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Variance linearity for real Gaussian zeros

Raphaël Lachièze-Rey*

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Abstract We investigate the zero set of a stationary Gaussian process on the real line, and in particular give lower bounds for the variance of the number of points on a large interval, in all generality. We prove that this point process is never hyperuniform, i.e. the variance is at least linear, and give necessary conditions to have linear variance, which are close to be sharp. We study the class of symmetric Bernoulli convolutions and give an example where the zero set is super rigid, weakly mixing, and not hyperuniform.

1 Background and motivation

We study here the zero set of a Gaussian stationary process X:

$$\mathsf{Z}_{\mathsf{X}} := \{ t \in \mathbb{R} : \mathsf{X}(t) = 0 \}.$$

The study of Gaussian zeros has emerged in the fifties, with pioneering works of Kac and Rice [15], and subsequent applications in the fields of telecommunications and signal processing, ocean waves, and random mechanics. Rigorous results can be found in the book of Cramer and Leadbetter in the sixties [5], along with the first second order results. In the following decades, most significant second order analyses have been made by Cuzick [6], Slud [23], Kratz & Léon [17], see the survey [16] and references therein. The weakest available condition for non-degenerate asymptotic variance is the square integrability of the covariance, in [23]. We give here a variance lower bound, under virtually no hypotheses, that implies that the variance is always at least linear, and give a necessary condition very close to known sufficient conditions for it to be actually linear. We also study the class of symmetric Bernoulli convolutions and focus on an example whose spectral measure has a Cantor nature.

Feldheim [9] has also studied zeros of Gaussian analytic functions (GAFs) whose law is invariant only under horizontal shifts, reduced to a horizontal strip of the form $\mathbb{R} \times [a, b]$, with $-\infty < a < b < \infty$. In a related work, Buckley and Feldheim [2] study the winding number of a GAF X: $\mathbb{R} \mapsto \mathbb{C}$. A strong motivation of their work is the analogy between the winding number and the number of zeros. Some convenient identities coming from complex analysis can provide explicit expressions for these indexes and their moments, whereas the number of zeros of a real Gaussian process is mostly studied in the literature through its Wiener-Ito expansion, which is sometimes not amenable to analysis; one can also use directly Kac-Rice formula [1] but this does not ease the task of estimating the variance. In particular, we cannot obtain an expression as explicit as (6) in [2] or Section 3.2 in [9] for the number of zeros. It will be interesting to observe in Section 2.1 that the results of these two works about GAFs are similar to the results we obtain here: the variance of the number of zeros (or the winding number) is always at least linear, it is

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not linear if some square integrability conditions related to the covariance functions are not satisfied, and it is quadratic if the spectral measure has atoms.

Another motivation is to explore the behaviour of Gaussian zeros in the light of the concepts of hyperuniformity and rigidity. Let us introduce more formally point processes before going further. Let $\mathcal{N}(\mathbb{R})$ be the space of locally finite subsets of \mathbb{R} , endowed with the σ -algebra $\mathcal{B}(\mathcal{N}(\mathbb{R}))$ generated by the mappings

$$\varphi_A: \mathsf{Z} \in \mathscr{N}(\mathbb{R}) \mapsto \#\mathsf{Z} \cap A$$

for A a Borel subset of \mathbb{R} , where # denotes the cardinality of a set. A point process Z is a measurable mapping from an underlying probability space $(\Omega; \mathscr{A}, \mathbf{P})$ to $(\mathscr{N}(\mathbb{R}), \mathcal{B}(\mathscr{N}(\mathbb{R})))$;

Z is furthermore stationary if $Z + x \stackrel{(d)}{=} Z$ for $x \in \mathbb{R}$, and we denote by $V_Z(T) = Var(\#Z \cap [-T,T])$. In physics of condensed matters and statistical physics, a great attention has been given recently to the phenomenon of *fluctuation suppression* of particle systems, which can be incarnated by the property of *hyperuniformity*: a stationary point process Z on \mathbb{R} is *hyperuniform*, or *superhomegeneous*, if

$$\liminf_{T \to \infty} \frac{V_{\mathsf{Z}}(T)}{T} = 0.$$

A striking property of many hyperuniform systems with suppressed fluctuations is their rigidity, i.e. the complete determination of a feature of the point process in a bounded domain given its configuration outside the domain: for a measurable mapping $\varphi : \mathcal{N}(\mathbb{R}) \to \mathbb{R}$, we say a point process Z is φ -rigid if for all T > 0,

$$\varphi(\mathsf{Z}\cap[-T,T])\in\sigma(\mathsf{Z}\cap[-T,T]^c).$$

In particular, we say it is *super rigid* if it is φ -rigid for every bounded measurable φ , i.e. if $\mathsf{Z} \cap [-T,T]$ is completely determined by $\mathsf{Z} \cap [-T,T]^c$ for T>0.

This phenomenon has been much documented, with many rigourous results proven in dimension 1 and 2, and recently in dimension 3 by Chatterjee on hierarchichal Coulomb gases [3]. The link between rigidity and hyperuniformity is the topic of many works in the recent years. Most notable examples of stationary hyperuniform models are determinantal processes whose kernels are projection operators [12], zeros of planar Gaussian analytic functions [13], some Coulomb systems [12], the β -sine processes for $\beta > 0$ [7], perturbed lattices [21], and also many non-stationary models, with a more general definition of hyperuniformity. In all these models, the point process is hyperuniform and φ -rigid, at least for $\varphi(Z) = \#Z$. It is widely believed that hyperuniformity is required to have rigidity [12], and the usual method of proof of rigidity, introduced in [13], involves linear statistics and implictly assumes hyperuniformity (see [7] for an alternative method involving DLR equations). The phenomenon of rigidity, although mathematically intriguing, has also practical uses in some percolation models related to the underlying point process, see [11].

We exploit the variance results derived in this paper to exhibit a non artificial example of a point process in dimension 1 that experiences the strongest possible form of rigidity, without being hyperuniform. Such examples are actually not so hard to build on \mathbb{R}^d , take for instance

$$\mathsf{Z} = L(\mathbb{Z}^d + U)$$

where L is some non-trivial non-negative L^2 real random variable and U is independent and uniformly distributed on $[0,1)^d$. It is also clear that there is absolutely no form of asymptotic independence in Z, in the sense that given two events A, B, the events $Z \in A$ and $Z + t \in B$ have no reason to be uncorrelated for large T. For this reason, we might ask for additional desirable properties of such counter-examples.

Definition 1. We say that a stationary point process Z of \mathbb{R} is weakly mixing if for all events $A, B \in \mathcal{B}(\mathcal{N}(\mathbb{R}))$,

$$\lim_{T\to\infty}T^{-1}\int_0^T|\mathbf{P}(\mathsf{Z}\in A,(\mathsf{Z}+t)\in B)-\mathbf{P}(\mathsf{Z}\in A)\mathbf{P}(\mathsf{Z}\in B)|dt\to 0.$$

The randomly scaled grid Z in \mathbb{R} does not satisfy this property because it is not even ergodic. We give here a stationary Gaussian process X whose spectral measure belongs to the class of *symmetric Bernoulli convolutions* and whose zero set Z_X is weakly mixing, super rigid, and not hyperuniform. See Propositions 9 and 12.

Regarding rigidity and hyperuniformity, stationary point processes obtained as the zeros of a random Gaussian functions have mainly been studied for *Gaussian analytic functions* in the complex plane, see [14]. The joint requirements of complex Gaussianity, analycity, and stationarity have reduced the class of zeros to essentially one model, up to rescaling, which happens to be rigid and hyperuniform [14]. The link between Gaussian zeros and rigidity is yet essentially to discover, especially in higher dimensions.

We present in Section 2 some results and definitions about Gaussian processes and their zeros, variance lower bounds, and we give a condition for super rigidity. We also properly explore the analogies with variance results for GAFs. In Section 3 we introduce the class of symmetric Bernoulli convolutions, give some properties and examples, in particular an example yielding super rigid non hyperuniform zeros. In Section 4 we prove the lower bounds on the variance by using chaotic decompositions of the number of zeros.

2 Gaussian zeroes

We first give basic results about Gaussian processes. See the monograph of Azaïs & Wschebor [1] for background material and advanced results about Gaussian processes and their level sets. A real Gaussian process is a random function $X : \mathbb{R} \to \mathbb{R}$ for which any finite dimensional distribution $(X(t_1), \ldots, X(t_n)), t_1, \ldots, t_n \in \mathbb{R}$, is a Gaussian vector with covariance matrix $(\mathcal{C}(t_i, t_j))_{1 \le i,j \le n}$, where \mathcal{C} is the covariance function of the process X. The law of X is entirely determined by \mathcal{C} , and the stationarity of X, defined by $X(\cdot + t) \stackrel{(d)}{=} X$ for $t \in \mathbb{R}$, is characterised by the existence of a function C_X , called reduced covariance function, such that $\mathcal{C}(t,s) = C_X(t-s), t,s \in \mathbb{R}$. Such a function is known to admit a spectral representation of the form

$$C_{\mathsf{X}}(t) = \int_{\mathbb{R}} e^{ixt} \mu_{\mathsf{X}}(dx) = 2 \int_{0}^{\infty} \cos(xt) \mu_{\mathsf{X}}(dx), t \in \mathbb{R},$$

where μ_X is the (symmetric) spectral measure of X on \mathbb{R} . We say that X is degenerate if $C_X(x) = A\cos(\varphi x)$ for some $A \ge 0, \varphi \in \mathbb{R}$, or equivalently if μ_X is formed by two symmetric atoms with same mass.

2.1 Zero set

We are primarily interested in the variance $V_{\mathsf{Z}_{\mathsf{X}}}(T)$ of $\#(\mathsf{Z}_{\mathsf{X}} \cap [-T, T])$, denoted by $V_{\mathsf{X}}(T)$ for short, where X is a Gaussian stationary process. The Kac-Rice formula, established in the fourties [15], is the starting point of most rigourous subsequent works. Cramer and Leadbetter [5] proved that it is necessary for $V_{\mathsf{X}}(T)$ to be finite for every T > 0 that C_{X} is twice differentiable in 0 and for some $\delta > 0$,

$$\int_0^{\delta} \frac{1}{t^2} (C_{\mathsf{X}}'(t) - C_{\mathsf{X}}''(0)t) dt < \infty. \tag{1}$$

Geman [10] then proved that the condition is actually necessary on any non-negligible interval, and this condition is thus referred to as *Geman's condition*. When this condition is satisfied, we use without loss of generality the convention $C''_{\mathsf{X}}(0) = -1$, equivalent to replacing $\mathsf{X}(t)$ by $\mathsf{X}(\alpha t)$ with $\alpha = (-C''_{\mathsf{X}}(0))^{-1/2}$.

By using approximation by m-dependent fields, Cuzick [6] proves that the variance is not super-linear, that is $\lim_{T\to\infty} T^{-1}V_{\mathsf{X}}(T)$ exists and is finite, if Geman's condition holds, as well as

$$\int C_{\mathsf{X}}(t)^2 dx < \infty, \int C_{\mathsf{X}}''(t)^2 dx < \infty. \tag{2}$$

He actually proves a central limit theorem under the additional assumption that

$$\lim_{T\to\infty} T^{-1}V_{\mathsf{X}}(T) > 0.$$

Computing the decomposition of the number of zeros with respect to Wiener-Ito integrals based on the field X, Slud [23] proves that (2) is sufficient for the asymptotic variance to be finite and positive, and for the central limit theorem to hold. Kratz & Léon [17] compute the chaos decomposition with respect to the joint process formed by the field X and its derivative X' and ease some computations related to the decomposition, allowing them to generalise the results to other crossing problems, this is the approach we will use at Section 4. We present here our results regarding the variance lower bound for the number of zeros. Say that μ_X is regular in $x \in \mathbb{R}$ if it admits a L^2 density on some neighbourhood of x.

Theorem 2. Let X be a non-degenerate stationary Gaussian field on \mathbb{R} . Then

(i) The variance is at least linear, i.e. Z_X is not hyperuniform:

$$\liminf_{T>0} T^{-1}V_{\mathsf{X}}(T) \in (0,\infty].$$

(ii) the condition $C'' - C \in L^2(\mathbb{R})$ is necessary for linear fluctuations, i.e. for

$$\liminf_{T \to \infty} T^{-1} V_{\mathsf{X}}(T) < \infty.$$

If μ_X is regular in ± 1 , this condition is equivalent to (2), which is sufficient for linear fluctuations.

(iii) If μ_X has an atom, $\inf_{T>0} T^{-2}V_X(T) > 0$.

The proof is at Section 4. We also give an example of a singular nonatomic spectral measure μ_X for which $T^{2-\varepsilon} = o(V_X(T))$ for $\varepsilon > 0$. Regarding the quadratic variance, as soon as the variance is finite on some interval with non-empty interior (Geman's condition), then the variance cannot be more than quadratic, see [1].

It is not clear which of these results already appeared in full generality, and in particular when μ does not have a density. We could not locate another necessary condition for linear variance in the literature, see also Remark 1.2 in [9].

For positive asymptotic variance to hold (i), we can notice that Cuzick [6] needs

$$\int \frac{(C'(t))^2}{1 - C(t)^2} dt < \frac{\pi}{2} \sqrt{\int x^2 \mu(dx)}$$

for $\lim_{T\to\infty} T^{-1}V_{\mathsf{X}}(T)$ to be strictly positive, a requirement still present in [16]. Other works aimed at proving central limit theorems with linear variance give various types of conditions: $C\in L^2$ in [23], $C\in L^1$, $C^{''}\in L^2$ in [17].

It is interesting to compare the results above to those obtained by Feldheim [9], who studied the zero set $\mathsf{Z}_{\mathsf{x}}^{[a,b]}$ of a Gaussian analytic function (GAF) X on the infinite strip

$$\mathbb{S}_{a,b} = \{ x + iy : x \in \mathbb{R}, a \leqslant y \leqslant b \},\$$

for some $-\infty < a < b < \infty$, and those of Buckley and Feldheim [2] who study the winding number W_T of a GAF X : $[0,T] \mapsto \mathbb{C}$. A GAF X is also characterised by its reduced covariance C_X and its spectral measure ρ_X along the real line:

$$C_{\mathsf{X}}(t-s) = \mathbf{E}(\mathsf{X}(t)\overline{\mathsf{X}}(s)) = \int_{\mathbb{R}} e^{-2\pi i(t-s)x} \rho_{\mathsf{X}}(dx), t, s \in \mathbb{R}.$$

Let us reproduce some of their results regarding the variances

$$V_{\mathsf{X}}^{[a,b]}(T) = \mathbf{Var}(\#\mathsf{Z}_{\mathsf{X}}^{[a,b]} \cap ([-T,T] \times [a,b])), \qquad V_{\mathsf{X}}'(T) = \mathbf{Var}(W_T).$$

The field X is called *degenerate* if ρ_X consists of a single atom.

Theorem 3 ([9, 2]). If X is a non-degenerate GAF,

(i) If the differentiability condition (3) in [2] holds, for T > 0

$$T^{-1}V'_{\mathsf{X}}(T) \geqslant C$$
 for some $C > 0$.

This expression has a limit in $(0, \infty)$ if $C_X, C_X' \in L^2$.

(ii) The limit

$$\lim_{T \to \infty} T^{-1} V_{\mathsf{X}}^{[a,b]}(T)$$

exists and is strictly positive. It is finite if

$$\int |C_{\mathsf{X}}(t+iy)|^2 dt < \infty, \int |C_{\mathsf{X}}''(t+iy)|^2 dt < \infty, y \in \{a,b\}$$

where C_X is analytically continued to $S_{a,b}$,

(ii) the limits

$$\lim_{T \to \infty} T^{-2} V_{\mathsf{X}}^{[a,b]}(T), \quad \lim_{T \to \infty} T^{-2} V_{\mathsf{X}}'(T)$$

exist and are finite, and are strictly positive iff ρ_X has an atom.

Hence the results are similar to ours regarding the quadratic variance and the linear variance. Feldheim [9] also gives a lower bound which yields a sufficient condition for asymptotic super linear variance, but leaves linearity undecided for some measures with a singularity, such as for $\rho_X(dx) = \mathbf{1}_{\{|x| \le 1\}} |x|^{-1/2} dt, x \in \mathbb{R}$. This is analoguous to our "dark spot" regarding linear variance, at the difference that we allow singularities as long as they are not located at ± 1 (under the convention C(0) = -C''(0) = 1.)

Remark 4. A degenerate GAF is, up to rescaling and multiplication, $X(t) = \gamma e^{2i\pi t}$, where $\gamma = \alpha + i\beta$ is a Gaussian standard complex variable (i.e. α, β are i.i.d real Gaussian variables), whereas a degenerate real Gaussian process in the sense of this paper is of the form $X(t) = \alpha \cos(t) + \beta \sin(t)$, up to rescaling and multiplication. In particular, the latter process, which can be extended to a GAF on \mathbb{C} , satisfies: $V_X(T)$ is bounded, but $V_X^{[a,b]}(T)$ and $V_X'(T)$ behave in T^2 .

2.2 Predictability and rigidity

We will actually consider a property stronger than super rigidity.

Definition 5. Say that a point process Z is completely predictable if

$$\sigma(\mathsf{Z}) = \bigcap_{t \in \mathbb{R}} \sigma(\mathsf{Z} \cap (-\infty, t)).$$

This property yields that the knowledge of Z on any half line determines the whole process. We have the obvious fact that if Z is completely predictable, it is super rigid.

The problem of wether the data of a stationary Gaussian process on a half line $\{X(t); t \leq 0\}$ determines the whole process has been widely investigated, see Szegö's alternative ([8, Section 4.2]). Such a behaviour does not necessarily imply that the knowledge of the zeros on a half line determines all the zeros. One might also try to reconstruct the whole process from the zeros, but Weierstrass factorisation theorem yields uncountably many functions with a given infinite set of zeros (without accumulation point). It would still be interesting to try to characterise processes whose zero set is completely predictable, or exhibit a weaker form of rigidity. We give here a simple condition on C_X for the predictability of Z_X .

Proposition 6. Assume X is almost surely of class C^1 , and there is a sequence $t_n \to \infty$ such that $|C_X(t_n)| \to 1$. Then Z_X is predicable.

Proof. Let us denote the shifted process by $X_n = X(t_n + \cdot)$. We have for $t \in \mathbb{R}$, $\pm \text{Cov}(X(t), X_n(t)) \to 1$, hence $\pm X_n(t) \to X(t)$ since these are Gaussian variables.

Let now I be a bounded interval of \mathbb{R} . Denote by \dot{I} the interior of I and by \bar{I} its closure. Since $\mathsf{Z}_\mathsf{X}, \mathsf{Z}_\mathsf{X}_n$ are stationary point processes, with probability 1, they don't have zeros on $\partial I = \bar{I} \setminus \dot{I}$. Hence a.s., for simultaneously all n,

$$(\mathsf{Z}_{\mathsf{X}_n} \cap \bar{I} \neq \emptyset) = (\mathsf{Z}_{\mathsf{X}_n} \cap \dot{I} \neq \emptyset); \ (\mathsf{Z}_{\mathsf{X}} \cap \bar{I} \neq \emptyset) = (\mathsf{Z}_{\mathsf{X}} \cap \dot{I} \neq \emptyset).$$

Since (X(0), X'(0)) is not degenerate, the zeros of X are almost all regular, i.e. X and X' don't vanish simultaneously (see [16, Th. 1.1]). Hence for each zero z of X we have X(t)X(s) < 0 for some s, t arbitrarily close from z.

Then we have almost surely

$$\mathsf{Z}_\mathsf{X} \cap \dot{I} \neq \emptyset \Rightarrow \mathsf{X}(t)\mathsf{X}(s) < 0 \text{ for some } t,s \in \dot{I} \cap \mathbb{Q}$$

 $\Rightarrow \exists t,s \in \dot{I} \cap \mathbb{Q} : \mathsf{X}_n(s)\mathsf{X}_n(t) < 0 \text{ for } n \text{ sufficiently large}$

because we have a.s. the simultaneous convergences $X_n(t) \to X(t), t \in \mathbb{Q}$, hence we can conclude that if X vanishes in \dot{I} , so does X_n for n sufficiently large.

Let us prove the converse statement. We have by Fatou's lemma, and using the stationarity again,

$$\mathbf{E}(\liminf_n \|\mathsf{X}'_n\|_I) \leqslant \liminf_n \mathbf{E}(\|\mathsf{X}'_n\|_I) = \mathbf{E}(\|\mathsf{X}'\|_I)$$

and the last quantity is finite since X' is a.s. bounded on I (see [1]). Hence

$$\mathbf{P}(\liminf_n |\mathsf{X}'_n| = \infty) = 0$$

and there is a.s. $L < \infty$ and $N \subset \mathbb{N}$ infinite such that $\|X'_n\|_I \leq L$ for $n \in \mathbb{N}$.

So let us assume that $X_n(t_n) = 0$ with $t_n \in I$ for n sufficiently large. Then it also vanishes for $n \in N$ sufficiently large, hence there is a subsequence $n' \in N$ such that $X_{n'}(t_{n'}) = 0$. Let us take another subsequence n'' such that $t_{n''} \to t \in \overline{I}$. We have

$$|X(t)| \le |X(t) - X_{n''}(t)| + |X_{n''}(t) - X_{n''}(t_{n''})| \le |X(t) - X_{n''}(t)| + L|t - t_{n''}|,$$

hence the latter a.s. converges to 0, and X vanishes in $t \in \overline{I}$. since X does not vanish on ∂I a.s., we have proved the a.s. equivalence,

$$\mathsf{Z}_{\mathsf{X}} \cap I \neq \emptyset \Leftrightarrow \mathsf{Z}_{\mathsf{X}_n} \cap I \neq \emptyset$$
 for n sufficiently large.

In particular, $Z_X \cap I \neq \emptyset$ is completely determined by

$$\bigcap_n \sigma(\mathsf{Z}_{\mathsf{X}_n} \cap I) = \bigcap_n \sigma((\mathsf{Z}_{\mathsf{X}} + t_n) \cap I) = \bigcap_n \sigma(\mathsf{Z}_{\mathsf{X}} \cap (I - t_n)) \subset \bigcap_t \sigma((\mathsf{Z}_{\mathsf{X}} \cap [t, \infty))),$$

which concludes the proof.

Let us give simple examples of Gaussian processes that satisfy this assumption, see [4] for the details. Say that a continuous function $f: \mathbb{R} \to \mathbb{R}$ is (Bohr)-almost periodic if for each $\varepsilon > 0$ there exists $T_0(\varepsilon) > 0$ such that every interval of length $T_0(\varepsilon)$ contains a number τ with the following property: $|f(t+\tau) - f(t)| < \varepsilon$ for each $t \in \mathbb{R}$. Bochner came up with the following characterisation: f is almost periodic iff for every sequence $t_n \to \infty$, there is a subsequence $t_{n'}$ such that

$$||f - f(t_{n'} + \cdot)|| \to 0.$$

Such functions are equivalently characterised as uniform limits of trigonometric polynomials on the real line.

Proposition 7. Let $a = (a_i)$ be a summable sequence of non-negative numbers, $\varphi = (\varphi_i)$ a sequence of real numbers, and

$$C(x) = \sum_{i} a_{i} \cos(\varphi_{i} x). \tag{3}$$

Then C is almost periodic, and it is the reduced covariance $C=C_X$ of a Gaussian stationary process X which is a.s. bounded, almost periodic, and whose zero set Z_X is completely predictable.

Proof. C is obviously the uniform limit of trigonometric polynomials. Then X can be represented by

$$X(x) = \sum_{i} a_{i}(\alpha_{i} \cos(\varphi_{i} x) + \beta_{i} \sin(\varphi_{i} x))$$

where the α_i, β_i form a sequence of independent standard Gaussian variables, hence X is a.s. periodic for the same reasons as C. Finally, Z_X is completely predictable because C(x) gets arbitrarily close to 1 as $x \to \infty$, hence the previous proposition applies. \square

Even if such processes could be useful for modelling pseudo-crystalline arrangments, these periodicity properties prevent them from having the slightest form of asymptotic independence, ergodicity, let alone mixing properties. We shall give in the next section examples of fields whose zero set is weakly mixing and completely predictible.

The excursion volume of almost periodic fields have interesting properties. It is proved in the forthcoming paper [18] that $\mathcal{V}_{\mathsf{X}}(T) = \mathbf{Var}(|[-T,T] \cap \{\mathsf{X}>0\}|)$ strongly depends on the diophantine properties of the φ_i , where $|\cdot|$ denotes Lebesgue measure. If for instance $C_{\mathsf{X}}(x) = \cos(x) + \cos(\varphi x)$, $\mathcal{V}_{\mathsf{X}}(T)$ can either be bounded, or grow as an arbitrarily high power of T, depending on weather φ has "good" or "bad" approximations by rational numbers.

3 Symmetric Bernoulli convolutions

Let $\lambda = (\lambda_k)$ be a square summable family of \mathbb{R}_+ , and

$$C^{\lambda}(t) = \prod_{k=1}^{\infty} \cos(\lambda_k t), t \in \mathbb{R}.$$

This function converges for every t ([19, (3.7.9)]), and it is the reduced covariance of some stationary Gaussian field denoted by X^{λ} . The order of multiplication does not matter in the value $C^{\lambda}(t)$, hence we can assume throughout that $\lambda_k \geq \lambda_{k+1}$. By studying the uniform convergence on every compact, it is clear that X^{λ} is analytic, with an analytic continuation to \mathbb{C} .

The spectral measure μ^{λ} of X^{λ} is called a *symmetric Bernoulli convolution* because it is the law of the random variable

$$Y^{\lambda} := \sum_{k=1}^{\infty} \lambda_k \varepsilon_k$$

where the ε_k are iid Rademacher variables. Infinite products of characteristic functions have been studied by Lukacs [19, Chapter 3.7], where basic results stated here are proved. The main focus of study in the literature is the nature of μ^{λ} . It is known that μ^{λ} is pure, i.e. either continuous, purely discrete, or purely singular [19, Th. 3.7.7]. Bernoulli convolutions with geometric progression are a very active research subject, see for instance [24], they are one of the most studied examples of self-similar measures, and they are objects of great interest in fractal geometry.

Regarding ergodicity, let $\{X(t); t \in \mathbb{R}\}$ be a Gaussian stationary process with almost surely continuous paths, and $\mathcal{B} = \sigma(X)$ the σ -algebra generated by the random variables $X(t), t \in \mathbb{R}$. For $B \in \mathcal{B}, t \in \mathbb{R}$, let B + t be the event B translated by t. The process X is called weakly mixing if

$$\lim_{T \to \infty} T^{-1} \int_0^T |\mathbf{P}(A \cap (B+t)) - \mathbf{P}(A)\mathbf{P}(B)| dt = 0.$$

Proposition 8. Assume $\lambda_n > R_n > 0$ for every n. Then μ^{λ} has no atoms, its support has Cantor type (compact, completely disconnected and without isolated points), and

$$\lim_{T \to \infty} \frac{1}{T} \int_0^T C^{\lambda}(t) dt = 0,$$

hence X^{λ} and its zero set Z^{λ} are weakly mixing.

Proof. The support of μ is

$$\Sigma = \{ \sum_k \lambda_k \varepsilon_k : \varepsilon = (\varepsilon_k) \in \{-1, 1\}^{\mathbb{N}} \}.$$

Hence it is clear that given $t = \sum_{k} \lambda_k \varepsilon_k \in \Sigma, \eta = \lambda_n$ for some $n \in \mathbb{N}$,

$$(\Sigma \setminus \{t\}) \cap [t - \eta, t + \eta] \supset \{\sum_{k=1}^{n} \lambda_k \varepsilon_k + \sum_{k>n} \lambda_k \varepsilon'_{k-n} : \varepsilon' \in \{-1, 1\}^{\mathbb{N}}\}$$

is infinite. It follows that t is an accumulation point of Σ . To prove that Σ is completely disconnected, let $t = \sum_k \varepsilon_k \lambda_k$, $s = \sum_k \varepsilon_k' \lambda_k \in \Sigma$ with t < s, and $k_0 = \min\{k : \varepsilon_k \neq \varepsilon_k'\}$. Then let

$$t_{+} = \sum_{k < k_{0}} \varepsilon_{k} \lambda_{k} + \sum_{k \geq k_{0}} \lambda_{k}, \ s_{-} = \sum_{k < k_{0}} \lambda_{k} \varepsilon_{k} - \sum_{k \geq k_{0}} \lambda_{k}.$$

We have $t \leq t_+ < s_- \leq s$, where the strict inequality comes from $R_{k_0} < \lambda_{k_0}$. Noticing that there is no point of Σ in (t_+, s_-) , t and s are not connected through Σ .

To prove that there is no atom, symmetry considerations yield that $\mu(\{t\}) = \mu(\{s\})$ for every $t, s \in \Sigma$; since μ is finite, it follows that $\mu(t) = 0$ for every $t \in \Sigma$. Hence according to Maruyama's theorem [20], X is ergodic, and according to Zak & Rosiński [22], it is hence weakly mixing.

Since X is a.s. continuous, for $I \in \mathcal{B}(\mathbb{R})$ open, we have a.s.

$$\#\mathsf{Z}_{\mathsf{X}}\cap I = \lim_{n\geqslant 1} \sum_{k=-\infty}^{\infty} \mathbf{1}_{\{\mathsf{X}(k/n)\mathsf{X}((k+1)/n)\leqslant 0\}},$$

hence any event $A \in \sigma(\mathsf{Z}_\mathsf{X})$ is $\sigma(\mathsf{X})$ -measurable. It readily follows that since X is weakly mixing, so is Z_X .

There are also many cases where C^{λ} has stronger mixing properties. Most available results concern coefficients of the form $\lambda_k = a^k$ for some $a \in (0,1)$. The case a = 1/2, i.e. $\lambda_k = 2^{-k}$, actually yields the *sine process* thanks to Vieta's formula:

$$C(t) = \frac{\sin(t)}{t} = \frac{2\sin(t/2)\cos(t/2)}{t} = \frac{4\sin(t/4)\cos(t/4)\cos(t/2)}{t}$$

$$= \dots$$

$$= \lim_{n} \frac{2^{n}\sin(2^{-n}t)}{t} \prod_{k} \cos(2^{-k}t) = C^{\lambda}(t).$$

This example is central, as it is the only a for which $R_n = \lambda_n$, and it is the smallest a for which $\mu^{\lambda}(dx) = \mathbf{1}_{\{[-1,1]\}}(x)dx$ is continuous. Since $C^{\lambda}(t) \to 0$ as $t \to \infty$, X^{λ} and its zero set Z^{λ} are (strongly) mixing. For a < 1/2, the previous considerations yield that μ^{λ} is Cantor-like with no atoms, hence weakly mixing. If a is not of the form m^{-k} for some $m \in \mathbb{N}$, $C^{\lambda}(t) = O(|\log(t)|^{-\gamma})$ for some $\gamma > 0$ as $t \to \infty$, whereas for $a = m^{-k}, m \geqslant 3$, $\limsup_{t \to \infty} C^{\lambda}(t) > 0$, see [19].

There are many other examples of such products, with diverse behaviours. If for instance $\lambda_k = k^{-1}$,

$$C^{\lambda}(t) = \prod_k \cos\left(\frac{t}{k}\right) = \prod_n \operatorname{sinc}\left(\frac{t}{2n+1}\right)$$

decreases faster than any polynomial. Let us give now an example of symmetric Bernoulli convolution which has rigid zeros. As a further proof of asymptotic independence, we show that this process is a.s. unbounded, at the contrary of almost periodic Gaussian fields mentioned in the previous section.

Proposition 9. Assume $\lambda_k = k!^{-1}\pi$. Then $(-1)^n C^{\lambda}(n!) \to 1$ as $n \to \infty$. In particular, the zero set Z of the Gaussian process X with reduced covariance C^{λ} is completely predictable, hence super rigid. Furthermore, μ_{λ} is purely singular without atoms, hence X and Z are weakly mixing. Also, $\sup_{t \in \mathbb{R}} |\mathsf{X}_{\lambda}(t)| = \infty$ a.s.

Proof. Let $t_n = n!$. We have

$$\prod_{k=1}^{n} \cos(t_n \lambda_k) = \prod_{k=1}^{n} \cos(\pi n! / k!) = (-1)^n.$$

For $k \ge n+1$,

$$\cos\left(\frac{\pi n!}{k!}\right) \geqslant 1 - \left(\frac{\pi n!}{k!}\right)^2 \geqslant 1 - \left(\frac{\pi^2}{k^2}\right)$$

hence the first result is proved with

$$1 \geqslant (-1)^n C(t_n) \geqslant \prod_{k=n+1}^{\infty} \left(1 - \left(\frac{\pi}{k}\right)^2\right) \geqslant 1 - c/n.$$

To prove the unboundedness of X^{λ} , consider the pseudo-metric on \mathbb{R}

$$d(t,s) = \sqrt{1 - C(t-s)}.$$

We will prove that (\mathbb{R}, d) cannot be covered by finitely many d-balls of some radius $\varepsilon > 0$. Let $t \in B_d(0, \varepsilon)$, and for $k \in \mathbb{Z}$, let $n_k = n_k(t)$ be the integer such that for some $\varepsilon_k = \varepsilon_k(t) \in [-1/2, 1/2)$

$$k!^{-1}t = n_k + \varepsilon_k.$$

Then for some c > 0

$$d(t,0)^2 \geqslant 1 - \exp(\sum_k \ln(|\cos(\varepsilon_k \pi))|) \geqslant c \sum_k \varepsilon_k^2$$

hence $c\sum_k \varepsilon_k^2 < \varepsilon^2$, and in particular each ε_k is smaller than $\delta := \sqrt{c^{-1}}\varepsilon$. Up to diminishing ε , assume $\delta < 1/4$. Let us now cover \mathbb{R}_+ by the sets

$$A_{m,k} = \{t = mk! + k! \varepsilon_k : \varepsilon_k \in [-1/2, 1/2]\}$$

for $k \in \mathbb{N}, 0 < m \le k$. We have for $t \in A_{m,k}, 0 \le l \le k$, with $k^{(l)} = k(k-1) \dots (k-l+1)$,

$$\frac{t}{(k-l)!} = mk^{(l)} + k^{(l)}\varepsilon_k = n_{k-l}(t) + \varepsilon_{k-l}(t),$$

hence $\varepsilon_k \in C_{k^{(l)}}$, where

$$C_K = \{t \in [-1/2, 1/2] : d(Kt, \mathbb{Z}) \leq \delta\}, K > 0.$$

Let us then define $E_0 = [-\delta, \delta]$ and for $0 \le l < k$,

$$E_{l+1} = E_l \cap C_{k^{(l+1)}}.$$

Let n_l be the number of maximal segments forming E_l . We prove by induction that $n_l \leq (2\delta)(4\delta)^{l-1}k^{(l)}$ for $k-l \geq \delta^{-1}$, i.e. $l \leq l^* := [k-\delta^{-1}]$. Indeed, through the intersection with $C_{k^{(l)}}$, E_l is made up of segments of length at most $a_l := 2\delta k^{-(l)}$, and each segment of length a_l is intersected by at most

$$a_l k^{(l+1)} + 2 = 2\delta(k-l) + 2 \le 4\delta(k-l)$$

segments of $C_{k(l+1)}$. It follows that

$$n_{l+1} \leqslant n_l 4\delta(k-l) \leqslant 2\delta(4\delta)^l k^{(l)}(k-l);$$

which proves the induction. Hence the Lebesgue measure of E_l for $l \leq l^*$ is smaller than $2\delta(4\delta)^l k^{(l)} 2\delta k^{-(l+1)} \leq (4\delta)^{l+1}$. We have for $l \geq l^*$, $E_l \subset E_{l^*}$. It follows that

$$|B_d(0,\varepsilon)\cap A_{m,k}|\leqslant |E_k|\leqslant |E_{l^*}|\leqslant (4\delta)^{l^*+1}.$$

It follows that for $\delta < 1/4$

$$|B_d(0,\sqrt{\varepsilon})| \leqslant \sum_k \sum_m (4\delta)^{k-\delta^{-1}-1} < \infty.$$

By stationarity, $|B_d(t,\varepsilon)| = |B_d(0,\varepsilon)| < \infty$ for $t \in \mathbb{R}$, hence \mathbb{R} cannot be covered by a finite number of such balls, hence the sample paths of X^λ are a.s. unbounded (see for instance [1, Th. 1.19]).

4 Variance lower bound for the number of zeros

In this section we take X as a stationary Gaussian process whose reduced covariance function C satisfies Geman's condition (1). In particular, C' and C''(0) exist, hence Fatou's lemma yields that the spectral measure μ has a finite second moment, and

$$C''(t) = -\int t^2 e^{itx} \mu(dx). \tag{4}$$

Recall that C is normalised so that C''(0) = -1. If Geman's condition is not satisfied, the variance is infinite and the result of the theorem holds true. We establish the following proposition, which readily yields (iii) in Theorem 2.

Proposition 10. We have for every T > 0

$$\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \frac{\sin(1)^2}{\pi} T^2 \int (1 - xy)^2 \mathbf{1}_{\{|x - y| < T^{-1}\}} \mu(dx) \mu(dy).$$

Let us now prove both this proposition and Theorem 2-(i),(ii). The method of Kratz and Léon is to compute the chaotic decomposition in X, X':

$$N_{\mathsf{X}}(T) = \mathbf{E}(N_{\mathsf{X}}(T)) + \sum_{q=1}^{\infty} N_{T,q}$$

where

$$N_{T,q} = \sum_{k=0}^{q} a_k d_{q-k} \int_0^T H_k(X(t)) H_{q-k}(X'(t)) dt$$

where $H_m, m \in \mathbb{N}$ is the m-th Hermite polynomial, and $a_m, d_m, m \in \mathbb{N}$ are coefficients that vanish for m odd. In particular,

$$a_0 = \frac{1}{\sqrt{2\pi}}, a_2 = \frac{1}{2\sqrt{2\pi}}$$

 $d_0 = 1, d_2 = -\frac{1}{2}$

and $N_{T,q}$ vanishes for odd q. It is proved in [1] that this decomposition is valid as long as $C(t) \neq \pm 1$ for $t \neq 0$.

If $C(\tau) = \pm 1$ for some $\tau > 0$, then X is 2τ -periodic a.s., and we have

$$N_{\mathsf{X}}(T) = N_{\mathsf{X}}(2\tau) |T/2\tau| + N_T'.$$

where $\mathbf{E}(N_T'^2)$ is bounded by $\mathbf{E}(N_{2\tau}^2)$. Also, $\mathbf{Var}(N_{\mathsf{X}}(2\tau)) = 0$ iff $C(t) = \cos(k\pi t/\tau)$ for some $k \in \mathbb{N}$. It follows that under the convention C''(0) = -1, $\mathbf{Var}(N_{\mathsf{X}}(T))$ is bounded iff $C(t) = \cos(t)$, and otherwise it is proportionnal to T^2 , hence the result of Theorem 2 holds true in this case.

Going back to the general case,

$$N_{T,2} = -\frac{1}{2\sqrt{2\pi}} \left[\int_{-T}^{T} H_0(\mathsf{X}(t)) H_2(\mathsf{X}'(t)) dt - \int_{-T}^{T} H_2(\mathsf{X}(t)) H_0(\mathsf{X}'(t)) dt \right].$$

Then, with $H_0 = 1$, $H_2(x) = x^2 - 1$, since for two standard Gaussian variables α, β with correlation ρ , $Cov(\alpha^2, \beta^2) = 2\rho^2$,

$$\begin{aligned} \mathbf{Var}(N_{\mathsf{X}}(T)) &= \sum_{q=1}^{\infty} \mathbf{Var}(N_{T,q}) \\ &\geqslant \mathbf{Var}(N_{T,2}) \\ &= \frac{1}{4\pi} \int_{-T}^{T} \int_{-T}^{T} \left[\mathrm{Cov}(\mathsf{X}(t)^2, \mathsf{X}(s)^2) + \mathrm{Cov}(\mathsf{X}'(t)^2, \mathsf{X}'(s)^2) - 2\mathrm{Cov}(\mathsf{X}(t)^2, \mathsf{X}'(s)^2) dt ds \right] dt ds. \end{aligned}$$

We put $\gamma_T(z) = \mathbf{1}_{\{|z| < T\}}$, hence with z = t - s, w = -s,

$$\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \frac{1}{4\pi} \int_{\mathbb{R}^{2}} [C(z)^{2} + C''(z)^{2} - 2C'(z)^{2}] \left[\int_{\mathbb{R}} \gamma_{T}(z - \omega) \gamma_{T}(\omega) dw \right] dz$$
$$= \frac{1}{4\pi} \int_{\mathbb{R}^{2}} [C(z)^{2} + C''(z)^{2} - 2C'(z)^{2}] \gamma_{T}^{\star 2}(z) dz.$$

We recall that for T > 0, $x \in \mathbb{R}$,

$$\widehat{\gamma_T}(x) := \int \gamma_T(x) e^{ixt} dt = T\widehat{\gamma_1}(Tt) = 2T \operatorname{sinc}(Tt).$$

By (4),

$$\begin{split} \int_{\mathbb{R}} C(z)^2 \gamma_T^{\star 2}(z) &= \int_{\mathbb{R}} \gamma_T^{\star 2}(z) \int_{\mathbb{R}^2} e^{ixz} e^{iyz} \mu(dx) \mu(dy) \\ &= 4T^2 \int_{\mathbb{R}} \mathrm{sinc}(T(x+y))^2 \mu(dx) \mu(dy) \\ \int_{\mathbb{R}} C'(z)^2 \gamma_T^{\star 2}(z) &= \int_{\mathbb{R}} \gamma_T^{\star 2}(z) \int_{\mathbb{R}^2} (ix) (iy) e^{ixz} e^{iyz} \mu(dx) \mu(dy) \\ &= -4T^2 \int_{\mathbb{R}} xy \, \mathrm{sinc}(T(x+y))^2 \mu(dx) \mu(dy) \\ \int_{\mathbb{R}} C''(z)^2 \gamma_T^{\star 2}(z) &= \int_{\mathbb{R}} \gamma_T^{\star 2}(z) \int_{\mathbb{R}^2} x^2 y^2 e^{ixz} e^{iyz} \mu(dx) \mu(dy) \\ &= 4T^2 \int_{\mathbb{R}} x^2 y^2 \, \mathrm{sinc}(T(x+y))^2 \mu(dx) \mu(dy) \end{split}$$

whence finally, after changing y to -y,

$$\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \frac{T^2}{\pi} \int_{\mathbb{R}^2} (1 - xy)^2 \mathrm{sinc}^2(T(x - y)) \mu(dx) \mu(dy).$$

A straightforward study shows that for $x \in \mathbb{R}$

$$\operatorname{sinc}(Tx)^2 \geqslant \sin(1)^2 \mathbf{1}_{\{|Tx|<1\}} \geqslant \sin(1)^2 \Delta_T(x)$$

where $\Delta_T(x) = \mathbf{1}_{\{|Tx|<1\}}(1-|Tx|)$. The first inequality proves Proposition 10.

Let now $\varepsilon > 0, A = \mathbb{R} \setminus ([1 - \varepsilon, 1 + \varepsilon] \cup [-1 - \varepsilon, -1 + \varepsilon]), \mu_{\varepsilon} = \mu \mathbf{1}_{\{A\}}$ satisfying the following: $\mu_{\varepsilon} \neq 0$ (recall that μ was assumed to not be concentrated on $\{\pm 1\}$), and if μ does not have a density on $\mathbb{R} \setminus \{-1, 1\}$, μ_{ε} does not have a density either. Also define

$$C_{\varepsilon}(t) = \int_{\mathbb{R}} e^{ixt} \mu_{\varepsilon}(dx).$$

For $T > 2\varepsilon^{-1}$, we have for some $c_{\varepsilon} > 0$

$$(1-xy)^2 \operatorname{sinc}(T(x-y))^2 \mathbf{1}_{\{x,y\notin[\pm 1\pm\varepsilon]\}} \geqslant \sin(1)^2 (1-xy)^2 \Delta_T(x-y) \mathbf{1}_{\{x,y\notin[\pm 1\pm\varepsilon]\}}$$
$$\geqslant c_{\varepsilon} \Delta_T(x-y) \mathbf{1}_{\{x,y\notin[\pm 1\pm\varepsilon]\}}.$$

Then the Fourier inversion formula yields for $t \in \mathbb{R}$

$$\Delta_T(x/2) = \frac{T}{2} \gamma_{T^{-1}} \star \gamma_{T^{-1}}(x) = \frac{T}{2\sqrt{2\pi}} \int \widehat{\gamma_{T^{-1}}}(t)^2 e^{ixt} dt = \frac{2T^{-1}}{\sqrt{2\pi}} \int \operatorname{sinc}(T^{-1}t)^2 e^{itx} dt,$$

hence for some c' > 0

$$\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \frac{c_{\varepsilon}T^{2}}{\pi} \iint \Delta_{T}(x+y)\mu_{\varepsilon}(dx)\mu_{\varepsilon}(dy)$$

$$= c'T \iiint e^{2ixt}e^{2iyt}\operatorname{sinc}(T^{-1}t)^{2}\mu_{\varepsilon}(dx)\mu_{\varepsilon}(dy)$$

$$= c'T \int_{\mathbb{R}}\operatorname{sinc}^{2}(T^{-1}t)C_{\varepsilon}(2t)^{2}dt$$

$$\geqslant c'\sin(1)^{2}T \int_{\mathbb{R}}\mathbf{1}_{\{|t|< T\}}C_{\varepsilon}(2t)^{2}dt \qquad (5)$$

and we have used the finiteness of μ_{ε} and the integrability of sinc². Hence Theorem 2-(i) is proved.

In view of proving (ii), we assume that the variance is linear, hence we need to show that $C'' - C \in L^2(\mathbb{R})$. By (5), if the variance is linear, C_{ε} is L^2 , hence μ_{ε} has a L^2 density, called f. Proposition 10 yields

$$\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \frac{\sin(1)^{2}T^{2}}{\pi} \iint (1 - xy)^{2} \mathbf{1}_{\{|x - y| < T^{-1}\}} f(x) f(y) dx dy$$

$$\liminf_{T} T^{-1} \mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \frac{\sin(1)^{2}}{\pi} \int_{\mathbb{R} \setminus \{-1,1\}} \left(\liminf_{T} T \int_{x - 1/T}^{x + 1/T} (1 - xy)^{2} f(y) dy \right) f(x) dx$$

$$= \frac{\sin(1)^{2}}{\pi} \int (1 - x^{2})^{2} f(x)^{2} dx$$

$$= \frac{\sin(1)^{2}}{\pi} \int (C'' - C)^{2}$$

by Parseval's identity, hence $C'' - C \in L^2$.

To conclude the proof of (ii), let us assume now that μ has a L^2 density f on a neighbourhood I of $\{-1,1\}$, and show that (2) is equivalent to $C''-C\in L^2$. If μ does not have a density on all \mathbb{R} , then $C\notin L^2$, otherwise \widehat{C} would be a L^2 density of μ . Then, if $C''-C\in L^2$, Fourier's inversion formula yields

$$C''(t) - C(t) = \frac{1}{\sqrt{2\pi}} \widehat{\int (C'' - C)(x)} e^{itx} dx = \int (1 - x^2)^2 e^{itx} \mu(dx),$$

hence μ has density $\widehat{C''-C}(x)(1-x^2)^{-2}$, which leads to a contradiction. Hence none of the conditions hold if μ does not have a density.

Let us assume finally that μ 's density f can be extended to all \mathbb{R} . Then

$$C(x) = \int e^{ixt} f(x) dx, \quad C''(x) = -\int x^2 f dx, \quad (C'' - C)(x) = \int (1 - x^2) f(x) dx.$$

The integrands are square intagrable on I, so we need to prove the equivalence on $B := \mathbb{R} \setminus I$. We have

$$C \in L^2(B) \Leftrightarrow f \in L^2(B), C'' \in L^2(B) \Leftrightarrow \int_B x^4 f(x)^2 dx < \infty, \quad C'' - C \in L^2(B) \Leftrightarrow \int_B (1 - x^2)^2 f^2(x) dx < \infty,$$

hence the equivalence stems from the existence of c>0 such that for $x\notin I$

$$c(1+x^4) \le (1-x^2)^2 \le 1+x^4$$
.

Remark 11. One might think that going for higher order chaoses could help replace (ii) by the sufficient condition $C, C'' \in L^2$, but the gap present at ± 1 is necessarily still present at higher order chaoses because the zeros of $C(t) = \cos(t)$ satisfy the same chaotic decomposition (see [16, Prop. 2.2]), hence the variance of every chaos is bounded in T. It remains open to find a necessary and sufficient condition for linear variance.

Proposition 12. If

$$C(t) = \prod_{k} \cos(t/k!),$$

then for all $\varepsilon > 0$,

$$T^{\varepsilon-2}\mathbf{Var}(N_{\mathsf{X}}(T))\to\infty$$

as $T \to \infty$.

Proof. Let

$$R_N = \sum_{k>N} k!^{-1}.$$

Recall that the spectral measure μ is the law of $Y = \lim_N Y_N$ where $Y_N = \sum_{k=1}^N \frac{\varepsilon_k}{k!}$, for random iid Rademacher variables (ε_k) . In particular, $|Y - Y_N| < R_N$ a.s. Two consecutive points of the support of Y_N have distance $2/n! > 2R_N$ (proved by induction) and they all have mass 2^{-N} . Let $Y', Y_N', N \in \mathbb{N}$ be iid copies of Y, Y_N . Since Y's support is in (-1, 1), there is c > 0 such that for T sufficiently large, using Proposition 10

$$\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant cT^2 \mathbf{P}(|Y - Y'| < 1/T).$$

Let $N = N_T$ be the integer such that $R_{N-1} \geqslant \frac{1}{4T} > R_N$.

$$\mathbf{P}(|Y - Y'| < 1/T) \geqslant \mathbf{P}(|Y_N - Y'_N| < 2R_N)$$

$$= \sum_{x \in \text{Supp}(Y_N)} \mathbf{P}(Y_N = Y'_N = x)$$

$$= 2^{-N}.$$

Then we have by convexity

$$\ln(4T) \geqslant \ln(R_{N-1}^{-1}) > \ln((N-1)!) = \sum_{k=1}^{N-1} \ln(k) > (N-1)\ln((N-1)/2).$$

It follows that

$$(4T)^{\varepsilon-2}\mathbf{Var}(N_{\mathsf{X}}(T)) \geqslant \exp[\varepsilon(N-1)\ln((N-1)/2))]2^{-N} \to \infty.$$

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