w \downarrow w_{n}} \uparrow \varphi_{1}(w) \leq 1+n,
$$

by definition of $\mathcal{A}_{1}$ and $w_{n}$. Therefore, since $\Phi_{1}$ and $\varphi_{1}$ are nonincreasing, we infer that with probability one, for $n$ large enough and for $w \in\left[w_{n+1}, w_{n}\right]$,

$$
\frac{n}{n+2}\left(1-\frac{1}{m}\right) \leq \frac{\Phi_{1}(w)}{\varphi_{1}(w)} \leq \frac{n+2}{n}\left(1+\frac{1}{m}\right),
$$

and (36) with $\ell=1$ follows in a straightforward manner.

### 11.3 Integrability with respect to a Poisson measure

The third preliminary lemma is a direct consequence of the second one and deals with the integrability of a gauge function with respect to a Poisson random measure.
Lemma 6 Let $\Pi$ be a Poisson random measure on $(0,1]$ with intensity $\pi \in \mathcal{P}$. Then, with probability one, for any gauge function $g \in \mathfrak{D}_{1}$,

$$
\int_{r \in(0,1]} g(r) \Pi(\mathrm{d} r)=\infty \Longleftrightarrow \int_{r \in(0,1]} g(r) \pi(\mathrm{d} r)=\infty
$$

Proof Assume that the almost sure event on which the statement of Lemma 5 holds occurs, and consider $g \in \mathfrak{D}_{1}$. The gauge function $g$ is continuous in the neighborhood of the origin, but not necessarily on the whole interval [ 0,1$]$. However, we may assume that this is the case, because the integrability of $g$ with respect to $\Pi$ or $\pi$ depends on its behavior near zero only. Now, as $g$ is continuous and nondecreasing on $[0,1]$, we may consider its associated Lebesgue-Stieltjes measure $\gamma$, characterized by $\gamma([0, r])=g(r)$ for $r \in[0,1]$. Given that $\Pi([w, 1])$ is equivalent to $\pi([w, 1])$ as $w \rightarrow 0$, we have

$$
\int_{w \in(0,1]} \Pi([w, 1]) \gamma(\mathrm{d} w)=\infty \Longleftrightarrow \int_{w \in(0,1]} \pi([w, 1]) \gamma(\mathrm{d} w)=\infty
$$

We conclude by remarking that, owing to the Fubini-Tonelli theorem, the above integrals are equal to those appearing in the statement of the lemma.

### 11.4 Approximation by homogeneously distributed hyperplanes

The next lemma is a general result on Diophantine approximation by hyperplanes, under the assumption that the hyperplanes are homogeneously distributed, in a specific sense that we now introduce. Given $\mathfrak{s} \in \mathbb{S}^{d-1}$, let $\mathrm{Hyp}_{\mathfrak{s}}$ be the set of all hyperplanes $h$ represented by a pair $(\rho, s) \in \mathbb{H}_{d}$ with $\langle s, \mathfrak{s}\rangle \neq 0$. Such a hyperplane $h$ is not parallel to $\mathfrak{s}$. So, for any $t \in \mathbb{R}^{d}$, the line $t+\mathbb{R} \mathfrak{s}$ meets $h$ at a single point $t+\xi_{t}^{\mathfrak{s}}(h) \mathfrak{s}$, where

$$
\xi_{t}^{\mathfrak{s}}(h)=\frac{\rho-\langle s, t\rangle}{\langle s, \mathfrak{s}\rangle} .
$$

Now, let $\mathfrak{H}=\left(\mathfrak{H}_{n}\right)_{n \geq 1}$ be a sequence in $\operatorname{Hyp}_{\mathfrak{s}}$ and let $W$ be a nonempty open subset of $\mathbb{R}^{d}$. Given $t \in W, \delta \in(0,1)$ and $j \geq 0$, let us consider the first $\left\lfloor 2^{j} / \delta\right\rfloor$ hyperplanes $\mathfrak{H}_{n}$ and focus on those which intersect the line $t+\mathbb{R} \mathfrak{s}$ at a single point located at a distance less than $\delta$ from $t$ or, equivalently, those for which $\left|\xi_{t}^{\mathfrak{S}}\left(\mathfrak{H}_{n}\right)\right|<\delta$. The fact that the sequence $\mathfrak{H}$ is homogeneously distributed in $W$ basically means that the resulting intersection points are dispersed in a regular manner around almost every point $t \in W$ in the direction $\mathfrak{s}$, in the sense that the set

$$
\mathrm{Q}_{t, \delta, j}^{\mathfrak{S}}(\mathfrak{H})=\left\{q \in\left\{0, \ldots, 2^{j}-1\right\} \left\lvert\, q=\left\lfloor\frac{2^{j}}{\delta}\left|\xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)\right|\right\rfloor\right. \text { for some } n \leq \frac{2^{j}}{\delta}\right\}
$$

has cardinality of the order of $2^{j}$, asymptotically. Accordingly, we let $\operatorname{Hom}_{\mathfrak{s}}(W)$ be the set of sequences $\mathfrak{H}$ in $\operatorname{Hyp}_{\mathfrak{s}}$ such that for $\mathcal{L}^{d}$-almost every $t \in W$,

$$
\limsup _{\delta \rightarrow 0} \underline{Q}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})>0 \quad \text { with } \underline{\mathrm{Q}}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})=\liminf _{j \rightarrow \infty} \frac{1}{2^{j}} \# \mathrm{Q}_{t, \delta, j}^{\mathfrak{F}}(\mathfrak{H}),
$$

where \# stands for cardinality, and we say that such a sequence is homogeneously distributed in $W$. Here, $\mathcal{L}^{d}$ denotes the Lebesgue measure in $\mathbb{R}^{d}$.

Lemma 7 Let $\mathfrak{H}=\left(\mathfrak{H}_{n}\right)_{n \geq 1}$ be a sequence in $\operatorname{Hom}_{\mathfrak{s}}(W)$ and let $\mathfrak{R}=\left(\mathfrak{R}_{n}\right)_{n \geq 1}$ be a nonincreasing sequence of positive real numbers such that $\sum_{n} \mathfrak{R}_{n}=\infty$. Then, the set

$$
\mathfrak{F}(\mathfrak{H}, \mathfrak{R})=\left\{t \in \mathbb{R}^{d} \mid \mathrm{d}\left(t, \mathfrak{H}_{n}\right)<\mathfrak{R}_{n} \text { for i.m. } n \geq 1\right\}
$$

has full Lebesgue measure in $W$.

Proof For any $t \in \mathbb{R}^{d}$, let $\mathfrak{E}_{t}(\mathfrak{H}, \mathfrak{R})$ be the set of $\tau \in \mathbb{R}$ such that $\left|\tau-\xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)\right|<\mathfrak{R}_{n}$ for infinitely many $n \geq 1$. It suffices to prove that

$$
\begin{equation*}
\forall t \in \mathbb{R}^{d} \quad \forall \delta \in(0,1) \quad \mathcal{L}^{1}\left(\mathfrak{E}_{t}(\mathfrak{H}, \mathfrak{R}) \cap(-\delta, \delta)\right) \geq \frac{\delta}{18} \underline{Q}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})^{2} \tag{39}
\end{equation*}
$$

Indeed, letting $(\mathbb{R} \mathfrak{s})^{\perp}$ be the orthogonal complement of $\mathbb{R} \mathfrak{s}$ in $\mathbb{R}^{d}$, we deduce from (39) that for all $b \in(\mathbb{R} \mathfrak{s})^{\perp}$ and $\tau \in \mathbb{R}$,

$$
\forall \delta \in(0,1) \quad \frac{1}{2 \delta} \int_{\tau-\delta}^{\tau+\delta} \mathbb{1}_{\mathfrak{E}_{b}(\mathfrak{H}, \mathfrak{R})}(v) \mathrm{d} v \geq\left(\frac{\underline{\mathrm{Q}}_{b+\tau \mathfrak{s}, \delta}^{\mathfrak{s}}(\mathfrak{H})}{6}\right)^{2} .
$$

By virtue of Lebesgue's density theorem [32, Corollary 2.14], this implies that for every $b \in(\mathbb{R} \mathfrak{s})^{\perp}$ and $\mathcal{L}^{1}$-almost every $\tau \in \mathbb{R}$,

$$
\mathbb{1}_{\mathfrak{E}_{b}(\mathfrak{H}, \mathfrak{R})}(\tau) \geq\left(\frac{1}{6} \limsup _{\delta \rightarrow 0} \underline{\mathrm{Q}}_{b+\tau \mathfrak{s}, \delta}^{\mathfrak{s}}(\mathfrak{H})\right)^{2} .
$$

As the sequence $\mathfrak{H}$ is in $\operatorname{Hom}_{\mathfrak{s}}(W)$, we infer that for $\mathcal{H}^{d-1}$-almost every $b \in(\mathbb{R} \mathfrak{s})^{\perp}$ and $\mathcal{L}^{1}$-almost every $\tau \in \mathbb{R}$,

$$
b+\tau \mathfrak{s} \in W \Longrightarrow \mathbb{1}_{\mathfrak{E}_{b}(\mathfrak{H}, \mathfrak{R})}(\tau)>0 \Longrightarrow b+\tau \mathfrak{s} \in \mathfrak{F}(\mathfrak{H}, \mathfrak{R})
$$

and the result follows. It now remains to establish (39). To this end, we may clearly assume that $\underline{\mathrm{Q}}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})>0$, and then show that for every $\eta>1$,

$$
\begin{equation*}
\mathcal{L}^{1}\left(\mathfrak{E}_{t}(\mathfrak{H}, \mathfrak{R}) \cap(-\delta, \delta)\right) \geq \frac{\delta}{18}\left(\frac{\underline{\underline{Q}}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})}{\eta}\right)^{2} . \tag{40}
\end{equation*}
$$

In order to prove (40), let us observe that, by definition of $\underline{\mathrm{Q}}_{t, \delta}^{\mathfrak{F}}(\mathfrak{H})$, there is an integer $\underline{j} \geq 0$ such that for any $j \geq \underline{j}$, there exists a set $\mathcal{N}_{j} \subseteq\left\{1, \ldots,\left\lfloor 2^{j} / \delta\right\rfloor\right\}$ with:
$-2^{j-1} \underline{\mathrm{Q}}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H}) / \eta \leq \# \mathcal{N}_{j} \leq 2^{j}$;

- $\left|\xi_{t}^{\mathfrak{S}}\left(\mathfrak{H}_{n}\right)\right|<\delta$ for any $n \in \mathcal{N}_{j}$;

In addition, for each $n \geq 1$, let $\mathfrak{R}_{n}^{\prime}=\mathfrak{R}_{n} \wedge(2(n+1))^{-1}$. Then, $\left(\Re_{n}^{\prime}\right)_{n \geq 1}$ is nonincreasing and $\sum_{n} \Re_{n}^{\prime}=\infty$. Indeed, the convergence of $\sum_{n} \Re_{n}^{\prime}$ would imply that $\mathfrak{R}_{n}^{\prime}=\mathrm{o}(1 / n)$ as $n \rightarrow \infty$, thus contradicting the divergence of $\sum_{n} \Re_{n}$.

For every $j \geq \underline{j}$, let us consider the set

$$
U_{j}=\bigcup_{n \in \mathcal{N}_{j}}\left(\xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)-\mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}, \xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)+\mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}\right) .
$$

All the points in the lim sup of the sets $U_{j}$, except maybe those of the form $\xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)$, belong to both $[-\delta, \delta]$ and $\mathfrak{E}_{t}(\mathfrak{H}, \mathfrak{R})$. Therefore,

$$
\mathcal{L}^{1}\left(\limsup _{j \rightarrow \infty} U_{j}\right) \leq \mathcal{L}^{1}\left(\mathfrak{E}_{t}(\mathfrak{H}, \mathfrak{R}) \cap(-\delta, \delta)\right)
$$

Lemma 5 in [40] enables us to obtain an appropriate lower bound on the left-hand side above. To be more specific, this result ensures that

$$
\mathcal{L}^{1}\left(\limsup _{j \rightarrow \infty} U_{j}\right) \geq \limsup _{J \rightarrow \infty} \frac{\left(\sum_{\underline{j} \leq j \leq J} \mathcal{L}^{1}\left(U_{j}\right)\right)^{2}}{\sum_{\underline{j} \leq j, j^{\prime} \leq J} \mathcal{L}^{1}\left(U_{j} \cap U_{j^{\prime}}\right)}
$$

with the proviso that $\sum_{j} \mathcal{L}^{1}\left(U_{j}\right)$ diverges. To check this last condition, it is crucial to observe that for each $j \geq \underline{j}$, the intervals forming the set $U_{j}$ are disjoint, so that

$$
\begin{equation*}
\mathcal{L}^{1}\left(U_{j}\right)=2 \mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime} \# \mathcal{N}_{j} \geq \frac{\underline{\mathrm{Q}}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})}{\eta} 2^{j} \mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime} \tag{41}
\end{equation*}
$$

Due to the divergence of $\sum_{n} \mathfrak{R}_{n}^{\prime}$ and the Cauchy condensation test, $\sum_{j} 2^{j} \mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}$ diverges too. The divergence of $\sum_{j} \mathcal{L}^{1}\left(U_{j}\right)$ then follows from (41). As a consequence, in order to obtain (40), it now suffices to show that for $J$ large enough,

$$
\begin{equation*}
\sum_{\underline{j} \leq j, j^{\prime} \leq J} \mathcal{L}^{1}\left(U_{j} \cap U_{j^{\prime}}\right) \leq \frac{18}{\delta}\left(\frac{\eta}{\underline{\mathrm{Q}}_{t, \delta}^{\mathfrak{s}}(\mathfrak{H})} \sum_{\underline{j} \leq j \leq J} \mathcal{L}^{1}\left(U_{j}\right)\right)^{2} \tag{42}
\end{equation*}
$$

To this end, let us derive an upper bound on the Lebesgue measure of $U_{j} \cap U_{j^{\prime}}$, where $\underline{j} \leq j<j^{\prime}$. This set is the union over $n \in \mathcal{N}_{j}$ and $n^{\prime} \in \mathcal{N}_{j^{\prime}, n}$ of the sets
$\left(\xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)-\mathfrak{R}_{\left\lfloor 2{ }^{j} / \delta\right\rfloor}^{\prime}, \xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n}\right)+\mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}\right) \cap\left(\xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n^{\prime}}\right)-\mathfrak{R}_{\left\lfloor 2^{j^{\prime}} / \delta\right\rfloor}^{\prime}, \xi_{t}^{\mathfrak{s}}\left(\mathfrak{H}_{n^{\prime}}\right)+\mathfrak{R}_{\left\lfloor 2^{\left.j^{\prime} / \delta\right\rfloor}\right.}^{\prime}\right)$,
where $\mathcal{N}_{j^{\prime}, n}$ is the collection of all $n^{\prime} \in \mathcal{N}_{j^{\prime}}$ such that this last intersection is nonempty. The cardinality of $\mathcal{N}_{j^{\prime}, n}$ is clearly bounded above by the number of integers of the form
 Along with the fact that $\# \mathcal{N}_{j} \leq 2^{j}$, this yields

$$
\mathcal{L}^{1}\left(U_{j} \cap U_{j^{\prime}}\right) \leq 2^{j}\left(2+\frac{2^{j^{\prime}+2}}{\delta} \mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}\right)\left(2 \mathfrak{R}_{\left\lfloor 2^{\left.j^{\prime} / \delta\right\rfloor}\right.}^{\prime}\right) .
$$

As a consequence, for any integer $J>\underline{j}$, the left-hand side of (42) is at most

$$
2 \sum_{\underline{j} \leq j \leq J} 2^{j} \Re_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}+8 \sum_{\underline{j} \leq j<j^{\prime} \leq J} 2^{j} \Re_{\left\lfloor 2^{\left.j^{\prime} / \delta\right\rfloor}\right.}^{\prime}+\frac{16}{\delta} \sum_{\underline{j} \leq j<j^{\prime} \leq J} 2^{j+j^{\prime}} \mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime} \Re_{\left\lfloor 2 j^{\prime} / \delta\right\rfloor}^{\prime}
$$

The third sum is smaller than half the sum over all $j, j^{\prime} \in\{j, \ldots, J\}$, and the second sum is smaller than the first one. Thus, the left-hand side of (42) is at most

$$
10 \sum_{\underline{j} \leq j \leq J} 2^{j} \mathfrak{R}_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}+\frac{8}{\delta}\left(\sum_{\underline{j} \leq j \leq J} 2^{j} \Re_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}\right)^{2} \leq \frac{18}{\delta}\left(\sum_{\underline{j} \leq j \leq J} 2^{j} \Re_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}\right)^{2}
$$

where the last bound holds for $J$ large enough, due to the divergence of the series $\sum_{j} 2^{j} \Re_{\left\lfloor 2^{j} / \delta\right\rfloor}^{\prime}$. We conclude using (41).
11.5 Poisson hyperplanes are homogeneously distributed

The next lemma asserts essentially that the random hyperplanes arising in the definition of $K_{\nu}(\alpha)$ are homogeneously distributed. To be more specific, we need to introduce some additional notations. Given $\alpha>0$, let $\mathrm{N}_{\alpha}$ denote the image under the mapping $(\rho, s, x) \mapsto\left(\rho, s,|x|^{1 / \alpha}\right)$ of the restriction to $(0, \infty) \times \mathbb{S}^{d-1} \times([-1,1] \backslash\{0\})$ of the measure $\mathbf{N}$. Then, $\mathbf{N}_{\alpha}$ is a Poisson random measure with intensity $\mathcal{L}_{+}^{1} \otimes \nu_{\alpha}$, where $\nu_{\alpha}$ is defined at the beginning of Sect. 10.

Given an arbitrary orthonormal basis $\left(e_{1}, \ldots, e_{d}\right)$ of $\mathbb{R}^{d}$, for each $i \in\{1, \ldots, d\}$, let $\mathcal{C}_{i}$ be the set of all $s \in \mathbb{S}^{d-1}$ satisfying $\left|\left\langle s, e_{i}\right\rangle\right| \geq d^{-1 / 2}$. Furthermore, let $\mathcal{A}_{i, 0}$ be the set of all $r \in(0,1]$ such that $v_{\alpha}\left(\mathcal{C}_{i} \times\{r\}\right) \geq 1$, and let $\mathcal{A}_{i, 1}=(0,1] \backslash \mathcal{A}_{i, 0}$. Then, for each $\ell \in\{0,1\}$, let $v_{\alpha, i, \ell}$ be the restriction of $v_{\alpha}$ to $\mathcal{C}_{i} \times \mathcal{A}_{i, \ell}$.

Let us assume that $v_{\alpha, i, \ell}$ has infinite total mass. Thus, $\nu_{\alpha, i, \ell}$ belongs to $\mathcal{P}$, in view of (1). Given an integer $A \geq 1$, the restriction $\mathrm{N}_{\alpha}^{A, i, \ell}$ of $\mathrm{N}_{\alpha}$ to $(0, A) \times \mathcal{C}_{i} \times \mathcal{A}_{i, \ell}$ may be written almost surely as

$$
\begin{equation*}
\mathrm{N}_{\alpha}^{A, i, \ell}=\sum_{n=1}^{\infty} \delta_{\left(\mathrm{P}_{n}^{A, i, \ell}, S_{n}^{A, i, \ell}, R_{n}^{A, i, \ell}\right)} \tag{43}
\end{equation*}
$$

for some sequence $\left(\mathrm{P}_{n}^{A, i, \ell}, S_{n}^{A, i, \ell}, R_{n}^{A, i, \ell}\right)_{n \geq 1}$ in $(0, A) \times \mathcal{C}_{i} \times \mathcal{A}_{i, \ell}$. Given that $v_{\alpha, i, \ell}$ is in $\mathcal{P}$, the following holds with probability one: for any $\varepsilon>0$, only finitely many $n \geq 1$ satisfy $R_{n}^{A, i, \ell} \geq \varepsilon$. So, up to a reordering, we may assume that the sequence $\left(R_{n}^{A}, i, \ell\right)_{n \geq 1}$ is nonincreasing and converges to zero. Last, for any $n \geq 1$, let $H_{n}^{A, i, \ell}$ denote the hyperplane defined in terms of $\mathrm{P}_{n}^{A, i, \ell}$ and $S_{n}^{A, i, \ell}$ as in (9).

Lemma 8 Almost surely, for any $A \geq 1, i \in\{1, \ldots, d\}$ and $\ell \in\{0,1\}$ such that $v_{\alpha, i, \ell}$ has infinite total mass, $H^{A, i, \ell}=\left(H_{n}^{A}, i, \ell\right)_{n \geq 1}$ is in $\operatorname{Hom}_{e_{i}}\left(\mathrm{~B}_{0}(A)\right)$.

Proof Given $A \geq 1, i \in\{1, \ldots, d\}$ and $\ell \in\{0,1\}$, it is clear that $H_{n}^{A, i, \ell} \in \operatorname{Hyp}_{e_{i}}$ for all $n \geq 1$. Assume that $v_{\alpha, i, \ell}$ has infinite total mass. In view of Fubini's theorem, it suffices to let $t \in \mathrm{~B}_{0}(A)$ and show that almost surely, $\lim \sup _{\delta \rightarrow 0} \underline{\mathrm{Q}}_{t, \delta}^{e_{i}}\left(H^{A, i, \ell}\right)>0$. To this end, given $\delta \in(0,1)$ and $j \geq 1$, suppose that $\# \mathrm{Q}_{t, \delta, j}^{e_{i}}\left(H^{A, i, \ell}\right) \leq \theta 2^{j}$, for some $\theta \in(0,1)$. So, there is a subset Q of $\left\{0, \ldots, 2^{j}-1\right\}$ with cardinality $\left\lfloor\theta 2^{j}\right\rfloor$ which satisfies $\mathrm{Q} \supseteq \mathrm{Q}_{t, \delta, j}^{e_{i}}\left(H^{A, i, \ell}\right)$. We have

$$
\begin{equation*}
\forall n \leq\left\lfloor\frac{2^{j}}{\delta}\right\rfloor\left|\xi_{t}^{e_{i}}\left(H_{n}^{A, i, \ell}\right)\right|=\left|\frac{\mathrm{P}_{n}^{A, i, \ell}-\left\langle S_{n}^{A, i, \ell}, t\right\rangle}{\left\langle S_{n}^{A, i, \ell}, e_{i}\right\rangle}\right| \in[\delta, \infty] \cup \bigcup_{q \in \mathrm{Q}} \lambda_{q}, \tag{44}
\end{equation*}
$$

where $\lambda_{q}$ denotes the interval $\left[q \delta 2^{-j},(q+1) \delta 2^{-j}\right)$.
Let us first assume that $\ell=0$ and derive an upper bound on the probability that (44) happens. There exists a decreasing sequence $\left(a_{p}^{i}\right)_{p \geq 1}$ of positive reals that converges to zero and whose terms form the set $\mathcal{A}_{i, 0}$. Then, basic properties of Poisson random measures enable us to write that almost surely,

$$
\mathrm{N}_{\alpha}^{A, i, 0}=\sum_{p=1}^{\infty} \sum_{m=1}^{N_{p}} \delta_{\left(\widetilde{\mathrm{P}}_{p, m}^{A, i, 0}, \widetilde{S}_{p, m}^{A, i, 0}, a_{p}^{i}\right)},
$$

where each $N_{p}$ is Poisson distributed with mean $A v_{\alpha, i, 0}\left(\mathcal{C}_{i} \times\left\{a_{p}^{i}\right\}\right)$, each $\widetilde{\mathrm{P}}_{p, m}^{A, i, 0}$ is uniformly distributed on $(0, A)$ and each $\widetilde{S}_{p, m}^{A, i, 0}$ is distributed according to the probability measure $\mu_{i, a_{p}^{i}}=v_{\alpha, i, 0}\left(\cdot \times\left\{a_{p}^{i}\right\}\right) / v_{\alpha, i, 0}\left(\mathcal{C}_{i} \times\left\{a_{p}^{i}\right\}\right)$ on $\mathcal{C}_{i}$, and all these variables are independent, see e.g. [34]. This means that, in (43) above, we have the property that, conditional on the $\sigma$-algebra generated by the sequence $\left(R_{n}^{A, i, 0}\right)_{n \geq 1}$, the variables $\mathrm{P}_{n}^{A, i, 0}$ and $S_{n}^{A, i, 0}$, for $n \geq 1$, are independent and distributed according to the uniform law on $(0, A)$ and the law $\mu_{i, R_{n}^{A, i, 0},}$, respectively. Hence, the conditional probability that (44) holds given that $\sigma$-algebra is

$$
\left.\prod_{n=1}^{\substack{\left.2^{j} / \delta\right\rfloor}} \int_{\substack{\rho \in(0, A) \\ s \in \mathcal{C}_{i}}} \mathbb{1}_{\left\{\left|\frac{\rho-\langle s, t\rangle}{\left\langle s, e_{i}\right\rangle}\right| \in[\delta, \infty] \cup\right.}^{\cup} \bigcup_{q \in \mathrm{Q}} \lambda_{q}\right\},
$$

Exploiting the symmetry of $v$, it is easy to check that this is equal to

$$
\prod_{n=1}^{\left\lfloor 2^{j} / \delta\right\rfloor}\left(1-\frac{\delta-\left\lfloor\theta 2^{j}\right\rfloor \delta 2^{-j}}{A} \int_{s \in \mathcal{C}_{i}}\left|\left\langle s, e_{i}\right\rangle\right| \mu_{i, R_{n}^{A, i, 0}}(\mathrm{~d} s)\right) \leq \exp \left(-\frac{(1-\theta) \delta}{A \sqrt{d}}\left\lfloor\frac{2^{j}}{\delta}\right\rfloor\right)
$$

for $\delta<A-\|t\|$. In order to derive the upper bound, we have used the fact that $1-z \leq \mathrm{e}^{-z}$ for any $z \in \mathbb{R}$ and that $\left|\left\langle s, e_{i}\right\rangle\right| \geq d^{-1 / 2}$ for any $s \in \mathcal{C}_{i}$. It follows that

$$
\mathbb{P}\left(\forall n \leq\left\lfloor\frac{2^{j}}{\delta}\right\rfloor\left|\xi_{t}^{e_{i}}\left(H_{n}^{A, i, 0}\right)\right| \in[\delta, \infty] \cup \bigcup_{q \in \mathrm{Q}} \lambda_{q}\right) \leq \exp \left(-\frac{1-\theta}{A \sqrt{d}} 2^{j-1}\right)
$$

Now, recall that the set Q is a subset of $\left\{0, \ldots, 2^{j}-1\right\}$ with cardinality $\left\lfloor\theta 2^{j}\right\rfloor$. We finally infer that

$$
\mathbb{P}\left(\# \mathrm{Q}_{t, \delta, j}^{e_{i}}\left(H^{A, i, 0}\right) \leq \theta 2^{j}\right) \leq\binom{ 2^{j}}{\left\lfloor\theta 2^{j}\right\rfloor} \exp \left(-\frac{1-\theta}{A \sqrt{d}} 2^{j-1}\right) .
$$

Making use of Stirling's formula, we get

$$
\begin{equation*}
\limsup _{j \rightarrow \infty} \frac{1}{2^{j}} \log \mathbb{P}\left(\# \mathrm{Q}_{t, \delta, j}^{e_{i}}\left(H^{A, i, 0}\right) \leq \theta 2^{j}\right) \leq-\Gamma(\theta)-\frac{1-\theta}{2 A \sqrt{d}}, \tag{45}
\end{equation*}
$$

with $\Gamma(\theta)=\theta \log \theta+(1-\theta) \log (1-\theta)$. Clearly, there exists a unique $\theta_{0} \in(0,1)$ at which the right-hand side of (45) vanishes. Moreover, this right-hand side is negative for any $\theta \in\left(0, \theta_{0}\right)$. Using the Borel-Cantelli lemma, we see that for any $\delta \in$ $(0, A-\|t\|)$ and any such $\theta$, almost surely, $\underline{Q}_{t, \delta}^{e_{i}}\left(H^{A, i, 0}\right) \geq \theta$. The result follows.

Now, if $\ell=1$, we may apply Lemma 5 with the image under $(\rho, s, r) \mapsto r$ of the measure $\mathcal{L}_{\mid(0, A)}^{1} \otimes v_{\alpha, i, 1}$, where $\mathcal{L}_{\mid(0, A)}^{1}$ is the restriction of $\mathcal{L}^{1}$ to $(0, A)$. Consequently, we infer that $\Phi_{i, 1}^{A} \sim_{0} A \varphi_{i, 1}$ almost surely, where

$$
\forall w \in(0,1] \quad\left\{\begin{array}{l}
\Phi_{i, 1}^{A}(w)=\mathrm{N}_{\alpha}^{A, i, 1}\left((0, A) \times \mathcal{C}_{i} \times[w, 1]\right) \\
\varphi_{i, 1}(w)=v_{\alpha, i, 1}\left(\mathcal{C}_{i} \times[w, 1]\right)
\end{array}\right.
$$

This directly implies that with probability one,

$$
C=\sup _{w \in(0,1]} \frac{\Phi_{i, 1}^{A}(w)}{A\left(1+\varphi_{i, 1}(w)\right)} \in[1, \infty)
$$

Furthermore, let us consider an integer $k \geq 1$ and an integer $j$ large enough to ensure that $1+\varphi_{i, 1}(w) \leq 2^{j} /(A \delta k)$ for some $w \in(0,1]$. As $\left(R_{n}^{A, i, 1}\right)_{n \geq 1}$ is nonincreasing, we have $n \leq \Phi_{i, 1}^{A}(w) \leq C 2^{j} /(\delta k)$ for any $n \geq 1$ with $R_{n}^{A, i, 1} \geq w$. Therefore, if $C \leq k$ and (44) hold simultaneously, then no integer $n \geq 1$ can satisfy both $R_{n}^{A, i, 1} \geq w$ and $\left|\xi_{t}^{e_{i}}\left(H_{n}^{A, i, 1}\right)\right| \in[0, \delta) \backslash \bigcup_{q \in \mathrm{Q}} \lambda_{q}$. This happens with probability $\mathrm{e}^{-I(w)}$, where

$$
I(w)=\int_{\substack{(\rho, s) \in(0, A) \times \mathcal{C}_{i} \\ r \in[w, 1]}} \mathbb{1}_{\left\{\left|\frac{\rho-\langle s, t\rangle}{\left\langle s, e_{i}\right\rangle}\right| \in[0, \delta) \backslash \bigcup_{\substack{q \in \mathbb{Q}}} \lambda_{q}\right\}} \mathrm{d} \rho v_{\alpha, i, 1}(\mathrm{~d} s, \mathrm{~d} r) \geq \frac{(1-\theta) \delta}{\sqrt{d}} \varphi_{i, 1}(w) .
$$

Here, the lower bound holds for $\delta<A-\|t\|$ and results from standard computations that exploit the symmetry of $v$. This leads to an upper bound on the probability that $C \leq k$ and (44) hold simultaneously, which may be optimized by letting $w \downarrow w_{*}$,
where $w_{*}$ is the infimum of all $w \in(0,1]$ with $1+\varphi_{i, 1}(w) \leq 2^{j} /(A \delta k)$. By definition of $w_{*}$ and $\mathcal{A}_{i, 1}$, and as $\varphi_{i, 1}$ is left-continuous with right limits, we have

$$
\lim _{w \downarrow w_{*}} \varphi_{i, 1}(w) \geq \varphi_{i, 1}\left(w^{*}\right)-1 \geq \frac{2^{j}}{A \delta k}-2 \geq \frac{2^{j-1}}{A \delta k}
$$

where the last inequality holds for $j$ large enough. We deduce that
$\mathbb{P}\left(C \leq k\right.$ and $\left.\forall n \leq\left\lfloor\frac{2^{j}}{\delta}\right\rfloor\left|\xi_{t}^{e_{i}}\left(H_{n}^{A, i, 1}\right)\right| \in[\delta, \infty] \cup \bigcup_{q \in \mathrm{Q}} \lambda_{q}\right) \leq \exp \left(-\frac{1-\theta}{A k \sqrt{d}} 2^{j-1}\right)$.
Just as in the previous case, this implies that for any $k \geq 1$ with $\mathbb{P}(C \leq k)>0$,

$$
\limsup _{j \rightarrow \infty} \frac{1}{2^{j}} \log \mathbb{P}\left(\# \mathrm{Q}_{t, \delta, j}^{e_{i}}\left(H^{A, i, 1}\right) \leq \theta 2^{j} \mid C \leq k\right) \leq-\Gamma(\theta)-\frac{1-\theta}{2 A k \sqrt{d}},
$$

where the right-hand side vanishes for a unique $\theta_{k} \in(0,1)$. Using the Borel-Cantelli lemma as above, we see that for any $\delta \in(0, A-\|t\|)$ and any $\theta \in\left(0, \theta_{k}\right)$, conditionally on $C \leq k$, with probability one, $\underline{\mathrm{Q}}_{t, \delta}^{e_{i}}\left(H^{A, i, 1}\right) \geq \theta$. Therefore, with probability one, conditionally on $C \leq k$, we have lim sup $\sin _{\delta \rightarrow 0} \underline{\mathrm{Q}}_{t, \delta}^{e_{i}}\left(H^{A, i, 1}\right) \geq \theta_{k}$. The result now follows from the fact that $C<\infty$ almost surely.

### 11.6 Ubiquity

Last, the proof of Theorem 7 calls upon the next lemma, which is a straightforward consequence of Theorem 3.6 in [17]. (The hypotheses of that theorem are verified here because the diameter of the set of $t \in \mathbb{R} e_{i}$ such that $\mathrm{d}\left(t, H_{n}^{A, i, \ell}\right)<1$ is bounded above by $2 d^{1 / 2}$, since $S_{n}^{A, i, \ell} \in \mathcal{C}_{i}$.) For every $g \in \mathfrak{D}_{1}$, let

$$
F^{A, i, \ell}(g)=\left\{t \in \mathbb{R}^{d} \mid \mathrm{d}\left(t, H_{n}^{A, i, \ell}\right)<g\left(R_{n}^{A, i, \ell}\right) \text { for i.m. } n \geq 1\right\}
$$

Lemma 9 Let $A \geq 1, i \in\{1, \ldots, d\}$ and $\ell \in\{0,1\}$ such that $v_{\alpha, i, \ell}$ has infinite total mass, and let $g \in \mathfrak{D}_{1}$. If the set $F^{A, i, \ell}(\underline{g})$ has full Lebesgue measure in $\mathrm{B}_{0}(A)$ for some $\underline{g} \in \mathfrak{D}_{1}$ with $g \prec \underline{g}$, then $F^{A, i, \ell}(r \mapsto r) \in \mathrm{G}^{d-1, g}\left(\mathrm{~B}_{0}(A)\right)$.

This lemma falls into the category of the ubiquity results obtained in $[16,17,19]$, which enable one to deduce the large intersection properties of a set, such as $F^{A, i, \ell}(r \mapsto r)$ in the present situation, from the sole knowledge of the Lebesgue measure of a corresponding enlarged set, which is $F^{A, i, \ell}(\underline{g})$ here.
11.7 End of the proof when $h_{v}(g)>\alpha$

To begin with, let us recall that $\mathrm{P}_{n}$ and $X_{n}$ are defined in terms of the atoms of the Poisson random measure N , see Sect. 4. Then, let $R_{n}=\left|X_{n}\right|^{1 / \alpha}$ for any $n \geq 1$ and,
given an integer $A \geq 1$, let $\mathcal{N}_{A}$ denote the set of all $n \geq 1$ such that $\mathrm{P}_{n}<A$ and $R_{n} \leq 1$.

For any $A \geq 1$, we may apply Lemma 6 with the image under $(\rho, s, r) \mapsto r$ of the measure $\mathcal{L}_{\mid(0, A)}^{1} \otimes v_{\alpha}$. Consequently, in view of (19), we deduce that with probability one, for any $A \geq 1$ and any $g \in \mathfrak{D}_{1}$ with $h_{v}(g)>\alpha$, the series $\sum_{n \in \mathcal{N}_{A}} g\left(R_{n}\right)$ converges. In particular, for any $\delta>0$, there necessarily exists an integer $n_{0} \geq 1$ such that $R_{n}<\delta$ for any $n \in \mathcal{N}_{A}$ with $n \geq n_{0}$. Then, for any $n_{1} \geq n_{0}$,

$$
K_{\nu}(\alpha) \cap \mathrm{B}_{0}(A-1) \subseteq \bigcup_{\substack{n \in \mathcal{N}_{A} \\ n \geq n_{1}}}\left\{t \in \mathrm{~B}_{0}(A) \mid \mathrm{d}\left(t, H_{n}\right)<R_{n}\right\}
$$

where $H_{n}$ is the hyperplane defined by (9). Moreover, each set in the union above may be covered by $\left(3\left\lfloor 2 A \sqrt{d} / R_{n}\right\rfloor\right)^{d-1}$ open balls with radius $2 R_{n}$. Therefore,

$$
\begin{aligned}
\mathcal{H}_{\delta}^{r \mapsto r^{d-1} g(r)}\left(K_{v}(\alpha) \cap \mathrm{B}_{0}(A-1)\right) & \leq \sum_{\substack{n \in \mathcal{N}_{A} \\
n \geq n_{1}}}\left(3\left\lfloor\frac{2 A \sqrt{d}}{R_{n}}\right\rfloor\right)^{d-1}\left(4 R_{n}\right)^{d-1} g\left(4 R_{n}\right) \\
& \leq 4(24 A \sqrt{d})^{d-1} \sum_{\substack{n \in \mathcal{N}_{A} \\
n \geq n_{1}}} g\left(R_{n}\right) .
\end{aligned}
$$

Letting $n_{1} \rightarrow \infty$ and $\delta \rightarrow 0$, we deduce that $\mathcal{H}^{d-1, g}\left(K_{\nu}(\alpha) \cap \mathrm{B}_{0}(A-1)\right)=0$. This holds for all integers $A \geq 1$, so the result follows.

### 11.8 End of the proof when $h_{\nu}(g)<\alpha$

By Lemma 6, the next statement holds almost surely: for all integers $A \geq 1, i \in$ $\{1, \ldots, d\}$ and $\ell \in\{0,1\}$, and all $g \in \mathfrak{D}_{1}$,

$$
\begin{equation*}
\int_{\substack{s \in \mathbb{S}^{d-1} \\ r \in(0,1]}} g(r) v_{\alpha, i, \ell}(\mathrm{~d} s, \mathrm{~d} r)=\infty \Longrightarrow \sum_{n=1}^{\infty} g\left(R_{n}^{A, i, \ell}\right)=\infty . \tag{46}
\end{equation*}
$$

Moreover, the statement of Lemma 8 holds with probability one as well. We shall work in what follows on the almost sure event on which these two statements hold. Let us consider a gauge function $g \in \mathfrak{D}_{1}$ such that $h_{\nu}(g)<\alpha$. Due to (19) and the fact that $\mathbb{S}^{d-1} \times(0,1]$ is covered by the sets $\mathcal{C}_{i} \times \mathcal{A}_{i, \ell}$, the integral in (46) is infinite for some $(i, \ell) \in\{1, \ldots, d\} \times\{0,1\}$.

Let us assume that $g \prec(r \mapsto r)$. Borrowing ideas from the proof of [18, Proposition 5], we may build $\underline{g} \in \mathfrak{D}_{1}$ such that $g \prec \underline{g}$ and the integral in (46) with $g$ replaced by $\underline{g}$ is infinite as well. Therefore, $\sum_{n} \underline{g}\left(R_{n}^{A, i, \ell}\right)$ diverges, so that $F^{A, i, \ell}(\underline{g})$ has full Lebesgue measure in $\mathrm{B}_{0}(A)$, by virtue of Lemma 7. Thanks to Lemma 9, we deduce that $F^{A, i, \ell}(r \mapsto r) \in \mathrm{G}^{d-1, g}\left(\mathrm{~B}_{0}(A)\right)$. The same result holds if $g \nprec(r \mapsto r)$, in view of the divergence of $\sum_{n=1}^{\infty} R_{n}^{A, i, \ell}$, combined with Lemma 7 and Theorem 5(4).

To conclude, it is crucial to note that the set $F^{A, i, \ell}(r \mapsto r)$ is contained in $K_{\nu}(\alpha)$. Given that $F^{A, i, \ell}(r \mapsto r)$ is in $\mathrm{G}^{d-1, g}\left(\mathrm{~B}_{0}(A)\right)$, the $G_{\delta}$-set $K_{v}(\alpha)$ belongs to the same class. Now, let $W$ be a bounded open subset of $\mathbb{R}^{d}$. Then, $W \subseteq \mathrm{~B}_{0}(A)$ for A large enough, so that $\mathcal{M}_{\infty}^{f}\left(K_{\nu}(\alpha) \cap W\right)=\mathcal{M}_{\infty}^{f}(W)$ for every $f \in \mathfrak{D}_{d}$ with $f \prec\left(r \mapsto r^{d-1} g(r)\right)$. Lemma 12 in [16] finally ensures that $K_{\nu}(\alpha) \in \mathrm{G}^{d-1, g}\left(\mathbb{R}^{d}\right)$.

## 12 Proofs concerning the jump component

Throughout the section, we assume that the measure $v$ is admissible. We now establish Proposition 1 and the results of Sect. 6.

### 12.1 Proof of Proposition 1

Let $A \in \mathbb{N}$ and $U_{A}=\overline{\mathrm{B}}_{0}(A+1) \cap\left(d^{-1 / 2} \mathbb{Z}^{d}\right)$. For any fixed $t \in \mathbb{R}^{d}$, the series $\sum_{j \geq 0} L_{\nu, j}(t)$ defining $L_{v}(t)$ converges almost surely, see Sect. 4. Thus, the event $\mathcal{E}_{A}$ consisting in the fact that the series $\sum_{j \geq 0} L_{v, j}(u)$, for $u \in U_{A}$, converge simultaneously has probability one. Furthermore, it follows from Lemma 2 that the event $\mathcal{E}_{A}^{\prime}=\left\{Z_{v}(A+1)<\infty\right\}$ has probability one too.

Let us now assume that the almost sure event $\mathcal{E}_{A} \cap \mathcal{E}_{A}^{\prime}$ happens, and let $\varepsilon>0$ and $t \in \overline{\mathrm{~B}}_{0}(A)$. Then, $\|t-u\| \leq 1 / 2$ for some $u \in U_{A}$. The sum $\chi_{v}$ defined by (3) is finite and the series $\sum_{j \geq 0} L_{v, j}(u)$ converges, so that

$$
\sum_{j=j_{1}}^{j_{2}} 2^{-j}\left(j+\left(j v_{j}\right)^{1 / 2}\right) \leq \varepsilon \quad \text { and } \quad\left|\sum_{j=j_{1}}^{j_{2}} L_{v, j}(u)\right| \leq \varepsilon
$$

for all integers $j_{1}$ and $j_{2}$ such that $j_{2} \geq j_{1} \geq j_{0}$, and some $j_{0} \geq 1$. Therefore,

$$
\left|\sum_{j=j_{1}}^{j_{2}} L_{v, j}(t)\right| \leq\left|\sum_{j=j_{1}}^{j_{2}} L_{v, j}(u)\right|+\sum_{j=j_{1}}^{j_{2}} \zeta_{v}(A+1, j, 1) \leq\left(1+\mathrm{Z}_{v}(A+1)\right) \varepsilon .
$$

The partial sums of $\sum_{j \geq 0} L_{v, j}(t)$ form a Cauchy sequence. So, for any $A \in \mathbb{N}$, with probability one, $\sum_{j \geq 0} L_{v, j}(t)$ converges for any $t \in \overline{\mathrm{~B}}_{0}(A)$, and the result follows.

### 12.2 Proof of Theorems 4 and 6

The statement of Corollary 5 holds with probability one, and that of Theorem 7 holds with probability one for all rational numbers $\alpha \in \mathbb{Q} \cap(0, \infty)$ simultaneously. Let us assume that the almost sure event on which these statements hold occurs. Theorems 4 and 6 follow from a series of propositions that we now state and establish. Throughout, we consider a real $h \in\left[0,1 / \beta_{\nu}\right)$, a gauge function $g \in \mathfrak{D}_{1}$ and a nonempty open set $W \subseteq \mathbb{R}^{d}$.

Proposition 11 If $h<h_{v}(g)$, then $\mathcal{H}^{d-1, g}\left(E_{v}(h)\right)=\mathcal{H}^{d-1, g}\left(E_{v}^{\prime}(h)\right)=0$.
Proof There exists a rational number $\alpha \in \mathbb{Q} \cap\left(h, h_{\nu}(g)\right)$ such that $E_{v}^{\prime}(h) \subseteq K_{\nu}(\alpha)$. Hence, $\mathcal{H}^{d-1, g}\left(E_{v}^{\prime}(h)\right) \leq \mathcal{H}^{d-1, g}\left(K_{v}(\alpha)\right)=0$. Furthermore, if $h>0$, then $E_{v}(h) \subseteq$ $E_{v}^{\prime}(h)$, so $\mathcal{H}^{d-1, g}\left(E_{v}(h)\right)=0$ as well. This result still holds for $h=0$, because $E_{v}(0)=E_{v}^{\prime}(0) \cup J_{v}$ and $J_{v}$ is a countable union of hyperplanes.

Proposition 12 If $h<h_{v}(g)$, then $E_{v}^{\prime}(h) \notin \mathrm{G}^{d-1, g}(W)$.
Proof Adapting the method developed in the proof of [18, Proposition 3], we may build a gauge function $\bar{g} \in \mathfrak{D}_{1}$ satisfying both $\bar{g} \prec g$ and $h_{v}(\bar{g}) \geq h_{v}(g)$. Then, applying Proposition 11 with $\bar{g}$ instead of $g$, we infer that $\mathcal{H}^{d-1, \bar{g}}\left(E_{v}^{\prime}(h) \cap W\right)=0$. We conclude by Theorem 5(3).

Proposition 13 If $h \geq h_{\nu}(g)$, then $E_{\nu}^{\prime}(h) \in \mathrm{G}^{d-1, g}(W)$.
Proof The mapping $\alpha \mapsto K_{\nu}(\alpha)$ is nondecreasing; it follows that

$$
E_{\nu}^{\prime}(h)=\left(\mathbb{R}^{d} \backslash J_{\nu}\right) \cap \bigcap_{\substack{h<\alpha \leq 1 / \beta_{\nu} \\ \alpha \in \mathbb{Q}^{2}}} K_{\nu}(\alpha)
$$

Furthermore, each set $K_{\nu}(\alpha)$ arising in this last intersection belongs to $\mathrm{G}^{d-1, g}(W)$. The set $\mathbb{R}^{d} \backslash J_{v}$ belongs to this class as well by virtue of Theorem 5 (4), because $J_{v}$ is the union of countably many hyperplanes. We conclude using Theorem 5 (1).

Proposition 14 If $h \geq h_{v}(g)$, then $\mathcal{H}^{d-1, g}\left(E_{v}^{\prime}(h) \cap W\right)=\infty$.
Proof The assumption of the proposition implies that $g \prec(r \mapsto r)$. Indeed, otherwise, we would clearly have $h_{v}(g) \geq 1 / \beta_{v}$, which is in contradiction with the fact that $h_{\nu}(g) \leq h<1 / \beta_{\nu}$. Therefore, borrowing ideas from the proof of [18, Proposition 5], we may build a gauge function $\underline{g} \in \mathfrak{D}_{1}$ satisfying both $g \prec \underline{g}$ and $h_{v}(\underline{g}) \leq h_{v}(g)$. Then, applying Proposition 13 with the gauge function $\underline{g}$ instead of $g$, we see that $E_{v}^{\prime}(h)$ belongs to $\mathrm{G}^{d-1, \underline{g}}(W)$. We conclude by Theorem 5 (3).

Proposition 15 If $h=h_{v}(g)$, then $\mathcal{H}^{d-1, g}\left(E_{v}(h) \cap W\right)=\infty$.
Proof In the case where $h=0$, the result follows directly from Proposition 14, because $E_{\nu}(0)$ contains $E_{v}^{\prime}(0)$. In the case where $h>0$, it suffices to make use of Proposition 14 again, together with the observation that

$$
E_{v}(h)=E_{v}^{\prime}(h) \backslash \bigcup_{\substack{0<\alpha<h \\ \alpha \in \mathbb{Q}}} K_{v}(\alpha),
$$

because the mapping $\alpha \mapsto K_{\nu}(\alpha)$ is nondecreasing, and that each set $K_{\nu}(\alpha)$ in the union above has Hausdorff measure zero for the gauge function $r \mapsto r^{d-1} g(r)$.

### 12.3 Proof of Corollary 2

Let us assume that the almost sure event on which the statement of Theorem 4 holds occurs. Let $h \in\left[0,1 / \beta_{v}\right), g \in \mathfrak{D}_{1}^{*}$ and $W$ be a nonempty open set. Most of the first part of the corollary, which gives $\mathcal{H}^{d-1, g}\left(E_{\nu}(h) \cap W\right)$ and $\mathcal{H}^{d-1, g}\left(E_{\nu}^{\prime}(h) \cap W\right)$, follows directly from Theorem 4 . The only new property is that $\mathcal{H}^{d-1, g}\left(E_{v}(h) \cap W\right)=\infty$ when $h>h_{\nu}(g)$. To prove this, let $\gamma_{g}=\lim _{r \rightarrow 0}(\log g(r)) / \log r$. Then, let us assume that $h>h_{\nu}(g)$ and remark that $h_{\nu}(g)=\gamma_{g} / \beta_{\nu}$. It follows that $\gamma_{g}<\beta_{\nu} h$, so that $g(r) \geq r^{\beta_{v} h}$ for $r \geq 0$ small enough. Thus, the Hausdorff measure of $E_{v}(h) \cap W$ for the gauge function $r \mapsto r^{d-1} g(r)$ is larger than or equal to its $\left(d-1+\beta_{v} h\right)$ dimensional measure, which is infinite as a result of Theorem 4 and the fact that $h_{v}\left(r \mapsto r^{\beta_{v} h}\right)=h$.

The proof of the second part of the corollary, which gives the value of the Hausdorff dimension of $E_{v}(h) \cap W$ and $E_{v}^{\prime}(h) \cap W$, is a consequence of the first part, together with the following observations: for $h>0$, note that $h_{v}\left(r \mapsto r^{s}\right)=s / \beta_{\nu}$ for any $s \in(0,1]$; for $h=0$, note that $h_{v}(g)=0$ for some $g \in \mathfrak{D}_{1}$ (which may be built by borrowing ideas from the proof of [18, Proposition 5]).

### 12.4 Proof of Proposition 5

As mentioned at the beginning of Sect. 11, if $v$ has finite total mass, then with probability one, $K_{\nu}(\alpha)=\emptyset$ for all $\alpha>0$. The result now follows from Corollary 5.

### 12.5 Proof of Proposition 4

The case where $h>1 / \beta_{\nu}$ follows from (25). To treat the case where $h=1 / \beta_{\nu}$, let us use of Corollary 5 in order to write that

$$
E_{\nu}^{\prime}\left(1 / \beta_{\nu}\right)=\mathbb{R}^{d} \backslash J_{v} \quad \text { and } \quad E_{v}\left(1 / \beta_{\nu}\right)=\mathbb{R}^{d} \backslash\left(J_{\nu} \cup \bigcup_{\substack{0<\alpha<1 / \beta_{\nu} \\ \alpha \in \mathbb{Q}}} K_{\nu}(\alpha)\right)
$$

In the union above, we may restrict $\alpha$ to being rational, because $\alpha \mapsto K_{\nu}(\alpha)$ is nondecreasing. Now, applying Theorem 7 with $g: r \mapsto r$, we infer that with probability one, the sets $K_{v}(\alpha)$, for $\alpha \in \mathbb{Q} \cap\left(0,1 / \beta_{\nu}\right)$, all have Lebesgue measure zero. Moreover, the set $J_{v}$ is a countable union of hyperplanes, thereby having Lebesgue measure zero as well. It follows that almost surely, the set $E_{v}\left(1 / \beta_{v}\right)$ has full Lebesgue measure in the whole space $\mathbb{R}^{d}$.

## 13 Estimates of the increments of $L_{v, j}$

The purpose of this section is to establish Lemma 2, that is, to prove the almost sure finiteness of $Z_{\nu}(A)$, for any fixed integer $A \geq 1$.

### 13.1 A net argument

Given an integer $k \geq 1$, let $\sigma_{k}=2^{-k} /\left\lfloor d^{1 / 2}\right\rfloor$. There exists $U_{A, k} \subseteq \sigma_{k} \mathbb{Z}^{d}$ with cardinality at most $\left(2^{k+2} A d^{1 / 2}\right)^{d}$ such that $\left[-A-2^{-k}, A+2^{-k}\right]^{d}$ is covered by the cubes $u+\left[0, \sigma_{k}\right)^{d}$, for $u \in U_{A, k}$. Then, for any $t \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}(A)$ and $\tau \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}\left(2^{-k}\right)$, there are two points $u$ and $u^{\prime}$ in $U_{A, k}$ such that $t-u$ and $t+\tau-u^{\prime}$ belong to $\left[0, \sigma_{k}\right)^{d}$. Moreover, writing $u=p \sigma_{k}$ and $u^{\prime}=p^{\prime} \sigma_{k}$ with $p, p^{\prime} \in \mathbb{Z}^{d}$, we see that the $\ell^{1}$-norm of $p-p^{\prime}$ is at most $3 d$, so there exists a finite sequence $\left(p_{i}\right)_{0 \leq i \leq n}$ in $\mathbb{Z}^{d}$ such that $n \leq 3 d, p_{0}=p, p_{n}=p^{\prime}$, and $\left\|p_{i+1}-p_{i}\right\|=1$ and $u_{i}=p_{i} \sigma_{k} \in U_{A, k}$ for all $i$. As a result, for any $j \geq 1$, the increment $\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right|$ is at most

$$
\left|L_{v, j}(t+\tau)-L_{v, j}\left(u^{\prime}\right)\right|+\left|L_{v, j}(t)-L_{v, j}(u)\right|+\sum_{i=0}^{n-1}\left|L_{v, j}\left(u_{i+1}\right)-L_{v, j}\left(u_{i}\right)\right|,
$$

Making use of (27) and letting

$$
\zeta_{v}(t, j, k)=\sup _{\tau \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}\left(2^{-k}\right)}\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right|
$$

for any $j, k \geq 1$ and $t \in \mathbb{R}^{d}$, it follows that

$$
\begin{equation*}
\zeta_{\nu}(A, j, k) \leq(3 d+2) \sup _{t \in U_{A, k}} \zeta_{\nu}(t, j, k) \tag{47}
\end{equation*}
$$

### 13.2 Estimates of $\zeta_{v}(t, j, k)$

Let us now derive an appropriate upper bound on $\zeta_{v}(t, j, k)$ and a control on the tail distribution of this bound. To this end, for any relatively compact Borel set $V \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$, let

$$
M_{v, j}(V)=\int_{\substack{(\rho, s) \in V \\|x| \in \mathcal{I}_{j}}} \mathrm{~N}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x) \text { and } m_{v, j}(V)=\int_{\substack{(\rho, s) \in V \\|x| \in \mathcal{I}_{j}}} \mathrm{~d} \rho \nu(\mathrm{~d} s, \mathrm{~d} x) .
$$

Recall that $\mathfrak{L}_{v, j}(V)$ is given by (10), so that we clearly have

$$
\begin{equation*}
\left|\mathfrak{L}_{v, j}(V)\right| \leq 2^{-j+1}\left(M_{\nu, j}(V)+m_{\nu, j}(V)\right) . \tag{48}
\end{equation*}
$$

Moreover, for any $t \in \mathbb{R}^{d}$ and any $\delta>0$, let $V_{t, \delta}^{\circ}=V_{t, \delta}^{+} \backslash V_{t, \delta}^{-}$, where

$$
V_{t, \delta}^{+}=\bigcup_{\tau \in \mathbb{Q}^{d} \cap \overline{\mathbf{B}}_{0}(\delta)} V_{t+\tau} \text { and } V_{t, \delta}^{-}=\bigcap_{\tau \in \mathbb{Q}^{d} \cap \overline{\mathbf{B}}_{0}(\delta)} V_{t+\tau}
$$

Note that $V_{t, \delta}^{\circ} \in \mathcal{B}_{0}\left(\mathbb{H}_{d}\right)$ and that $|\rho-\langle s, t\rangle| \leq \delta$ for any $(\rho, s) \in V_{t, \delta}^{\circ}$. Therefore, exploiting the symmetry of $\nu$, we have

$$
\begin{equation*}
m_{v, j}\left(V_{t, \delta}^{\circ}\right) \leq \frac{1}{2} \int_{\substack{s \in \mathbb{S}^{d-1} \\|x| \in \mathcal{I}_{j}}} \int_{\rho \in \mathbb{R}} \mathbb{1}_{\{|\rho-\langle s, t\rangle| \leq \delta\}} \mathrm{d} \rho v(\mathrm{~d} s, \mathrm{~d} x)=2 \delta v_{j} \tag{49}
\end{equation*}
$$

Our approach now depends on whether or not $2^{-k} v_{j} \leq \eta j k$, where $\eta \geq 1$ is a real constant to be fixed later.

### 13.2.1 Case where $2^{-k} v_{j} \leq \eta j k$

Here, the suitable bound on $\zeta_{v}(t, j, k)$ and an estimate of its tail distribution are given by the next two results.

Lemma 10 For any $t \in \mathbb{R}^{d}$ and any $\eta, j, k \geq 1$ with $2^{-k} \nu_{j} \leq \eta j k$,

$$
\zeta_{v}(t, j, k) \leq 2^{-j+1}\left(M_{v, j}\left(V_{t, 2^{-k}}^{\circ}\right)+2 \eta j k\right) .
$$

Proof For every $\tau \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}\left(2^{-k}\right)$, the increment $\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right|$ is equal to

$$
\begin{aligned}
& \left|\left(\mathfrak{L}_{v, j}\left(V_{t+\tau} \backslash V_{t}\right)+\mathfrak{L}_{v, j}\left(V_{t+\tau} \cap V_{t}\right)\right)-\left(\mathfrak{L}_{v, j}\left(V_{t} \backslash V_{t+\tau}\right)+\mathfrak{L}_{v, j}\left(V_{t} \cap V_{t+\tau}\right)\right)\right| \\
& \quad \leq 2^{-j+1}\left(M_{v, j}\left(V_{t+\tau} \Delta V_{t}\right)+m_{v, j}\left(V_{t+\tau} \Delta V_{t}\right)\right),
\end{aligned}
$$

where the last bound is due to (48). Furthermore, $V_{t+\tau} \Delta V_{t} \subseteq V_{t, 2^{-k}}^{\circ}$, so that

$$
\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right| \leq 2^{-j+1}\left(M_{\nu, j}\left(V_{t, 2^{-k}}^{\circ}\right)+m_{\nu, j}\left(V_{t, 2^{-k}}^{\circ}\right)\right)
$$

The result follows from (49) and the assumption on $j$ and $k$.
Lemma 11 For any $t \in \mathbb{R}^{d}$ and any $\eta, j, k \geq 1$ with $2^{-k} \nu_{j} \leq \eta j k$,

$$
\mathbb{P}\left(M_{\nu, j}\left(V_{t, 2^{-k}}^{\circ}\right) \geq 5 \eta j k\right) \leq 2 \mathrm{e}^{-\eta j k}
$$

Proof We may clearly assume that $m_{\nu, j}\left(V_{t, 2^{-k}}^{\circ}\right)$ is positive. In view of (49), the fact that $M_{\nu, j}\left(V_{t, 2^{-k}}^{\circ}\right) \geq 5 \eta j k$ implies that

$$
\int_{\substack{(\rho, s) \in V^{\circ} \\|x| \in \mathscr{I}_{j}^{2-k}}} \mathrm{~N}^{*}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x)=M_{v, j}\left(V_{t, 2^{-k}}^{\circ}\right)-m_{v, j}\left(V_{t, 2^{-k}}^{\circ}\right) \geq 3 \eta j k,
$$

which may happen with probability at most $2 \mathrm{e}^{-\eta j k}$, by virtue of Lemma 4.
13.2.2 Case where $2^{-k} v_{j}>\eta j k$

Let $\delta_{j, k}=\left(2^{k} v_{j}\right)^{-1 / 2}$ and $\mathrm{T}_{j, k}$ be the set of points in $\mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}\left(2^{1-k}\right)$ of the form $\left.p / L\left(d 2^{k} v_{j}\right)^{1 / 2}\right\rfloor$ with $p \in \mathbb{Z}^{d}$. We clearly have $\# \mathrm{~T}_{j, k} \leq\left(25 d 2^{-k} v_{j}\right)^{d / 2}$. Here are the analogs of Lemmas 10 and 11 .

Lemma 12 For any $t \in \mathbb{R}^{d}$ and any $\eta, j, k \geq 1$ with $2^{-k} v_{j}>\eta j k$,

$$
\begin{aligned}
\zeta_{\nu}(t, j, k) \leq & 2 \sup _{\tau, \tau^{\prime} \in \mathrm{T}_{j, k}}\left|\mathfrak{L}_{\nu, j}\left(V_{t+\tau, \delta_{j, k}}^{+} \backslash V_{t+\tau^{\prime}, \delta_{j, k}}^{-}\right)\right| \\
& +2^{3-j} \sup _{\tau \in \mathrm{T}_{j, k}} M_{\nu, j}\left(V_{t+\tau, \delta_{j, k}}^{\circ}\right)+2^{4-j-k / 2} \nu_{j}^{1 / 2} .
\end{aligned}
$$

Proof Given $\tau \in \mathbb{Q}^{d} \cap \overline{\mathrm{~B}}_{0}\left(2^{-k}\right)$, there clearly exists a $\tau^{\prime} \in \mathrm{T}_{j, k}$ with $\left\|\tau-\tau^{\prime}\right\| \leq \delta_{j, k}$, so that $V_{t+\tau^{\prime}, \delta_{j, k}}^{-} \subseteq V_{t+\tau} \subseteq V_{t+\tau^{\prime}, \delta_{j, k}}^{+}$. Moreover, we also have $V_{t, \delta_{j, k}}^{-} \subseteq V_{t} \subseteq V_{t, \delta_{j, k}}^{+}$. Then, just as in the proof of Lemma 10,

$$
\left|L_{v, j}(t+\tau)-L_{v, j}(t)\right| \leq\left|\mathfrak{L}_{v, j}\left(V_{t+\tau} \backslash V_{t}\right)\right|+\left|\mathfrak{L}_{v, j}\left(V_{t} \backslash V_{t+\tau}\right)\right| .
$$

Splitting $V_{t+\tau^{\prime}, \delta_{j, k}}^{+} \backslash V_{t, \delta_{j, k}}^{-}$into its subset $V_{t+\tau} \backslash V_{t}$ and the complement, we get

$$
\left|\mathfrak{L}_{v, j}\left(V_{t+\tau} \backslash V_{t}\right)\right|=\left|\mathfrak{L}_{v, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{+} \backslash V_{t, \delta_{j, k}}^{-}\right)-\mathfrak{L}_{\nu, j}\left(\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{+} \backslash V_{t, \delta_{j, k}}^{-}\right) \backslash\left(V_{t+\tau} \backslash V_{t}\right)\right)\right| .
$$

Owing to the triangle inequality, the upper bound given by (48) and the observation that the complement of the set $V_{t+\tau} \backslash V_{t}$ in $V_{t+\tau^{\prime}, \delta_{j, k}}^{+} \backslash V_{t, \delta_{j, k}}^{-}$is included in the union of $V_{t+\tau^{\prime}, \delta_{j, k}}^{\circ}$ and $V_{t, \delta_{j, k}}^{\circ}$, the right-hand side above is smaller than or equal to

$$
\begin{aligned}
& \left|\mathfrak{L}_{\nu, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{+} \backslash V_{t, \delta_{j, k}}^{-}\right)\right|+2^{-j+1}\left(M_{v, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{\circ} \cup V_{t, \delta_{j, k}}^{\circ}\right)\right. \\
& \left.\quad+m_{v, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{\circ} \cup V_{t, \delta_{j, k}}^{\circ}\right)\right) .
\end{aligned}
$$

In addition, (49) ensures that the sets $V_{t+\tau^{\prime}, \delta_{j, k}}^{\circ}$ and $V_{t, \delta_{j, k}}^{\circ}$ have $m_{\nu, j}$-measure at most $2 \delta_{j, k} v_{j}$. Therefore, $\left|\mathfrak{L}_{v, j}\left(V_{t+\tau} \backslash V_{t}\right)\right|$ is at most
$\left|\mathfrak{L}_{v, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{+} \backslash V_{t, \delta_{j, k}}^{-}\right)\right|+2^{-j+1}\left(M_{v, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{\circ}\right)+M_{v, j}\left(V_{t, \delta_{j, k}}^{\circ}\right)+4 \delta_{j, k} v_{j}\right)$.
Likewise, $\left|\mathfrak{L}_{v, j}\left(V_{t} \backslash V_{t+\tau}\right)\right|$ is smaller than or equal to
$\left|\mathfrak{L}_{v, j}\left(V_{t, \delta_{j, k}}^{+} \backslash V_{t+\tau^{\prime}, \delta_{j, k}}^{-}\right)\right|+2^{-j+1}\left(M_{v, j}\left(V_{t+\tau^{\prime}, \delta_{j, k}}^{\circ}\right)+M_{v, j}\left(V_{t, \delta_{j, k}}^{\circ}\right)+4 \delta_{j, k} \nu_{j}\right)$,
and the result follows.

Lemma 13 For any $t \in \mathbb{R}^{d}$, any $\eta, j, k \geq 1$ with $2^{-k} \nu_{j}>\eta j k$, and any $\tau, \tau^{\prime} \in \mathrm{T}_{j, k}$,

$$
\left\{\begin{array}{l}
\left.\mathbb{P}\left(\mid \mathfrak{L}_{\nu, j}\left(V_{t+\tau, \delta_{j, k}}^{+}\right) V_{t+\tau^{\prime}, \delta_{j, k}}^{-}\right) \mid \geq 2^{3-j-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}\right) \leq 2 \mathrm{e}^{-\eta j k} \\
\mathbb{P}\left(M_{\nu, j}\left(V_{t+\tau, \delta_{j, k}}^{\circ}\right) \geq 2^{3-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}\right) \leq 2 \mathrm{e}^{-\eta j k}
\end{array}\right.
$$

Proof For the first bound, in view of (48), we may clearly assume that the $m_{\nu, j}$-measure of $V_{t+\tau, \delta_{j, k}}^{+} \backslash V_{t+\tau^{\prime}, \delta_{j, k}}^{-}$is positive. Moreover, this set is included in $V_{t, 3 \cdot 2-k}^{\circ}$, so its $m_{v_{j}}$-measure is at most $6 \cdot 2^{-k} v_{j}$, owing to (49). Thus, Lemma 4 and the fact that $2^{-\dot{k}} v_{j}>\eta j k$ imply that the probability under study is at most

$$
2 \exp \left(-\frac{3 \cdot 2^{6-2 j-k} \eta j k v_{j}}{2^{5-2 j-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}+36 \cdot 2^{2-2 j-k} v_{j}}\right) \leq 2 \mathrm{e}^{-\eta j k}
$$

The second inequality that we need to establish is an upper bound on the probability of an event which implies that

$$
\begin{aligned}
\int_{\substack{(\rho, s) \in V_{t+\tau, \delta_{j, k}}^{|x| \in \mathcal{I}_{j}}}} \mathrm{~N}^{*}(\mathrm{~d} \rho, \mathrm{~d} s, \mathrm{~d} x) & =M_{v, j}\left(V_{t+\tau, \delta_{j, k}}^{\circ}\right)-m_{v, j}\left(V_{t+\tau, \delta_{j, k}}^{\circ}\right) \\
& \geq 2^{3-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}-2 \delta_{j, k} v_{j} \geq 4\left(\eta j k 2^{-k} v_{j}\right)^{1 / 2}
\end{aligned}
$$

Thanks to Lemma 4 and (49) again, this may happen with probability at most

$$
2 \exp \left(-\frac{3 \eta j k 2^{4-k} v_{j}}{2^{3-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}+12\left(2^{-k} v_{j}\right)^{1 / 2}}\right) \leq 2 \mathrm{e}^{-\eta j k}
$$

where the last bound follows from the fact that $2^{-k} v_{j}>\eta j k \geq 1$.

### 13.3 End of the proof

Let us consider a real $\eta \geq 1$. If $2^{-k} \nu_{j} \leq \eta j k$, let $\mathcal{E}_{j, k}$ denote the event that the following does not hold:

$$
\begin{equation*}
\sup _{t \in U_{A, k}} M_{v, j}\left(V_{t, 2^{-k}}^{\circ}\right) \leq 5 \eta j k \tag{50}
\end{equation*}
$$

Thanks to Lemma 11, its probability satisfies

$$
\mathbb{P}\left(\mathcal{E}_{j, k}\right) \leq 2 \mathrm{e}^{-\eta j k} \# U_{A, k} \leq 2\left(4 A d^{1 / 2}\right)^{d} 2^{d k} \mathrm{e}^{-\eta j k}
$$

Now, if $2^{-k} v_{j}>\eta j k$, let $\mathcal{E}_{j, k}$ denote the event that the following does not hold:

$$
\left\{\begin{array}{l}
\sup _{\substack{t \in U_{A, k} \\
\tau, \tau^{\prime} \in \mathrm{T}_{j, k}}}\left|\mathfrak{L}_{j}\left(V_{t+\tau, \delta_{j, k}}^{+} \backslash V_{t+\tau^{\prime}, \delta_{j, k}}^{-}\right)\right| \leq 2^{3-j-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}  \tag{51}\\
\sup _{\substack{t \in U_{A, k} \\
\tau \in \mathrm{~T}_{j, k}}} M_{v, j}\left(V_{t+\tau, \delta_{j, k}}^{\circ}\right) \leq 2^{3-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}
\end{array}\right.
$$

Owing to Lemma 13, we have

$$
\begin{aligned}
\mathbb{P}\left(\mathcal{E}_{j, k}\right) & \leq 2 \mathrm{e}^{-\eta j k} \# U_{A, k} \# \mathrm{~T}_{j, k}\left(1+\# \mathrm{~T}_{j, k}\right) \\
& \leq 4\left(100 A d^{3 / 2}\right)^{d} v_{j}^{d} \mathrm{e}^{-\eta j k} \leq 4\left(100 A d^{3 / 2} c_{\nu}\right)^{d} 2^{2 d j} \mathrm{e}^{-\eta j k},
\end{aligned}
$$

where $c_{v}=\sum_{j \geq 1} 2^{-2 j} v_{j}$, which is finite owing to (1).
From now on, let us suppose that $\eta>2 d \log 2$. Using the above bounds, it is easy to check that $\sum_{(j, k) \in \mathbb{N}^{2}} \mathbb{P}\left(\mathcal{E}_{j, k}\right)<\infty$. Letting $D_{n}$ be the set of $(j, k) \in \mathbb{N}^{2}$ with $\max \{j, k\} \geq n$, we deduce that $\mathbb{P}\left(\bigcap_{n=1}^{\infty} \downarrow \bigcup_{(j, k) \in D_{n}} \mathcal{E}_{j, k}\right)=0$. So, with probability one, there exists an integer $n \geq 1$ such that (50) holds for any $(j, k) \in D_{n}$ with $2^{-k} v_{j} \leq \eta j k$, and (51) holds for any $(j, k) \in D_{n}$ with $2^{-k} v_{j}>\eta j k$. In the first case, it follows from Lemma 10 and (47) that $\zeta_{v}(A, j, k)$ is bounded by $14(3 d+2) \eta 2^{-j} j k$. In the second case, it is bounded by $96(3 d+2) 2^{-j-k / 2}\left(\eta j k v_{j}\right)^{1 / 2}$, owing to Lemma 12 and (47) again. Letting $\eta=2 d>2 d \log 2$, we finally get

$$
\text { a.s. } \exists n \geq 1 \quad \sup _{(j, k) \in D_{n}} \frac{\zeta_{\nu}(A, j, k)}{2^{-j} k\left(j+2^{-k / 2}\left(j v_{j}\right)^{1 / 2}\right)} \leq 192(3 d+2) d .
$$

In addition, for any $j, k \geq 1$, we deduce from (48) that the expectation of $\zeta_{\nu}(A, j, k)$ is at most $2^{4-j}\left(A+2^{-k}\right) v_{j}$, so that $\zeta_{v}(A, j, k)<\infty$ almost surely. Thus,

$$
\text { a.s. } \forall n \geq 2 \quad \sup _{(j, k) \in \mathbb{N}^{2} \backslash D_{n}} \frac{\zeta_{v}(A, j, k)}{2^{-j} k\left(j+2^{-k / 2}\left(j v_{j}\right)^{1 / 2}\right)}<\infty .
$$

Lemma 2 now clearly follows.

## References

1. Adler, R., Monrad, D., Scissors, R., Wilson, R.: Representations, decompositions and sample function continuity of random fields with independent increments. Stochastic Process. Appl. 15(1), 3-30 (1983)
2. Adler, R., Taylor, J.: Random Fields and Geometry. Springer Monographs in Mathematics, vol. 17. Springer, New York (2007)
3. Arneodo, A., Bacry, E., Muzy, J.F.: Random cascades on wavelet dyadic trees. J. Math. Phys. 39(8), 4142-4164 (1998)
4. Aubry, J.-M., Jaffard, S.: Random wavelet series. Comm. Math. Phys. 227(3), 483-514 (2002)
5. Aubry, J.-M., Maman, D., Seuret, S.: Local behavior of traces of Besov functions: Prevalent results (2010). arXiv:1002.3123
6. Ayache, A., Jaffard, S., Taqqu, M.S.: Wavelet construction of generalized multifractional processes. Rev. Mat. Iberoam. 23(1), 327-370 (2007)
7. Barral, J., Fournier, N., Jaffard, S., Seuret, S.: A pure jump Markov process with a random singularity spectrum. Ann. Probab. 38(5), 1924-1946 (2010)
8. Barral, J., Seuret, S.: Heterogeneous ubiquitous systems in $\mathbb{R}^{d}$ and Hausdorff dimension. Bull. Braz. Math. Soc. 38(3), 467-515 (2007)
9. Barral, J., Seuret, S.: Ubiquity and large intersections properties under digit frequencies constraints. Math. Proc. Cambridge Philos. Soc. 145(3), 527-548 (2008)
10. Ben Slimane, M.: Multifractal formalism and anisotropic selfsimilar functions. Math. Proc. Cambridge Philos. Soc. 124(2), 329-363 (1998)
11. Beresnevich, V., Dickinson, D., Velani, S.: Measure theoretic laws for limsup sets. Mem. Am. Math. Soc. 179(846), 1-91 (2006)
12. Bertoin, J.: Lévy Processes. Cambridge Tracts in Mathematics, vol. 121. Cambridge University Press, Cambridge (1996)
13. Blumenthal, R.M., Getoor, R.K.: Sample functions of stochastic processes with stationary independent increments. J. Appl. Math. Mech. 10, 493-516 (1961)
14. Chentsov, N.N.: Lévy's Brownian motion for several parameters and generalized white noise. Theory Probab. Appl. 2, 265-266 (1957)
15. Durand, A.: Random wavelet series based on a tree-indexed Markov chain. Comm. Math. Phys. 283(2), 451-477 (2008)
16. Durand, A.: Sets with large intersection and ubiquity. Math. Proc. Cambridge Philos. Soc. 144(1), 119144 (2008)
17. Durand, A.: Large intersection properties in Diophantine approximation and dynamical systems. J. Lond. Math. Soc. 79(2), 377-398 (2009)
18. Durand, A.: Singularity sets of Lévy processes. Probab. Theory Relat. Fields 143(3-4), 517-544 (2009)
19. Durand, A.: On randomly placed arcs on the circle. In: Barral, J., Seuret, S. (eds.) Recent Developments in Fractals and Related Fields, pp. 343-351. Applied and Numerical Harmonic Analysis, vol. 67. Birkhäuser, Boston (2010)
20. Dvoretzky, A.: On the oscillation of the Brownian motion process. Israel J. Math. 1, 212-214 (1963)
21. Falconer, K.J.: Sets with large intersection properties. J. Lond. Math. Soc. 49(2), 267-280 (1994)
22. Falconer, K.J.: Fractal Geometry: Mathematical Foundations and Applications. 2nd edn. Wiley, Chichester (2003)
23. Houdré, C., Privault, N.: Concentration and deviation inequalities in infinite dimensions via covariance representations. Bernoulli 8(6), 697-720 (2002)
24. Jaffard, S.: The multifractal nature of Lévy processes. Probab. Theory Relat. Fields 114(2), 207227 (1999)
25. Jaffard, S.: Wavelet techniques in multifractal analysis. In: Lapidus, M., van Frankenhuysen, M. (eds.) Fractal geometry and applications: a jubilee of Benoît Mandelbrot, pp. 91-151. Proceedings of Symposia in Pure Mathematics, vol. 72. American Mathematical Society, Providence (2004)
26. Jaffard, S.: Pointwise and directional regularity of nonharmonic Fourier series. Appl. Comput. Harmon. Anal. 28(3), 251-266 (2010)
27. Kingman, J.F.C.: Poisson Processes. Oxford Studies in Probability 3. The Clarendon Press, Oxford University Press, New York (1993)
28. Lagaize, S.: Hölder exponent for a two-parameter Lévy process. J. Multivariate Anal. 77(2), 270285 (2001)
29. Ledoux, M., Talagrand, M.: Probability in Banach Spaces: Isoperimetry and Processes. Ergebnisse der Mathematik und ihrer Grenzgebiete, vol. 23, no. 3. Springer, Berlin (1991)
30. Lévy, P.: Processus stochastiques et mouvement brownien. 2nd edn. Gauthier-Villars, Paris (1965)
31. Lifshits, M.A.: Gaussian Random Functions. Mathematics and its Applications, vol. 322. Kluwer, Dordrecht (1995)
32. Mattila, P.: Geometry of Sets and Measures in Euclidean Spaces: Fractals and Rectifiability. Cambridge Studies in Advanced Mathematics, vol. 44. Cambridge University Press, Cambridge (1995)
33. Mori, T.: Representation of linearly additive random fields. Probab. Theory Relat. Fields 92(1), 91115 (1992)
34. Neveu, J.: Processus ponctuels. In: École d'Été de Probabilités de Saint-Flour, VI-1976, pp. 249-445. Lecture Notes in Mathematics, vol. 598. Springer, Berlin (1977)
35. Reynaud-Bouret, P.: Adaptive estimation of the intensity of inhomogeneous Poisson processes via concentration inequalities. Probab. Theory Relat. Fields 126(1), 103-153 (2003)
36. Rogers, C.: Hausdorff Measures. Cambridge University Press, Cambridge (1970)
37. Samorodnitsky, G., Taqqu, M.S.: Stable Non-Gaussian Random Processes: Stochastic Models with Infinite Variance. Stochastic Modeling. Chapman \& Hall, New York (1994)
38. Sato, K.I.: Lévy Processes and Infinitely Divisible Distributions. Cambridge Studies in Advanced Mathematics. Cambridge University Press, Cambridge (1999)
39. Shieh, N.-R.: Sample functions of Lévy-Chentsov random fields. In: Probability Theory and Mathematical Statistics (Tokyo, 1995), pp. 450-459. World Scientific Publishing, River Edge (1996)
40. Sprindžuk, V.: Metric Theory of Diophantine Approximation. Wiley, New York (1979)
41. Straf, M.L.: Weak convergence of stochastic processes with several parameters. In: Proceedings of the 6th Berkeley Symposium on Mathematical Statistics and Probability (University of California, Berkeley, CA, 1970/1971), vol. II: Probability Theory, pp. 187-221. University of California Press, Berkeley (1972)
42. Vares, M.E.: Some results on local growth of two-parameter Lévy processes. Bol. Soc. Brasil. Mat. 12(2), 35-55 (1981)
43. Wendt, H., Abry, P., Jaffard, S., Ji, H., Shen, Z.: Wavelet leader multifractal analysis for texture classification. In: Proceedings of IEEE Conference on ICIP 2009, pp. 3829-3832 (2009)

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