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Asymptotic variance of stationary reversible and normal Markov processes

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Abstract

We obtain necessary and sufficient conditions for the regular variation of the variance of partial sums of functionals of discrete and continuous-time stationary Markov processes with normal transition operators. We also construct a class of Metropolis-Hastings algorithms which satisfy a central limit theorem and invariance principle when the variance is not linear in n.

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1 Introduction

Let $(\xi_n)_{n\in\mathbb{Z}}$ be a stationary Markov chain defined on a probability space $(\Omega,\mathcal{F},\mathbb{P})$ with values in a general state space (S,\mathcal{A}) and let the marginal distribution be denoted by $\pi(A)=\mathbb{P}(\xi_0\in A)$. We assume that there is a regular conditional distribution denoted by $Q(x,A)=\mathbb{P}(\xi_1\in A|\xi_0=x)$. Let Q also denote the Markov transition operator acting via $(Qg)(x)=\int_S g(s)Q(x,\mathrm{d}s)$, on $\mathbb{L}^2_0(\pi)$, the set of measurable functions on S such that $\int_S g^2(s)\pi(\mathrm{d}s)<\infty$ and $\int_S g(s)\pi(\mathrm{d}s)=0$. If $g,h\in\mathbb{L}^2_0(\pi)$, the integral $\int_S g(s)h(s)\pi(\mathrm{d}s)$ will sometimes be denoted by $\langle g,h\rangle$.

For some function $g \in \mathbb{L}_0^2(\pi)$, let

$$X_i = g(\xi_i), \quad S_n(X) = \sum_{i=1}^n X_i, \quad \sigma_n(g) = (\mathbb{E}S_n^2(X))^{1/2}.$$
 (1.1)

Denote by \mathcal{F}_k the σ -field generated by ξ_i with $i \leq k$.

For any integrable random variable X we denote by $\mathbb{E}_k X = \mathbb{E}(X|\mathcal{F}_k)$. With this notation, $\mathbb{E}_0 X_1 = Qg(\xi_0) = \mathbb{E}(X_1|\xi_0)$. We denote by $\|X\|_p$ the norm in $\mathbb{L}^p(\Omega, \mathcal{F}, \mathbb{P})$.

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The Markov chain is called *normal* when the transition operator Q is normal, that is it commutes with its adjoint Q^* , namely $QQ^* = Q^*Q$.

From the spectral theory of normal operators on Hilbert spaces (see for instance [34]), it is well known that for every $g \in L_0^2(\pi)$ there is a unique transition spectral measure ν supported on the spectrum of the operator $D := \{z \in \mathbb{C} : |z| \leq 1\}$, such that

$$cov(X_0, X_n) = cov((g(\xi_0), Q^n g(\xi_0))) = \langle g, Q^n g \rangle = \int_D z^n \nu(\mathrm{d}z). \tag{1.2}$$

and

$$cov(\mathbb{E}_0(X_i), \mathbb{E}_0(X_j)) = \langle Q^i g, Q^j g \rangle = \langle g, Q^i (Q^*)^j g \rangle = \int_{\mathbb{D}} z^i \bar{z}^j \nu(\mathrm{d}z).$$

In particular, the Markov chain is reversible if $Q=Q^*$. The condition of reversibility is equivalent to requiring that (ξ_0,ξ_1) and (ξ_1,ξ_0) have the same distribution. Furthermore, in the reversible case ν is concentrated on [-1,1].

Limit theorems for additive functionals of reversible Markov chains have received considerable attention in the literature not only for their intrinsic interest, but also for their great array of applications which range from interacting particle systems (see the seminal paper by Kipnis and Varadhan [22]) and random walks in random environments (see for example [37]), to the relatively recent applications in computational statistics with the advent of Markov Chain Monte Carlo algorithms (e.g. [18, 33]). Limit theorems have appeared under a great array of conditions, notably geometric ergodicity (for an overview see [24]), conditions on the growth of the conditional expectations $\mathbb{E}(S_n|X_1)$ (see e.g. [27], [30]), under spectral conditions (see [22, 15, 10]) or under conditions on the resolvent of the transition operator (see [37]), a method which is also applicable in the non-normal case where spectral calculus may not apply.

The variance of the partial sums plays a major role in limit theorems where it acts as a normalizer, and also in computational statistics where the asymptotic variance is used as a measure of the efficiency of an algorithm (see e.g. [36]). It is not surprising then, that in certain cases, conditions for the central limit theorem have been imposed directly on the growth of the variance. In fact in 1986 Kipnis and Varadhan [22] proved the functional form of the central limit theorem for functionals of stationary reversible ergodic Markov chains under the assumption that

$$\lim_{n \to \infty} \frac{\operatorname{var}(S_n)}{n} = \sigma_g^2,\tag{1.3}$$

and further established necessary and sufficient conditions for the variance of the partial sums to behave linearly in n in terms of the transition spectral measure ν . In particular they showed that for any reversible ergodic Markov chain the convergence in (1.3) is sufficient for the functional central limit theorem $S_{[nt]}/\sqrt{n} \Rightarrow |\sigma_g|W(t)$ (where W(t) is the standard Brownian motion, \Rightarrow denotes weak convergence and [x] denotes the integer part of x). Moreover, (1.3) is equivalent to the fact that the finite limiting variance is then given by

$$\sigma_g^2 = \int_{-1}^1 \frac{1+t}{1-t} \nu(\mathrm{d}t) < \infty. \tag{1.4}$$

Furthermore, according to Remark 4 on page 514 in [10] if in addition to (1.4) we assume $\rho(-1)=0$ then, we also have

$$\lim_{n \to \infty} \sum_{i=0}^{n} \text{cov}(X_0, X_i) = \int_{-1}^{1} \frac{1}{1-t} \nu(dt).$$

See also [18] for a discussion of when (1.3) and (1.4) are equivalent.

It is remarkable that in the reversible case, conditions (1.3) and (1.4) are equivalent, both sufficient for the central limit theorem and invariance principle, and conjectured to be sufficient for the almost sure conditional central limit theorem. It is an open problem whether (1.4) is also necessary for the central limit theorem with normalisation \sqrt{n} , even in the reversible case.

On the other hand, notice that any invertible transformation T generates a unitary, and thus normal, transition operator Qf(x)=f(T(x)), since $Q^*f(x)=f(T^{-1}(x))$ whence $QQ^*=Q^*Q=I$ is the identity operator. In particular, any stationary sequence ξ_i , can be treated as a functional of a normal Markov chain. Therefore for normal, non-reversible Markov chains, (1.3) and the central limit theorem and invariance principle are no longer equivalent without further assumptions (see e.g. Bradley [5] and Giraudo and Volný [14] for counterexamples).

For the non-reversible case, Gordin and Lifšic [15] applied martingale methods and stated, among other results, the central limit theorem for functionals of stationary ergodic Markov chains with normal transition operator, under the spectral condition

$$\int_{D} \frac{1}{|1-z|} \nu(\mathrm{d}z) < \infty. \tag{1.5}$$

If condition (1.5) holds then (1.3) also holds with

$$\sigma^2 := \int_D \frac{1 - |z|^2}{|1 - z|^2} \nu(\mathrm{d}z) < \infty. \tag{1.6}$$

One of our main results, Theorem 2.6, gives necessary and sufficient conditions for the existence of the limit $\text{var}(S_n)/n \to K < \infty$. We shall see that $\text{var}(S_n)/n \to K$ if and only if $\sigma^2 < \infty$ and $\nu(U_x)/x \to C$ as $x \to 0^+$, where

$$U_x = \{z = (1 - r)e^{iu} : |z| \le 1, \ 0 \le r \le |u| \le x\}.$$

In this case $K = \sigma^2 + \pi C$. Furthermore if (1.5) holds then C = 0.

Recently Zhao et al. [39] and Longla et al. [26], in the context of reversible Markov chains, studied the asymptotic behavior of S_n for the more general case when the variance of partial sums behaves as a regularly varying function $\sigma_n^2 = \text{var}(S_n) = nh(n)$ where h(x) is slowly varying, i.e. $h: (0,\infty) \to (0,\infty)$, continuous, and $h(st)/h(t) \to 1$ for all s>0. For this case the situation is different and in [39, 26] examples are given of stationary, reversible, and ergodic Markov chains that satisfy the CLT under a normalization different of σ_n , namely $S_n/\sigma_n \Rightarrow N(0,c^2)$ for a $c \in (0,1)$. Of course by a version of Fatou's lemma (see Theorem 3.4 [2]), the limiting variance is dominated by $\lim \inf \mathbb{E} S_n^2/\sigma_n^2 = 1$.

On the other hand, in a recent paper, Deligiannidis and Utev [8] have studied the relationship between the variance of the partial sums of weakly stationary processes and the spectral measure induced by the unitary shift operator. To be more precise, by the Birghoff-Herglotz Theorem (see e.g. Brockwell and Davis [6]), there exists a unique measure on the unit circle, or equivalently a non-decreasing function F, called the spectral distribution function on $[0,2\pi]$, such that

$$\operatorname{cov}(X_0, X_n) = \int_0^{2\pi} e^{\mathrm{i}n\theta} F(\mathrm{d}\theta), \quad \text{for all } n \in \mathbb{Z}.$$
 (1.7)

If F is absolutely continuous with respect to the normalized Lebesgue measure λ on $[0,2\pi]$, then the Radon-Nikodym derivative f of F with respect to the Lebesgue measure is called the *spectral density*; in other words $F(\mathrm{d}\theta)=f(\theta)\mathrm{d}\theta$. The main result of [8] is given below. In the sequel, the notation $a_n \sim b_n$ as $n \to \infty$ means that $\lim_{n \to \infty} a_n/b_n = 1$.

Theorem A. [Deligiannidis and Utev [8]] Let $S_n := X_1 + \cdots + X_n$ where $(X_i)_{i \in \mathbb{Z}}$ is a real weakly stationary sequence. For $\alpha \in (0,2)$, define $C(\alpha) := \Gamma(1+\alpha)\sin(\frac{\alpha\pi}{2})/[\pi(2-\alpha)]$, and let h be slowly varying at infinity. Then $\operatorname{var}(S_n) \sim n^{\alpha}h(n)$ as $n \to \infty$ if and only if $F(x) \sim \frac{1}{2}C(\alpha)x^{2-\alpha}h(1/x)$ as $x \to 0$.

In this paper we obtain necessary and sufficient conditions for the regular variation of the variance of partial sums of functionals of stationary Markov chains with normal operators. The necessary and sufficient conditions are based on several different representations in terms of:

- 1. the spectral distribution function in the sense of the Birghoff-Herglotz theorem,
- 2. the transition spectral measure of the associated transition operator,
- 3. the harmonic measure of Brownian motion in the disk,
- 4. a martingale decomposition.

In the case of stationary reversible Markov Chains we also construct a class of Metropolis-Hastings algorithms with non-linear growth of variance, for which we establish the invariance principle and conditional central limit theorem with normalization $\sqrt{nh(n)}$.

Continuous-time processes. In the continuous time setting, let $\{\xi_t\}_{t\geq 0}$ be a stationary Markov process with values in the general state space (S, \mathcal{A}) , defined on a probability space (Ω, \mathcal{F}, P) , with stationary measure π . We assume that the contraction semigroup

$$T_t g(x) := \mathbb{E}[g(\xi_t)|\xi_0 = x], \quad g \in L^2(\pi), \quad t \ge 0,$$

is strongly continuous on $L^2(\pi)$, and we let $\{\mathcal{F}_t\}_{t\geq 0}$ be a filtration on (Ω,\mathcal{F},P) with respect to which $\{\xi_t\}_t$ is progressively measurable and satisfies the Markov property $\mathrm{E}(g(\xi_t)|\mathcal{F}_u)=T_{t-u}g(\xi_u)$, for any $g\in L^2(\pi)$ and $0\leq u< t$. Furthermore we can write $T_t=\mathrm{e}^{Lt}$, where L is the infinitesimal generator of the process $\{\xi_t\}_t$, and $\mathcal{D}(L)$ its domain in $L^2(\pi)$. We assume T_t to be normal, that is $T_t^*T=T_tT_t^*$, which then implies that L is a normal, possibly unbounded operator, with spectrum supported in the left half-plane $\{z\in\mathbb{C}:\Re(z)\leq 0\}$ (see [34, Theorem 13.38]). In the reversible case the spectrum of L is supported on the left real half-axis(see [22, Remark 1.7]).

Similarly to the discrete case, with any $f \in L^2(\mu)$ we can associate a unique spectral measure $\nu(\mathrm{d}z) = \nu_f(\mathrm{d}z)$ supported on the spectrum of L such that

$$\langle f, T_t f \rangle = \operatorname{cov}(f(\xi_0), f(\xi_t)) = \int_{\Re(z) \le 0} e^{zt} \nu(\mathrm{d}z).$$

In the reversible case Kipnis and Varadhan [22] proved an invariance principle under the condition that $f \in \mathcal{D}((-L)^{-1/2})$, which in spectral form is equivalent to

$$\int_{\alpha=-\infty}^{0} \frac{-1}{\alpha} \nu(\mathrm{d}\alpha) < \infty. \tag{1.8}$$

Building on the techniques in [15, 22], Holzmann [19, 20] established the central limit theorem for processes with normal transition semi-groups (see also [29]), under the condition

$$\int_{\Re(z)<0} |z|^{-1} \nu(\mathrm{d}z) < \infty. \tag{1.9}$$

In this case

$$\lim_{T \to \infty} \frac{\operatorname{var}(S_T)}{T} = -2 \int_{\mathbb{H}^-} \Re(1/z) \nu(\mathrm{d}z) =: \varsigma^2.$$

On the other hand, using resolvent calculus, Toth [38, 37] treated general discrete and continuous-time Markov processes and obtained a martingale approximation, central limit theorem and convergence of finite-dimensional distributions to those of Brownian motion, under conditions on the resolvent operator which may hold even in the non-normal case. Similar conditions, albeit in the normal case, also appeared later in [19, 20].

Under any of the above conditions, it is clear that the variance of S_T is asymptotically linear in T. Similarly to the discrete case, we show in Theorem 3.2, that $\text{var}(S_n)/n \to K = \varsigma^2 + \pi C$ if and only if $\varsigma^2 < \infty$ and $\nu(U_x)/x \to C$, where

$$U_x = \{a + ib : 0 \le -a \le |b| \le x\}.$$

The rest of the paper is structured as follows. We provide our results for discrete time processes in Section 2 and for continuous time in Section 3. Section 4 contains the proofs, while the Appendix contains two standard Tauberian theorems to make the paper self-contained, and technical lemmas used in Section 3.

2 Results for Markov chains

2.1 Relation between the transition spectral measure and spectral distribution function

Our first result gives a representation of the spectral distribution function in terms of the transition spectral measure. This link makes it possible to use the results in [8] to analyze the variance of partial sums. Quite remarkably, if the transition spectral measure is supported on the open unit disk, the spectral distribution function is absolutely continuous with spectral density given by (2.1), and in this case the sequence $cov(X_0,X_n)$ converges to 0.

Lemma 2.1 (Representation Lemma). Let $(\xi_n)_{n\in\mathbb{Z}}$ be a stationary Markov chain, with normal transition operator Q. Let $g\in L^2_0(\pi)$, $X_i:=g(\xi_i)$ and write $\nu=\nu_g$ for the operator spectral measure with respect to g. Also denote the unit circle $\Gamma:=\{z:|z|=1\}$ and by $D_0:=\{z:|z|<1\}$. Denote by ν_Γ the restriction of the measure ν to Γ and by ν_0 denote the restriction of the measure ν to D_0 . Then

$$\operatorname{cov}(X_0, X_n) = \int_0^{2\pi} e^{\mathrm{i}tn} [\nu_{\Gamma}(\mathrm{d}t) + f(t)\mathrm{d}t],$$

where

$$f(t) = \frac{1}{2\pi} \int_{D_0} \frac{1 - |z|^2}{|1 - ze^{it}|^2} \nu_0(dz).$$
 (2.1)

Furthermore the spectral distribution function has the representation

$$F(\mathrm{d}t) = \nu_{\Gamma}(\mathrm{d}t) + f(t)\mathrm{d}t. \tag{2.2}$$

Remark 2.2. By integrating relation (2.2) we obtain

$$F(x) = x \int_{D_0} T_x(z) \nu_0(\mathrm{d}z) + \nu_{\Gamma}([0, x])$$
 (2.3)

where

$$T_x(z) := \frac{1}{2\pi} (1 - |z|^2) \int_0^1 \frac{dt}{|1 - z e^{itx}|^2}.$$

By combining Representation Lemma 2.1 with Theorem A we obtain the following corollary.

Corollary 2.3. Let $(\xi_n)_{n\in\mathbb{Z}}$ be as in Lemma 2.1 and let $\alpha\in(0,2)$. Then $\mathrm{var}(S_n)=n^\alpha h(n)$ as $n\to\infty$ if and only if $F(x)=\frac{1}{2}C(\alpha)x^{2-\alpha}h(1/x)$ as $x\to0^+$.

It should be obvious from the statement of Representation Lemma 2.1 that Theorem A is directly applicable to the measure $\mathrm{d}F$. The conditions on F, mentioned in Corollary 2.3, when expressed in terms of the operator spectral measure, become technical conditions on the growth of integrals of the Poisson kernel over the unit disk. To get further insight into this lemma we shall apply it to reversible Markov chains.

Our next result, a corollary of Representation Lemma 2.1 combined with Theorem A, provides this link and points out a set of equivalent conditions for regular variation of the variance for reversible Markov chains. In fact, as it turns out, if the spectral measure has no atoms at ± 1 , it follows that the spectral distribution function is absolutely continuous and we obtain an expression for the spectral density. Related ideas, under more restrictive assumptions have appeared in [12], while in [10] a spectral density representation was obtained for positive self-adjoint transition operators, in other words ν supported on [0,1).

Corollary 2.4. Assume that Q is self-adjoint and that the transition spectral measure ν does not have atoms at ± 1 . Then, the spectral distribution function F defined by (1.7) is absolutely continuous with spectral density given by

$$f(t) = \frac{1}{2\pi} \int_{-1}^{1} \frac{1 - \lambda^2}{1 + \lambda^2 - 2\lambda \cos t} d\nu(\lambda), \tag{2.4}$$

and for $\alpha \in [1, 2)$, the following are equivalent:

(i)
$$var(S_n) = n^{\alpha}h(n)$$
 as $n \to \infty$,

(ii)
$$F(x) = \frac{1}{2}C(\alpha)x^{2-\alpha}h(1/x)$$
 as $x \to 0^+$.

Moreover, if $h(x) \to \infty$ as $x \to 0^+$, then (i), (ii) are equivalent to $\int_0^1 \frac{r(\mathrm{d}y)}{x^2 + y^2} = \frac{\pi}{2} C(\alpha) x^{1-\alpha} h(1/x) + O(1) \text{ as } x \to 0^+; \text{ where } r(0,y] = \nu(1-y,1).$

2.2 Relation between spectral measure and planar Brownian motion

Our next result makes essential use of the Poisson kernel which appears in (2.1) to provide a fascinating interpretation of the spectral distribution function F in terms of the harmonic measure of planar Brownian motion started at a random point in the open unit disk.

Theorem 2.5. Let ν be the transition spectral measure, and let $(B_t^z)_{t\geq 0}$ be standard planar Brownian motion in $\mathbb C$, started at the point $z\in D$. Also let Z be a random point in D distributed according to ν and let $\tau_D^Z:=\inf\{t\geq 0:B_t^Z\notin D\}$. Let $\Gamma_x:=\{z:z=\mathrm{e}^{\mathrm{i}ty},|y|< x\}$ and $\alpha\in(0,2)$. Then, the following statements are equivalent:

(i)
$$var(S_n) \sim n^{\alpha}h(n)$$
 as $n \to \infty$;

(ii)
$$P\{B_{ au_D^Z}^Z \in \Gamma_x\} \sim C(\alpha) x^{2-\alpha} h(1/x) / \nu(D)$$
 as $x \to 0$.

2.3 Linear growth of variance for partial sums for normal Markov Chains

By applying martingale techniques we establish necessary and sufficient conditions for the asymptotic linear variance behavior for general normal Markov chains.

Theorem 2.6. With the notation of Lemma 2.1

(a) The limit, $var(S_n)/n \to K < \infty$ exists if and only if

$$\sigma^{2} := \int_{D} \frac{1 - |z|^{2}}{|1 - z|^{2}} \nu(\mathrm{d}z) < \infty, \quad \text{and}$$
 (2.5)

$$I_n := \frac{1}{n} \int_D \frac{|1 - z^n|^2}{|1 - z|^2} \nu(\mathrm{d}z) \to L,$$
 (2.6)

where $K = \sigma^2 + L$.

- (b) Moreover, under (2.5) the following are equivalent:
 - (i) (2.6) holds with $L = \pi C$.
 - (ii) $\nu(U_x)/x \to C$ as $x \to 0^+$, where

$$U_x = \{z = (1 - r)e^{iu} \in D : 0 \le r \le |u| \le x\}.$$

(iii) $n\nu(D_n)\to C$ as $n\to\infty$, where

$$D_n = \{ z = re^{2i\pi\theta}; \ 1 - \frac{1}{n} \le r \le 1, \ -\frac{1}{n} \le \theta \le \frac{1}{n} \}.$$

It should be noted that there are many sufficient conditions for the convergence

$$var(S_n)/n \to \sigma^2 < \infty. \tag{2.7}$$

(1) It is immediate from the proof of Theorem 2.6, that (2.7) is equivalent to $\sigma^2 < \infty$ and

$$\frac{1}{n} \int_{D} \frac{|1 - z^{n}|^{2}}{|1 - z|^{2}} \nu(\mathrm{d}z) \to 0.$$
 (2.8)

(2) In Corollary 7.1 in [4] it was shown that if we assume

$$\int_D \frac{1}{|1-z|} \nu(\mathrm{d}z) < \infty,$$

then both (2.5) and (2.8) are satisfied and so convergence (2.7) holds, the result attributed to Gordin and Lifšic [15] (see Theorem 7.1 in [4]; see also [10]).

(3) From Representation Lemma 2.1 and Ibragimov's version(see [21]) of the Hardy-Littlewood theorem

$$\operatorname{var}(S_n)/n \to 2\pi f(0) = \sigma^2.$$

- (4) On the other hand, from the Representation Lemma 2.1 and Theorem A, convergence (2.7) is equivalent to the uniform integrability of $T_x(z)$ with respect to ν_0 as $x \to 0^+$.
- (5) Motivated by the complex Darboux-Wiener-Tauberian approach (e.g. as in [9]), by analyzing

$$V(\lambda) = \sum_{n=1}^{\infty} \operatorname{var}(S_n) \lambda^n = \frac{\lambda}{(1-\lambda)^2} \int_D \frac{1+\lambda z}{1-\lambda z} \nu(dz),$$

a sufficient condition for (2.7) is

$$\int_D \frac{1+\lambda z}{1-\lambda z} \nu(dz) \to \sigma^2 = \int_D \frac{1+z}{1-z} \nu(dz) \text{ as } \lambda \to 1 \text{ with } |\lambda| < 1.$$

Here, since $\nu(\mathrm{d}z)=\nu(\mathrm{d}\bar{z})$, the integral is understood in the Cauchy sense i.e.

$$\int_D f(z)\nu(dz) = \frac{1}{2} \int_D [f(z) + f(\bar{z})]\nu(dz).$$

- (6) From Theorem 2.6, it follows that (2.7) is equivalent to $\sigma^2 < \infty$ and $\nu(U_x)/x \to 0$ as $x \to 0^+$.
- (7) Finally, again from Theorem 2.6, it follows that (2.7) is equivalent to $\sigma^2 < \infty$ and $n\nu(D_n) \to 0$ as $n \to \infty$. This result is also a corollary of the following inequality motivated by Cuny and Lin [7]

$$\frac{n\nu(D_n)}{36} \le \frac{1}{n} \mathbb{E}(\mathbb{E}_1(S_n))^2 \le \frac{4}{n} \sum_{j=1}^{n-1} j\nu(D_j).$$

Remark 2.7. Notice that when the transition spectral measure ν is concentrated on Γ (dynamical system), then $\sigma^2 = 0$ and so σ^2 cannot be the limiting variance, in general.

By inspecting the proof of the Theorem 2.6, under (2.5) we are able to characterize regular variation of $\text{var}(S_n)$ when $\liminf_{n\to\infty} \text{var}(S_n)/n > 0$. More exactly, we have the following proposition.

Proposition 2.8. Assume $\liminf_{n\to\infty} \operatorname{var}(S_n)/n>0$ and that (2.5) holds. Then, for $\alpha\in[1,2)$, and a positive function h, slowly varying at infinity, $\operatorname{var}(S_n)\sim n^\alpha h(n)$ as $n\to\infty$ if and only if $\nu(U_x)\sim C(\alpha)x^{2-\alpha}h(1/x)$ as $x\to 0^+$, with $C(\alpha)$ as defined in Theorem A. In particular, $\operatorname{var}(S_n)/n$ is slowly varying at $n\to\infty$ if and only if $\nu(U_x)/x$ is slowly varying as $x\to 0^+$.

2.4 Relation between the variance of partial sums and transition spectral measure of reversible Markov chains

We continue the study of stationary reversible Markov chains and provide further necessary and sufficient conditions for its variance to be regularly varying, in terms of the operator spectral measure by a direct approach, without the link with the spectral distribution function.

Theorem 2.9. Assume Q is self-adjoint, $\alpha \geq 1$, $var(S_n)/n \rightarrow \infty$, and let $c_\alpha := \alpha(2 - \alpha)/2\Gamma(3 - \alpha)$. Then

$$V(x) = \int_{-1}^{1-x} \frac{1}{1-t} \nu(\mathrm{d}t) \sim c_{\alpha} x^{1-\alpha} h(\frac{1}{x}) \text{ as } x \to 0^+$$

if and only if

$$var(S_n) = n^{\alpha}h(n) \text{ as } n \to \infty.$$

Furthermore if $\alpha > 1$ then $var(S_n) = n^{\alpha}h(n)$ as $n \to \infty$ iff

$$\nu(1-x,1] \sim d_{\alpha}x^{2-\alpha}h(\frac{1}{x}) \text{ as } x \to 0^+,$$

where $d_{\alpha} := \alpha(\alpha - 1)/2\Gamma(3 - \alpha)$.

Remark 2.10. It should be obvious from the statement in the above theorem, that regular variation of the variance is equivalent to regular variation of the transition spectral measure only in the case $\alpha > 1$. As the following example demonstrates, in the case $\alpha = 1$, there are reversible Markov chains whose variance of partial sums varies regularly with exponent 1 even though $\nu(1-x,1]$ is not a regularly varying function.

Example 2.11. Take a probability measure v on [-1,1] defined for 0 < a < 1/2 by

$$dv = \frac{1}{c}(1 - |x|) \Big(1 + a \sin[\ln(1 - |x|)] + a \cos[\ln(1 - |x|)] \Big) dx.$$
 (2.9)

where c is the normalizing constant. Then, the unique invariant measure is

$$d\pi = \frac{dv}{\theta(1-|x|)} = \frac{1}{2} \Big(1 + a \sin[\ln(1-|x|)] + a \cos[\ln(1-|x|)] \Big) dx.$$

We first compute the following integral

$$\begin{split} \int_0^{1-x} \frac{\mathrm{d}\pi}{1-y} &= \frac{1}{2} \int_x^1 \frac{1}{t} \Big(1 + a \sin[\ln(t)] + a \cos[\ln(t)] \Big) \mathrm{d}t \\ &= -\frac{1}{2} \ln x + O(1) \text{ as } x \to 0, \end{split}$$

whence, by Theorem 2.9

$$\lim_{n \to \infty} \frac{\operatorname{var}(S_n)}{n \log n} = c.$$

However, the covariances are not regularly varying because the spectral measure is not. To see why it is enough to show that r(x) is not regularly varying at 0. Indeed, if we take $y_k = e^{-2\pi k} \to 0^+$, and $y_k = e^{\pi/2 - 2\pi k} \to 0^+$, then $r(y_k) = y_k$ and $r(y_k)/y_k \to 1$. However, for the choice $r(z_k) = z_k(1+2\alpha)$, we have $r(z_k)/z_k \to 1+2\alpha$, and hence the spectral measure is not regularly varying.

Remark 2.12. Often in the literature, conditions for the linear growth of the variance are given in terms of the covariances (see for example [18]). As it turns out, one can construct positive covariance sequences such that $\sum_{k=0}^n \operatorname{cov}(X_0, X_k) = h(n)$ is slowly varying, and hence the variance is regularly varying, but $a_n = \operatorname{cov}(X_0, X_n) > 0$, is not slowly varying. To construct such a chain, suppose that ε_n is an oscillating positive sequence such that $\varepsilon_n \to 0$ and $\sum_k a_k = \infty$ where $a_k := \varepsilon_k/k$. Then $g_n = \sum_{k=1}^n a_k$ is slowly varying since

$$g_{[bn]} = g_n + \sum_{i=n+1}^{[bn]} \frac{\varepsilon_i}{i} = g_n + O(\varepsilon_{n+1} \log b).$$

So, a priori we have many situations when $var(S_n) = nh(n)$ even though the covariances (and hence the operator spectral measure) are not regularly varying.

The above proof was direct in the sense, that it relied only on the use of classical Tauberian theory without linking the transition spectral measure with the spectral distribution function, and thus without invoking the results of [8].

2.5 Examples of limit theorems with non-linear normalizer

As an application we construct a class of stationary irreducible and aperiodic Markov Chains, based on the Metropolis-Hastings algorithm, with $var(S_n) \sim nh(n)$. Markov chains of this type are often studied in the literature from different points of view, in Doukhan et al. [11], Rio ([31] and [32]), Merlevède and Peligrad [28], Zhao et al. [39]) and Longla et al. [26].

Let $E = \{|x| \leq 1\}$ and define the transition kernel of a Markov chain by

$$Q(x, A) = |x|\delta_x(A) + (1 - |x|)v(A),$$

where δ_x denotes the Dirac measure and v is a symmetric probability measure on [-1,1] in the sense that for any $A \subset [0,1]$ we have v(A) = v(-A). We shall assume that

$$\theta = \int_{-1}^{1} \frac{1}{(1 - |x|)} \nu(\mathrm{d}x) < \infty. \tag{2.10}$$

We mention that Q is a stationary transition function with the invariant distribution

$$\mu(\mathrm{d}x) = \frac{1}{\theta(1-|x|)} v(\mathrm{d}x).$$

Then, the stationary Markov chain $(\xi_i)_i$ with values in E and transition probability Q(x,A) and with marginal distribution μ is reversible and positively recurrent. Moreover, for any odd function g we have $Q^k(g)(x) = |x|^k g(x)$ and therefore

$$Q^{k}(g)(\xi_{0}) = \mathbb{E}(g(\xi_{k})|\xi_{0}) = |\xi_{0}|^{k} g(\xi_{0}) \text{ a.s.}$$
(2.11)

For the odd function $g(x) = \operatorname{sgn}(x)$, define $X_i := \operatorname{sgn}(\xi_i)$. Then for any positive integer k

$$\langle g, Q^k(g) \rangle = \int_{-1}^1 |x|^k \mu(\mathrm{d}x) = 2 \int_0^1 x^k \mu(\mathrm{d}x),$$

and so on [0,1], 2μ coincides with the transition spectral measure ν , associated to Q and g. Furthermore, the operator Q is of positive type $\nu[-1,0)=0$. In other words

$$\nu = \left\{ \begin{array}{l} 2\mu \text{ on } [0,1] \\ 0 \text{ on } [-1,0) \end{array} \right. . \label{eq:definition}$$

Therefore, by Theorem 2.9 applied with $\alpha = 1$, $var(S_n) = nh(n)$ with $h(n) \to \infty$ slowly varying at $n \to \infty$ if and only if

$$V(x) = \int_0^{1-x} \frac{\nu(\mathrm{d}y)}{1-y} = \int_0^{1-x} \frac{2\mu(\mathrm{d}y)}{1-y} \sim \frac{1}{2}h(\frac{1}{x}),$$

is slowly varying as $x \to 0^+$.

Our next result presents a large class of transition spectral measures for the model above which leads to functional central limit theorem.

Theorem 2.13. Let V(x) be slowly varying as $x \to 0^+$. Then, the central limit theorem, the functional central limit theorem and the conditional central limit theorem hold for partial sums of X_i defined above.

Next, we give a particular example of a Metropolis-Hastings algorithm in which a non-degenerate central limit theorem holds under a certain normalization. However, when normalized by the standard deviation we have degenerate limiting distribution.

Example 2.14. For $x \in (0,1)$ consider the slowly varying function $V(x) = \exp(\sqrt{\ln(1/x)})$. By Theorem 2.9, as $n \to \infty$

$$var(S_n) = 2nV(1/n)(1 + o(1)).$$

On the other hand, let us choose b_n such that $nE[\tau_1^2I(\tau_1\leq b_n)]/b_n^2\sim 1$ as $n\to\infty$. By Lemma A.3 it follows that $2\theta nV(1/b_n)\sim b_n^2$, with θ as defined in Eq. (2.10). Note now that

$$2\theta nV(1/b_n) = 2\theta n \exp(\sqrt{\ln(b_n)}) = 2\theta n \exp(\sqrt{(1/2)\ln(4nV(1/b_n))})$$
$$= 2\theta n \exp(o(1) + \sqrt{(1/2)\ln n} + (1/2)^{3/2}),$$

which implies that $b_n^2 \sim 2n\theta \exp(\sqrt{(1/2) \ln n})$, giving the following CLT:

$$S_n/[2\theta n \exp(\sqrt{(1/2)\ln n})^{1/2} \Rightarrow N(0,1).$$

However

$$\frac{b_n^2}{nV(1/n)} \to 0 \quad \text{and therefore } \frac{S_n}{\sqrt{\mathrm{var}(S_n)}} \to^{\mathbb{P}} 0.$$

3 Continuous-time Markov processes

Suppose we have a stationary Markov process $\{\xi_t\}_{t\geq 0}$, with values in the general state space (S,\mathcal{A}) , and for $g\in L^2_0(\pi)$ let $T_tg(x):=\mathbb{E}[g(\xi_t)|\xi_0=x]$. Further $T_t=\mathrm{e}^{Lt}$, where L is the infinitesimal generator which we assume to be normal, which then implies that its spectrum is supported on $\{z\in\mathbb{C}:\Re(z)\leq 0\}$, such that

$$\operatorname{cov}(f(\xi_t), f(\xi_0)) = \int_{\Re(z) < 0} e^{zt} \nu(\mathrm{d}z).$$

Finally define $S_T(g) := \int_{s=0}^T g(\xi_s) ds$.

The following result is a continuous time analogue of Theorem 2.5, linking the spectral distribution function with the harmonic measure of planar Brownian motion.

Theorem 3.1. Let $\{\xi_t\}_t$ be a stationary Markov process with normal generator L, invariant measure π , and let $g \in L^2(\pi)$ and $\nu = \nu_f$ be the transition spectral measure associated with L and g. Write $(B_t^z)_{t \geq 0}$ for a standard planar Brownian motion in $\mathbb C$, started at the point $z \in \mathbb H^- := \{z \in \mathbb C : \Re(z) \leq 0\}$. Also let Z be a random point in $\mathbb H^-$ distributed according to ν and let $\tau_{\mathbb H^-}^Z := \inf\{t \geq 0 : B_t^Z \notin \mathbb H^-\}$. For $\alpha \in (0,2)$, the following statements are equivalent:

- (i) $var(S_T(g)) \sim T^{\alpha}h(T)$ as $T \to \infty$;
- $\text{(ii)} \ \ \mathrm{P}\big\{B^Z(\tau^Z_{\mathbb{H}^-}) \in (-\mathrm{i} x, \mathrm{i} x)\big\} \sim C(\alpha) x^{2-\alpha} h(1/x)/\nu(\mathbb{H}^-) \ \text{as} \ x \to 0^+.$

The following theorem gives a necessary and sufficient condition in terms of the transition spectral measure ν . Define for x > 0,

$$U_x := \{a + ib : 0 \le -a \le |b| \le x\}.$$

and let

$$\varsigma^2 := -2 \int_{\mathbb{H}^-} \Re(1/z) \nu(\mathrm{d}z).$$

Theorem 3.2. With the notation of Theorem 3.1 the following are equivalent:

- (i) $var(S_T(g))/T \rightarrow L = \varsigma^2 + K$, where K > 0;
- (ii) $\varsigma^2 < \infty$ and $\nu(U_x)/x \to K/\pi$ as $x \to 0^+$.

In addition, if $\varsigma^2 < \infty$ and $\liminf_{T \to \infty} \text{var}(S_T)/T = \infty$, then $\text{var}(S_T) \sim T^{\alpha}h(T)$, for $\alpha \geq 1$ and h slowly varying at infinity, if and only if $\nu(U_x) \sim C(\alpha)x^{2-\alpha}h(1/x)$.

4 Proofs

Proof of Representation Lemma 2.1. For $t \in [-\pi, \pi]$ and $z \in D_0$ define the function

$$f(t,z) := \frac{1}{2\pi} \left[1 + \sum_{k=1}^{\infty} \left(z^k e^{itk} + \bar{z}^k e^{-itk} \right) \right] = \frac{1}{2\pi} \frac{1 - |z|^2}{|1 - ze^{it}|^2},$$

the Poisson kernel for the unit disk. Our approach is to integrate on D_0 with respect to $\nu_0(\mathrm{d}z)$, obtaining in this way a function defined on $[0,2\pi]$ as follows

$$f(t) = \frac{1}{2\pi} \int_{D_0} \left(1 + \sum_{k=1}^{\infty} \left(z^k e^{itk} + \bar{z}^k e^{-itk} \right) \right) \nu_0(dz) = \frac{1}{2\pi} \int_{D_0} \frac{1 - |z|^2}{|1 - ze^{it}|^2} \nu_0(dz).$$

The function is well defined since we are integrating the positive Poisson kernel over the open disk, and in fact, by using polar coordinates, we also have

$$\int_{0}^{2\pi} f(t) dt = \frac{1}{2\pi} \int_{s=0}^{2\pi} \int_{r=0}^{1^{-}} \int_{t=0}^{2\pi} \frac{1 - r^{2}}{1 - 2r \cos(s + t) + r^{2}} dt \, \nu_{0}(dr, ds)$$

$$\leq \frac{1}{2\pi} \int_{s=0}^{2\pi} \int_{r=0}^{1^{-}} \frac{2\pi (1 - r^{2})}{1 - r^{2}} \, \nu_{0}(dr, ds) = 2\pi < \infty.$$

Therefore, it is obvious that $f \in L^1(0,2\pi)$, and it makes sense to calculate

$$\int_{0}^{2\pi} e^{itn} f(t) dt = \frac{1}{2\pi} \int_{0}^{2\pi} e^{itn} \int_{D_{0}} \left[1 + \sum_{k=1}^{\infty} \left(z^{k} e^{itk} + \bar{z}^{k} e^{-itk} \right) \right] \nu_{0}(dz) dt$$

$$= \int_{D_{0}} z^{n} \nu_{0}(dz).$$
(4.1)

Because of the decomposition

$$\operatorname{cov}(X_0, X_n) = \int_D z^n d\nu(z) = \int_{D_0} z^n d\nu_0(z) + \int_0^{2\pi} e^{itn} \nu_{\Gamma}(dt),$$

by (4.1) we obtain that

$$\operatorname{cov}(X_0, X_n) = \int_0^{2\pi} e^{\mathrm{i}tn} f(t) dt + \int_0^{2\pi} e^{\mathrm{i}tn} \nu_{\Gamma}(dt).$$

Now, by (1.7) the spectral distribution function F associated with the stationary sequence $(X_i)_{i}$, is then given by (2.2).

Proof of Corollary 2.4. The result follows from Representation Lemma 2.1 and Theorem A. To obtain the last point in the theorem, by standard analysis, the spectral measure has the following useful asymptotic representation

$$F(x) = \int_0^x f(t)dt = \frac{1}{2\pi} \int_0^x \int_{-1}^1 \frac{1 - \lambda^2}{1 + \lambda^2 - 2\lambda \cos t} \nu(d\lambda) dt$$
$$= O(x) + \frac{x}{\pi} \int_0^1 r(y) \frac{dy}{x^2 + y^2} \quad \text{as } x \to 0^+.$$

So, we derive $F(x) \sim \frac{1}{2}C(\alpha)x^{2-\alpha}h(1/x)$ if and only if

$$\int_0^1 r(y) \frac{\mathrm{d}y}{x^2 + y^2} \sim \frac{\pi}{2} C(\alpha) x^{1-\alpha} h(1/x) + O(1) \quad \text{as } x \to 0^+.$$

Proof of Theorem 2.5. As usual, let D be the closed unit disk, and D_0 its interior. From Representation Lemma 2.1 the spectral density $f \in L_1([-\pi, \pi])$ is given by the formula

$$f(t) = \frac{1}{2\pi} \int_{D_0} \frac{1 - |z|^2}{|1 - ze^{it}|^2} \nu(dz).$$

Notice that D is regular for Brownian motion, in the sense that all points in $\Gamma = \partial D$ are regular, i.e. for all $z \in \partial D$ and for $\tilde{\tau}_D^z := \inf\{t > 0 : B_t^z \notin D\}$ we have $\mathrm{P}^z\{\tilde{\tau}_D^z = 0\} = 1$. The harmonic measure in D from z is the probability measure on ∂D , $\mathrm{hm}(z,D;\cdot)$ given by

$$hm(z, D; V) = P^{z} \{ B(\tau_D^z) \in V \},$$

where P^z denotes the probability measure of Brownian motion started at the point z, and V is any Borel subset of ∂D .

Since ∂D is piecewise analytic, $\operatorname{hm}(z,D;\cdot)$ is absolutely continuous with respect to Lebesgue measure (length) on ∂D and the density is the Poisson kernel (see for example [25]). In the case of the unit disk D the density for $\operatorname{hm}(z,D;\cdot)$ for $z\in D$, $w\in \partial D$ or $t\in [0,2\pi]$, is given by

$$H_D(z, w) = \frac{1}{2\pi} \frac{1 - |z|^2}{|w - z|^2} = \frac{1}{2\pi} \frac{1 - |z|^2}{|1 - e^{it}z|^2}.$$

Let Z be a D-valued random variable with probability measure ν properly normalized, independent of the Brownian motion. Then

$$\int_{-x}^{x} f(t)dt = \frac{1}{2\pi} \int_{D_0} \int_{-x}^{x} \frac{1 - |z|^2}{|1 - ze^{it}|^2} dt \, \nu_0(dz)$$
$$= \int_{D_0} P^z \{B(\tau_D^z) \in (-x, x)\} \nu_0(dz).$$

On the other hand, since $\Gamma=\partial D$ is regular for Brownian motion, we have for all $z\in\Gamma$, $P^z(B(\tau_D^z)\in\Gamma_x)=1$ if $z\in\Gamma_x$ and 0 otherwise. Thus

$$\nu(\Gamma_x) = \int_{\Gamma_x} \nu_{\Gamma}(\mathrm{d}z) = \int_{\Gamma} \mathrm{P}^z(B(\tau_D^z) \in \Gamma_x) \nu_{\Gamma}(\mathrm{d}z).$$

Therefore from Representation Lemma 2.1 we have

$$G(x) := \int_{-x}^{x} F(\mathrm{d}x) = \int_{D_0} \mathrm{P}^z \{ B(\tau_D^z) \in \Gamma_x \} \nu_0(\mathrm{d}z) + \int_{\Gamma} \mathrm{P}^z (B(\tau_D^z) \in \Gamma_x) \nu_{\Gamma}(\mathrm{d}z)$$
$$= \nu(D) \mathrm{P} \Big(B^Z(\tau_D^Z) \in \Gamma_x \Big) =: \nu(D) H_{\nu,D}(\Gamma_x).$$

The measure $H_{\nu,D}(\cdot)$, is essentially the harmonic measure when Brownian motion starts at a random point and stops when it hits ∂D . Finally, from Theorem A we conclude that $var(S_n)$ is regularly varying if and only if the measure $H_{\nu,D}$ is regularly varying at the origin.

Proof of Theorem 2.6. The first part is motivated by Gordin and Lifšic [15] (see Theorem 7.1 in [4]; see also [10]). We write the martingale type orthogonal decomposition:

$$S_n = \mathbb{E}_0(S_n) + \sum_{i=1}^n \mathbb{E}_i(S_n - S_{i-1}) - \mathbb{E}_{i-1}(S_n - S_{i-1}).$$

So

$$var(S_n) = \mathbb{E}(\mathbb{E}_0(S_n))^2 + \sum_{i=1}^n \mathbb{E}(\mathbb{E}_i(S_n - S_{i-1}) - \mathbb{E}_{i-1}(S_n - S_{i-1}))^2$$
$$= \mathbb{E}(\mathbb{E}_0(S_n))^2 + \sum_{i=1}^n \mathbb{E}(\mathbb{E}_1(S_i) - \mathbb{E}_0(S_i))^2.$$

By applying spectral calculus,

$$= \mathbb{E}(\mathbb{E}_{0}(S_{n}))^{2} + \sum_{j=1}^{n} \int_{D} |1 + z + \dots + z^{j-1}|^{2} (1 - |z|^{2}) \nu(\mathrm{d}z)$$

$$= \mathbb{E}(\mathbb{E}_{0}(S_{n}))^{2} + \sum_{j=1}^{n} \int_{D_{0}} \frac{|1 - z^{j}|^{2} (1 - |z|^{2})}{|1 - z|^{2}} \nu(\mathrm{d}z)$$

$$= \mathbb{E}(\mathbb{E}_{0}(S_{n}))^{2} + n \int_{D_{0}} \delta_{n}(z) \frac{(1 - |z|^{2})}{|1 - z|^{2}} \nu(\mathrm{d}z).$$

Note that $\delta_n(z) \leq 4$ and for all $z \in D_0$

$$\delta_n(z) = \frac{1}{n} \sum_{j=1}^n |1 - z^j|^2 \to 1 \text{ as } n \to \infty.$$

Thus, by the Lebesgue dominated theorem and our conditions

$$\lim_{n \to \infty} \int_{D_0} \delta_n(z) \frac{(1 - |z|^2)}{|1 - z|^2} \nu(\mathrm{d}z) = \int_{D_0} \frac{1 - |z|^2}{|1 - z|^2} \nu(\mathrm{d}z) = \sigma^2.$$

This, along with Fatou's lemma, proves that $\text{var}(S_n)/n \to K$ exists if and only if $\sigma^2 < \infty$ and

$$\frac{1}{n}\mathbb{E}(\mathbb{E}_0(S_n))^2 \to L = K - \sigma^2.$$

Now, let us introduce a new measure on D_0

$$\mu(z) = \frac{1 - |z|}{|1 - z|^2} \nu(\mathrm{d}z),$$

which is finite when $\sigma^2 < \infty$.

To complete the proof of the first part of the theorem, we notice that by the spectral calculus

$$\mathbb{E}(\mathbb{E}_{0}(S_{n}))^{2} = \int_{D} |z + \dots + z^{n}|^{2} \nu(\mathrm{d}z) = \int_{D} |z|^{2} \frac{|1 - z^{n}|^{2}}{|1 - z|^{2}} \nu(\mathrm{d}z)$$

$$= \int_{D} \frac{|1 - z^{n}|^{2}}{|1 - z|^{2}} \nu(\mathrm{d}z) - \int_{D} |1 - z^{n}|^{2} (1 + |z|) \mu(\mathrm{d}z)$$

$$= \int_{D} \frac{|1 - z^{n}|^{2}}{|1 - z|^{2}} \nu(\mathrm{d}z) + O(1),$$

since μ is a finite measure.

To prove the second part of this theorem, we show equivalence of (i) and (ii) and then of (ii) and (iii). In addition, throughout we use notation:

$$z = |z|e^{iArg(z)}, |z| = 1 - y, \theta = Arg(z).$$

We note first that

$$|1 - z^n|^2 = (1 - |z|^n)^2 + |z|^n \sin^2(n\theta/2).$$

The proof strategy consists in showing, several successive approximation steps, that

$$\frac{1}{n}\mathbb{E}(\mathbb{E}_0(S_n))^2 = \int_0^{\pi} \frac{\sin^2(n\theta/2)}{\sin^2(\theta/2)} G(\mathrm{d}\theta) + o(1),$$

for some appropriate measure G, and then to apply Theorem A. With this in mind we write

$$I_n = \frac{1}{n} \int_D \frac{|1 - z^n|^2}{|1 - z|^2} \nu(\mathrm{d}z)$$

$$= \frac{1}{n} \int_D \frac{|z|^n \sin^2(n \operatorname{Arg}(z)/2)}{|1 - z|^2} \nu(\mathrm{d}z) + \frac{1}{n} \int_D \frac{(1 - |z|^n)^2}{|1 - z|^2} \nu(\mathrm{d}z) =: I'_n + \Delta'_n.$$

Note that

$$\Delta'_n = \int_{D_0} \left(\frac{(1-|z|^n)}{n(1-|z|)} \right) (1-|z|^n) \frac{(1-|z|)}{|1-z|^2} \nu(\mathrm{d}z)$$
$$= \int_D \frac{1}{n} \left(\sum_{j=0}^{n-1} |z|^j \right) (1-|z|^n) \mu(\mathrm{d}z).$$

By Lebesgue dominated convergence theorem, since the bounded integral argument goes to 0 for each $|z| \le 1$, we have $\Delta'_n \to 0$ as $n \to \infty$.

Then, write

$$I'_n = \frac{1}{n} \int_D \frac{\sin^2(n\operatorname{Arg}(z)/2)}{|1 - z|^2} \nu(\mathrm{d}z) - \frac{1}{n} \int_D (1 - |z|^n) \frac{\sin^2(n\operatorname{Arg}(z)/2)}{|1 - z|^2} \nu(\mathrm{d}z)$$

=: $I''_n + \Delta''_n$,

and again

$$\Delta_n'' = \int_{D_0} \left(\frac{(1 - |z|^n)}{n(1 - |z|)} \right) \sin^2(n \operatorname{Arg}(z)/2) \frac{(1 - |z|)}{|1 - z|^2} \nu(\mathrm{d}z)$$
$$= \int_{D_0} \left(\frac{(1 - |z|^n)}{n(1 - |z|)} \right) \sin^2(n \operatorname{Arg}(z)/2) \mu(\mathrm{d}z).$$

Note that by Lebesgue dominated theorem $\Delta_n'' \to 0$ as $n \to \infty$, since the bounded integral argument goes to 0 for each |z| < 1.

Fix now a small positive a>0, recall that $z=(1-y)e^{i\theta}$ and define an auxiliary subset of D

$$D_a = \{z = (1 - y)e^{i\theta} : 0 < |\theta| < a, 0 < y < a\}.$$

Further, notice that by the dominated convergence theorem

$$\varepsilon_n = \int_{D_a} \left| \frac{\sin(n \operatorname{Arg}(z)/2)}{n \operatorname{Arg}(z)} \right| \mu(\mathrm{d}z) \to 0 \text{ as } n \to \infty,$$

since the bounded integral argument goes to 0 for each |z| < 1.

Let N be large enough, so that $|\epsilon_n| < 1$ for all $n \ge N$, and take $\delta_n = \max(\mathrm{e}^{-n}, \sqrt{\varepsilon_n})$ so that $\delta_n > 0$ for all n. In this way $\epsilon_n/\delta_n = 0$ is well-defined if $\epsilon_n = 0$, and $\epsilon_n/\delta_n \le \sqrt{\epsilon_n} \to 0$ as $n \to \infty$. Further define two auxiliary sequences of subsets of D_a

$$D_{a,n} = \{ z = (1 - y)e^{i\theta} : 0 < \delta_n | \theta | \le y \le a, |\theta| \le a \},$$

$$U_{a,n} = U_a \setminus D_{a,n} = \{ z = (1 - y)e^{i\theta} : 0 \le y < \delta_n |\theta|, |\theta| \le a \}.$$

Next, let

$$g_n(z) = \frac{\sin^2(n\text{Arg}(z)/2)}{n|1-z|^2}$$

and write

$$I_n'' = \frac{1}{n} \int_D g_n(z) \nu(dz) = \frac{O(1)}{n} + \int_{D_n} g_n(z) \nu(dz)$$

and

$$\int_{D_n} g_n(z)\nu(\mathrm{d}z) = \int_{D_n} g_n(z)\nu(\mathrm{d}z) + \int_{U_{n,n}} g_n(z)\nu(\mathrm{d}z) =: \Delta_n^{\prime\prime\prime} + I_n^{\prime\prime\prime}.$$

Notice that by construction ($\theta = Arg(z)$)

$$\begin{split} \Delta_n''' &\leq \int_{D_{a,n}} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| |\sin(n\theta/2)| \frac{|\theta|}{|1-z|^2} \nu(\mathrm{d}z) \\ &\leq \frac{1}{\delta_n} \int_{D_{a,n}} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| \frac{1-|z|}{|1-z|^2} \nu(\mathrm{d}z) \\ &= \frac{1}{\delta_n} \int_{D_{a,n}} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| \mu(\mathrm{d}z) \leq \frac{1}{\delta_n} \int_{D_a} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| \mu(\mathrm{d}z) \ = \frac{\varepsilon_n}{\delta_n} \to 0. \end{split}$$

In addition, on $U_{a,n}$

$$|1-z|^2 = y^2 + (1-y)\sin^2(\theta/2) = \sin^2(\theta/2)(1+\delta_{z,n}),$$

where $|\delta_{z,n}| \leq c\delta_n$ for some positive c and hence

$$I_n''' = \frac{1}{n} \int_{U_{a,n}} \frac{\sin^2(n\theta/2)}{|1-z|^2} \nu(\mathrm{d}z)$$

$$= \frac{1}{n} \int_{U_{a,n}} \frac{\sin^2(n\theta/2)}{\sin^2(\theta/2)} \nu(\mathrm{d}z) + \frac{1}{n} \int_{U_{a,n}} \frac{\sin^2(n\theta/2)}{\sin^2(\theta/2)} \left(1 - \frac{1}{1 + \delta_{z,n}}\right) \nu(\mathrm{d}z)$$

$$=: I_n^{(4)} + \Delta_n^{(4)}.$$

By construction $\Delta_n^{(4)} = o(1)I_n^{iv}$ and then $I_n''' = I_n^{(4)}(1+o(1))$. Write

$$I_n^{(4)} = \frac{1}{n} \int_{U_a} \frac{\sin^2(n\theta/2)}{\sin^2(\theta/2)} \nu(\mathrm{d}z) - \frac{1}{n} \int_{U_a \cap D_{a,n}} \frac{\sin^2(n\theta/2)}{\sin^2(\theta/2)} \nu(\mathrm{d}z).$$

= $I_n^{(5)} - \Delta_n^{(5)}$.

In a similar way as Δ_n''' has been estimated in Step 3, by construction ($\theta = \operatorname{Arg}(z)$) we obtain

$$\begin{split} &\Delta_n^{(5)} \leq \int_{U_a \cap D_{a,n}} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| |\sin(n\theta/2)| |(1+\delta_{z,n})| \frac{|\theta|}{|1-z|^2} \nu(\mathrm{d}z) \\ &\leq \frac{1+c\delta_n}{\delta_n} \int_{D_{a,n}} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| \frac{1-|z|}{|1-z|^2} \nu(\mathrm{d}z) \\ &\leq \frac{1+c\delta_n}{\delta_n} \int_{D_a} \Big| \frac{\sin(n\theta/2)}{n\theta} \Big| \mu(\mathrm{d}z) = \frac{(1+c\delta_n)\varepsilon_n}{\delta_n} \to 0 \text{ as } n \to \infty. \end{split}$$

Finally, let us define the set

$$D_{+} = \{z = (1 - y)e^{i\theta} : 0 < y < |\theta| < \pi\}.$$

Then, it follows that $I_n/n \to L$ if and only if

$$J_n = \int_{D_+} I_n(\operatorname{Arg}(z))\nu(\mathrm{d}z) \to 0,$$

where I_n is the Fejer kernel,

$$I_n(x) := \frac{\sin^2(nx/2)}{n\sin^2(x/2)}.$$

Define $G(x)=\nu(U_x)$ and notice that it is a non-negative, non-decreasing bounded function. In addition, for any step function $g(\theta)=I(u<|\theta|\leq v)=I(|\theta|\leq v)-I(|\theta|\leq w)$ with u,w being continuity points of G(x) we have

$$\int_{D_{+}} g(\operatorname{Arg}(z))\nu(dz) = \int_{D_{+}} I(g(\theta))\nu(dz) = \int_{D_{+}} I(|\theta| \le v)\nu(dz) - \int_{D_{+}} I(|\theta| \le w)\nu(dz)$$
$$= \nu(U_{v}) - \nu(U_{w}) = \int g(x)dG(x)$$

and then, by Caratheodory and Lebesgue theorem we have

$$J_n = \int_0^{\pi} I_n(x) dG(x).$$

By Theorem A, $J_n/n \to L$ if and only if $L = \pi C$ where $G(x)/x = \nu(U_x)/x \to C$ as $x \to 0^+$, completing the proof of equivalence of (i) and (ii) in Part (b).

To prove the equivalence of (ii) and (iii) in Part (b) under the finiteness of integral (2.5), let us define

$$W_x = \{ z = (1 - r)e^{iu} \in D : 0 \le |u| \le r \le x \},$$

$$D_{1/x} = \{ z = (1 - r)e^{iu} \in D : 0 \le |u|, r \le x \}.$$

Since $\nu(1)=0$, it follows that $\nu(W_x)\to 0$ as $x\to 0^+$.

On the other hand, on ${\cal W}_x$

$$\frac{1 - |z|}{|1 - z|^2} \ge \frac{1}{2x}$$

and hence by (2.5),

$$0 \leftarrow \int_{W_x} \frac{1-|z|}{|1-z|^2} \nu(\mathrm{d}z) \geq \frac{1}{2x} \nu(W_x),$$

which implies that $\nu(D_{1/x})/x \to C$ if and only if (ii) holds as $x \to 0^+$ and this completes the proof.

Remark 4.1. The sufficient part in Theorem 2.6 can be derived directly by performing coordinate transformation mapping the open unit disk to the upper half-plane and further careful analysis. Briefly, we change coordinates to the upper half-plane $\mathbb{H}:=\{(a,b):a\in\mathbb{R},b>0\}$, via the inverse Cayley transform $\phi(w):\mathbb{H}\to D_0$, where $\phi(w):=(1+\mathrm{i}w)/(1-\mathrm{i}w)$. The finite measure ν is transformed to a finite measure ρ on \mathbb{H} , which for simplicity we can assume to be a probability measure. Then, for $w:=a+\mathrm{i}b$ we have

$$\sigma^2 = \int_{\mathbb{H}} \frac{1 - |\phi(a + ib)|^2}{|1 - \phi(a + ib)|^2} dv(\phi^{-1}(z)) = \int_{a = -\infty}^{\infty} \int_{b=0}^{\infty} \frac{b}{a^2 + b^2} d\rho(a, b).$$

and a further change of variables $z = \tan(t/2)$ gives

$$\int_{t=0}^x f(t)\mathrm{d}t = \frac{1}{2\pi} \int_{t=0}^x \iint_{\mathbb{H}} \frac{b}{(a+t/2)^2 + b^2} \rho(\mathrm{d}a, \mathrm{d}b)\mathrm{d}t + o(x).$$

Proof of Theorem 2.9. Let $1 \le \alpha \le 2$ and denote $C(n) := \sum_{i=0}^{n-1} \operatorname{cov}(X_0, X_i)$. We start by the well known representation

$$var(S_n) = n\left[\frac{2}{n}\sum_{k=1}^n C(k) - E(X_0^2)\right].$$

It is clear then, since $var(S_n)/n \to \infty$, that $var(S_n)$ has the same asymptotic behavior as $2\sum_{k=1}^n C(k)$. Implementing the notations

$$a_j = \int_0^1 x^j \nu_1(\mathrm{d}x)$$
 and $b_j = \int_0^0 x^j \nu(\mathrm{d}x),$

where ν_1 coincides with ν on (0,1], and $\nu_1(\{0\})=0$, we have the representation

$$C(k) = \sum_{j=0}^{k-1} a_j + \sum_{j=0}^{k-1} b_j = C_1(k) + C_2(k).$$

We shall show that the terms $C_1(k)$ have a dominant contribution to the variance of partial sum. To analyze $C_2(k)$ it is convenient to make blocks of size 2, a trick that has also appeared in [18] where it is attributed to [16]. We notice that

$$c_l = b_{2l} + b_{2l+1} = \int_{-1}^{0} (x^{2l} + x^{2l+1}) d\nu > 0.$$

Furthermore, for all m

$$\sum_{l=0}^{m-1} c_l = \sum_{l=0}^{m-1} \int_{-1}^0 x^{2l} (1+x) \nu(\mathrm{d}x) = \int_{-1}^0 \frac{1-x^{2m}}{1-x} \nu(\mathrm{d}x) \le E(X_0^2).$$

Therefore $|C_2(k)| \leq 2E(X_0^2)$ and so, $\sum_{k=1}^n C_2(k) \leq 2nE(X_0^2)$. Because $\operatorname{var}(S_n)/n \to \infty$ we note that $\operatorname{var}(S_n)$ has the same asymptotic behavior as $2\sum_{k=1}^n C_1(k)$. Now, each $C_1(k) = \sum_{j=0}^{k-1} a_k$ with $a_k > 0$. So, because the sequence $(C_1(k))_k$ is increasing, by the monotone Tauberian theorem (see Theorem A.2 in [3, Cor 1.7.3]) for all $\alpha \geq 1$ we have

$$\operatorname{var}(S_n) \sim 2\sum_{k=1}^n C_1(k) \sim n^{\alpha} h(n) \quad \text{if and only if} \quad C_1(n) \sim \alpha n^{\alpha-1} h(n)/2. \tag{4.2}$$

Note now that

$$C_1(n) = \int_0^1 \frac{1 - x^n}{1 - x} \nu_1(\mathrm{d}x).$$

It is convenient to consider the transformation $T:[0,1]\to [0,1]$ defined by T(x)=1-x. For a Borelian A of [0,1] define the measure

$$r(A) = \nu_1(T(A)).$$
 (4.3)

Then

$$C_1(n) = \int_0^1 \frac{1 - (1 - y)^n}{y} r(\mathrm{d}y).$$

We shall integrate by parts. Denote

$$R(u) = \int_{u}^{1} \frac{1}{y} r(\mathrm{d}y)$$
 and $U_{n}(u) = [1 - (1 - u)^{n}].$

Let 0 < b < 1. By the definition of ν_1 we have $r(\{1\}) = 0$. Since U_n is continuous, by [1, Theorem 18.4],

$$\int_{1}^{1} \frac{1 - (1 - y)^{n}}{y} r(dy) = \left[1 - (1 - b)^{n}\right] R(b) + n \int_{1}^{1} (1 - u)^{n-1} R(u) du.$$

Note that

$$\lim \sup_{b \to 0^+} [1 - (1 - b)^n] R(b) \le nbR(b) \le n \mathbb{E}(X_0^2) = o(\sigma_n^2).$$

Therefore

$$C_1(n) = o(\sigma_n^2) + n \int_0^1 (1-u)^{n-1} R(u) du.$$

By the change of variables $1 - u = e^{-y}$ we have

$$\int_0^1 (1-u)^{n-1} R(u) du = \int_0^\infty R(1-e^{-y}) e^{-yn} dy.$$

It follows that

$$C_1(n) \sim \alpha n^{\alpha-1}h(n)/2 \quad \text{if and only if} \quad n \int_0^\infty R(1-\mathrm{e}^{-y})\mathrm{e}^{-yn}dx \sim \alpha n^{\alpha-1}h(n)/2 \quad \text{(4.4)}$$
 if and only if
$$\int_0^\infty R(1-\mathrm{e}^{-y})\mathrm{e}^{-yn}\mathrm{d}x \sim \alpha n^{\alpha-2}h(n)/2.$$

From here we shall apply Karamata's Tauberian Theorem A.1. Since $\alpha \leq 2$,

$$\int_0^\infty R(1-\mathrm{e}^{-y})\mathrm{e}^{-yn}\mathrm{d}x \sim \alpha n^{\alpha-2}h(n)/2 \text{ as } n\to\infty \tag{4.5}$$
 if and only if
$$\int_0^x R(1-\mathrm{e}^{-y})\mathrm{d}y \sim c_\alpha' x^{2-\alpha}h(1/x) \text{ as } x\to 0^+,$$

where $c_{\alpha}' = \alpha/[2\Gamma(3-\alpha)]$. Again by the monotone Karamata Theorem A.2 this happens if and only if

$$R(1 - e^{-y}) \sim (2 - \alpha)c'_{\alpha}x^{1-\alpha}h(1/x) \text{ as } x \to 0^+.$$
 (4.6)

Changing variables $x = 1 - e^{-y}$ and taking into account Karamata's representation for slowly varying functions we get

$$R(x) \sim c_{\alpha} \left[\ln(1/(1-x)) \right]^{1-\alpha} h(-1/\ln(1-x))$$

$$\sim c_{\alpha} x^{1-\alpha} h(1/x) \text{ as } x \to 0^{+},$$
(4.7)

where $c_{\alpha}=(2-\alpha)c_{\alpha}'$. By combining the results in relations (4.2)-(4.7) we have that

$$\operatorname{var}(S_n) \sim n^{\alpha} h(n)$$
 if and only if $R(x) \sim c_{\alpha} x^{1-\alpha} h(1/x)$ as $x \to 0^+$.

It remains to note that

$$R(x) = \int_{x}^{1} \frac{1}{y} r(\mathrm{d}y) = \int_{0}^{1-x} \frac{1}{1-y} v_1(\mathrm{d}y) \sim V(x) \text{ as } x \to 0^+.$$

When $1 < \alpha < 2$, one can say more: the distribution function induced by the spectral measure is regularly varying.

Note that again by Theorem A.2, since the sequence a_k is a monotone sequence of positive numbers and $\alpha - 1 > 0$ we have

$$C_1(n) = \sum_{k=0}^{n-1} a_k = \alpha n^{\alpha-1} h(n)/2 \quad \text{if and only if} \quad a_n \sim \alpha(\alpha-1) n^{\alpha-2} h(n)/2.$$

Now, by considering the mapping $T':[0,1]\to[0,\infty)$, given by T'^{-x} , we obtain

$$a_n = \int_0^1 t^n \ \nu_1(\mathrm{d}t) = \int_0^\infty e^{-nx} \psi(\mathrm{d}x),$$

where $\psi(A) = r(T'(A))$, for A Borelian in $[0, \infty)$. Letting $d_{\alpha} := \alpha(\alpha - 1)/2\Gamma(3 - \alpha)$, it follows by Theorem A.1 ([3, Thm 1.7.1']), that

$$a_n \sim \alpha(\alpha - 1)n^{\alpha - 2}h(n)/2$$
 as $n \to \infty$ iff $\psi[0, x] \sim d_{\alpha}x^{2-\alpha}h(1/x)$ as $x \to 0^+$.

Then, we obtain as before, by the properties of slowly varying functions

$$\psi[0,x] \sim d_{\alpha}x^{2-\alpha}h(1/x)$$
 as $x \to 0^+$ iff $\nu_1(1-x,1] \sim d_{\alpha}x^{2-\alpha}h(1/x)$ as $x \to 0^+$.

This last relation combined with (4.4) gives the last part of the theorem.

Proof of Theorem 2.13. We only prove the central limit theorem, with the other results following in a similar manner. Our approach is based on the regeneration process. Define

$$T_0 = \inf\{i > 0 : \xi_i \neq \xi_0\}$$

$$T_{k+1} = \inf\{i > T_k : \xi_i \neq \xi_{i-1}\},$$

 \Box

and let $\tau_k := T_{k+1} - T_k$. It is well known that $(\xi_{\tau_k}, \tau_k)_{k \geq 1}$ are i.i.d. random variables with ξ_{τ_k} having the distribution v. Furthermore,

$$\mathbb{P}(\tau_1 > n | \xi_{\tau_1} = x) = |x|^n.$$

Then, it follows that

$$\mathbb{E}(\tau_1|\xi_{\tau_1}=x) = \frac{1}{1-|x|} \quad \text{and} \quad \mathbb{E}(\tau_1) = \int_{-1}^1 \frac{1}{1-|x|} \nu(\mathrm{d}x) = \theta.$$

So, by the law of large numbers $T_n/n \to \theta$ a.s. Let us study the tail distribution of τ_1 . Since

$$\mathbb{P}(\tau_1|X_{\tau_1}| > u|\xi_{\tau_1} = x) = \mathbb{P}(\tau_1 > u|\xi_{\tau_1} = x) = |x|^u,$$

by integration we obtain

$$\mathbb{P}(\tau_1 > u) = \int_{-1}^1 |x|^u v(\mathrm{d}x) = 2 \int_0^1 x^u v(\mathrm{d}x). \tag{4.8}$$

Using now the relation between v(dx) and $\mu(dx)$ and symmetry we get

$$\mathbb{P}(\tau_1 > u) = 2\theta \int_0^1 x^u (1 - x) \mu(\mathrm{d}x) = \theta \int_0^1 x^u (1 - x) \nu(\mathrm{d}x),$$

where ν is spectral measure. By the fact that V(x) is slowly varying and Lemma A.3 in Section A it follows that

$$H(u) := E[\tau_1^2 I(\tau_1 \le u)]$$
 is slowly varying as $u \to \infty$. (4.9)

For each n, let m_n be such that $T_{m_n} \leq n < T_{m_n+1}$. Note that we have the representation

$$\sum_{k=1}^{n} X_k - \sum_{k=1}^{[n/\theta]} Y_k = (T_0 - 1)X_0 + (\sum_{k=1}^{m_n} \tau_k X_{\tau_k} - \sum_{k=1}^{[n/\theta]} \tau_k X_{\tau_k}) + \sum_{k=T_{m_n+1}}^{n} X_k, \tag{4.10}$$

where $Y_k = \tau_k X_{\tau_k}$ is a centered i.i.d. sequence, and by (4.9) is in the domain of attraction of a normal law (see Feller [13]). Therefore,

$$\frac{\sum_{k=1}^{[n/\theta]} Y_k}{b_{[n/\theta]}} \Rightarrow N(0,1). \tag{4.11}$$

where $b_n^2 \sim nH(b_n)$. The rest of the proof is completed on the same lines as in the proof of Example 12 in [26], the final result being that

$$\frac{\sum_{k=1}^{n} X_k}{b_{[n/\theta]}} \Rightarrow N(0,1).$$

Proof of Theorem 3.1. The proof is similar to that of Theorem 2.5 once one observes that for $z=a+\mathrm{i} b$, with $a\leq 0$ and $b\in\mathbb{R}$, t>0 and $x\in\mathbb{R}$ we have

$$e^{zt} = -\frac{1}{\pi} \int_{-\infty}^{\infty} e^{itx} \Re\left[\frac{1}{z - ix}\right] dx = \frac{1}{\pi} \int_{-\infty}^{\infty} e^{itx} \frac{-a}{a^2 + (b - x)^2} dx.$$

Therefore, by Fubini's theorem

$$\operatorname{cov}(f(\xi_0), f(\xi_t)) = \int_{\Re(z) \le 0} e^{zt} \nu(\mathrm{d}z) = \int_{x = -\infty}^{\infty} e^{\mathrm{i}tx} \left[\frac{-1}{\pi} \int_{\Re(z) \le 0} \Re\left(\frac{1}{z - \mathrm{i}x}\right) \nu(\mathrm{d}z) \right] \mathrm{d}x.$$

By letting $z=w{\rm i}$, where $w\in \mathbb{H}$, and the conformal invariance of Brownian motion, one can immediately deduce that

 $-\frac{1}{\pi}\Re\left(\frac{1}{z-\mathrm{i}x}\right)\mathrm{d}x,$

is the harmonic measure of Brownian motion in the left half-plane started at the point z

Proof of Theorem 3.2. First observe that

$$var(S_T) = 2 \int_{s=0}^{T} (T-s) \int_{\Re(z) \le 0} e^{zx} \nu(dz) ds = 2 \int_{\Re(z) \le 0} \Re\left[\frac{e^{zT} - zT - 1}{z^2}\right] \nu(dz).$$

By splitting $S_T = \mathbb{E}_0(S_T) + S_T - \mathbb{E}_0(S_T)$ we obtain

$$\operatorname{var}(S_T) = \mathbb{E}\left[(S_T - \mathbb{E}_0(S_T))^2 \right] + \mathbb{E}\left[\mathbb{E}_0(S_T)^2 \right]$$

$$= \int_{\Re(z) \le 0} \Re\left[2 \frac{e^{zT} - zT - 1}{z^2} - \frac{|1 - e^{zT}|^2}{|z|^2} \right] \nu(\mathrm{d}z) + \int_{\Re(z) \le 0} \frac{|1 - e^{zT}|^2}{|z|^2} \nu(\mathrm{d}z)$$

$$=: I_1 + I_2.$$

Assume $-\Re(1/z) \in L_1(d\nu)$. Careful calculation shows that

$$-\Re\left(\frac{1}{z}\right) \leftarrow \Re\left[2\frac{e^{zT} - zT - 1}{z^2} - \frac{|1 - e^{zT}|^2}{|z|^2}\right] \le \frac{C|x|}{x^2 + y^2} \in L_1(d\nu),$$

and thus $I_1/T \to -\int_{\Re(z) \le 0} \Re(1/z)\nu(\mathrm{d}z)$. Fatou's lemma also shows that if $\operatorname{var}(S_T)/T$ converges then $\Re(1/z) \in L_1(\mathrm{d}\nu)$.

Next we analyze I_2/T . Notice that

$$|1 - e^{zT}|^2 = (1 - e^{Tx})^2 + 4e^{Tx}\sin^2(Ty/2),$$

and since $(1 - e^{Tx})^2 \le T|x|$ for x < 0, we have that

$$\frac{1}{T} \int_{\Re(z) \le 0} \frac{|1 - e^{zT}|^2}{|z|^2} \nu(dz) = \frac{1}{T} \int_{\Re(z) \le 0} \frac{4e^{Tx} \sin^2(Ty/2)}{x^2 + y^2} \nu(dx, dy) + o(1)$$

$$= \frac{1}{T} \int_{\Re(z) \le 0} \frac{4\sin^2(Ty/2)}{x^2 + y^2} \nu(dx, dy) + o(1).$$

For a > 0 write

$$D_a := \{ x + iy : 0 \le -x \le a, 0 \le |y| \le a \}.$$

Notice that on $D_a^{(c)}$ the integrand is less than $1/2a^2$ and therefore

$$\int_{\Re(z) \le 0} \frac{4\sin^2(Ty/2)}{x^2 + y^2} \nu(\mathrm{d}x, \mathrm{d}y) = \int_{D_2} \frac{4\sin^2(Ty/2)}{x^2 + y^2} \nu(\mathrm{d}x, \mathrm{d}y) + o(1).$$

Let

$$\epsilon_T := \int_{D_a} \left| \frac{\sin(Ty/2)}{Ty} \right| \times \frac{|x|\nu(\mathrm{d}x,\mathrm{d}y)}{x^2 + y^2} \to 0,$$

since the bounded integrand vanishes and $\mu(\mathrm{d}z) := |x|\nu(\mathrm{d}z)/(x^2+y^2)$ is a finite measure. Let $\delta_T := \max(\mathrm{e}^{-T}, \sqrt{\epsilon_T})$, so that $\delta_T > 0$ and $\epsilon_T/\delta_T \to 0$, and define

$$U_a := \{x + iy : 0 \le -x \le |y| \le a\}, \quad D_{a,T} := \{x + iy : 0 \le \delta_T |y| \le -x, |y| \le a\}$$
$$U_{a,T} := \{x + iy : 0 \le -x \le \delta_T |y|, |y| \le a\}$$

Since on $D_{a,T}$ we have $|y| \leq |x|/\delta_T$

$$\frac{1}{T} \int_{D_{0,T}} \frac{\sin^2(Ty/2)}{x^2 + y^2} d\nu \le \frac{1}{\delta_T} \int_{D_{0,T}} \frac{|\sin(Ty/2)|}{Ty} \frac{|x|}{x^2 + y^2} d\nu \le \frac{\epsilon_T}{\delta_T} \to 0,$$

and thus since $U_{a,T}=D_a-D_{a,T}$ it follows that $y^2+x^2=y^2(1+O(\delta_T^2))$ and

$$\frac{1}{T} \int_{D_a} \frac{4 \sin^2(Ty/2)}{x^2 + y^2} \nu(\mathrm{d}x, \mathrm{d}y) = \frac{1}{T} \int_{U_{0,T}} \frac{\sin^2(Ty/2)}{(y/2)^2} \nu(\mathrm{d}x, \mathrm{d}y) \times (1 + o(1)) + o(1).$$

Notice that on $U_a \cap D_{a,T}$ we have $0 \le \delta_T |y| \le |x| \le |y| \le a$ and thus

$$\frac{|x|}{x^2 + y^2} \ge \frac{|x|}{x^2 + y^2} \ge \frac{\delta_T |y|}{2y^2} = \frac{\delta_T}{2|y|} \Longrightarrow \frac{1}{|y|} \le \frac{1}{\delta_T} \frac{2|x|}{x^2 + y^2}$$

Since $U_{a,T} = U_a - U_a \cap D_{a,T}$, from the above

$$\frac{1}{T} \int_{U_{a,T}} \frac{\sin^2(Ty/2)}{(y/2)^2} \nu(\mathrm{d} x, \mathrm{d} y) = \frac{1}{T} \int_{U_a} \frac{\sin^2(Ty/2)}{(y/2)^2} \nu(\mathrm{d} x, \mathrm{d} y) + o(1).$$

Therefore

$$\begin{split} \frac{1}{T} \int_{U_a \cap D_{a,T}} \frac{\sin^2(Ty/2)}{(y/2)^2} \nu(\mathrm{d}x,\mathrm{d}y) &\leq C \int_{U_a \cap D_{a,T}} \frac{|\sin(Ty/2)|}{T|y|} \times \frac{1}{|y|} \nu(\mathrm{d}x,\mathrm{d}y) \\ &\leq \frac{C}{\delta_T} \int_{U_a \cap D_{a,T}} \frac{|\sin(Ty/2)|}{T|y|} \times \frac{|x|}{x^2 + y^2} \nu(\mathrm{d}x,\mathrm{d}y) &\leq C \frac{\epsilon_T}{\delta_T} \to 0. \end{split}$$

Finally let

$$U:=\{x+\mathrm{i} y:0\leq -x\leq |y|\}.$$

From the above computation, since a>0 was arbitrary and all error terms depending on a vanish as $a\to\infty$ we have

$$\begin{split} \lim_{T \to \infty} \frac{1}{T} \int_{\Re(z) \le 0} \frac{|1 - \mathrm{e}^{zT}|^2}{|z|^2} \nu(\mathrm{d}z) &= \lim_{T \to \infty} \frac{1}{T} \int_{U_a} \frac{\sin^2(Ty/2)}{(y/2)^2} \nu(\mathrm{d}x, \mathrm{d}y) \\ &= \lim_{T \to \infty} \frac{1}{T} \int_{U} \frac{\sin^2(Ty/2)}{(y/2)^2} \nu(\mathrm{d}x, \mathrm{d}y) = \lim_{T \to \infty} \frac{1}{T} \int_{x=0}^{\infty} \frac{\sin^2(Tx/2)}{(x/2)^2} G(\mathrm{d}x), \end{split}$$

where $G(x) = \nu(U_x)$, using similar arguments to the proof of Theorem 2.6. Now we have

$$\frac{1}{T} \int_{x=0}^{\infty} \frac{\sin^2(Tx/2)}{(x/2)^2} G(\mathrm{d}x) = \frac{1}{T} \int_{x=0}^{\pi} \frac{\sin^2(Tx/2)}{\sin^2(x/2)} G(\mathrm{d}x) + O(1/T),$$

since

$$\frac{1}{T} \int_{x=0}^{\pi} \sin^2(Tx/2) \left| \frac{1}{(\frac{x}{2})^2} - \frac{1}{\sin^2(\frac{x}{2})} \right| G(\mathrm{d}x) \le \frac{1}{T} \int_{x=0}^{\pi} \frac{(x/2)^4}{\sin^2(x/2)(x/2)^2} G(\mathrm{d}x) = O(1/T).$$

The result then follows from Theorem A.

A Technical lemmas

A.1 Standard Tauberian Theorems

In order to make this paper more self-contained we state the following classical Tauberian theorems (Theorem 2.3 In [35] or Theorem 1.7. in [3], due to Feller).

Theorem A.1. Let U(x) be a monotone non-decreasing function on $[0,\infty)$ such that

$$w(x) = \int_{0^{-}}^{\infty} e^{-xu} dU(u)$$
 is finite for all $x > 0$.

Then if $\rho \geq 0$ and L is a slowly varying function, then

$$w(x) = x^{-\rho}L(x)$$
 as $x \to \infty$ iff $U(x) = x^{\rho}L(1/x)/\Gamma(\rho+1)$ as $x \to 0^+$.

We give the monotone Tauberian theorem (Theorem 2.4 in [35] or Theorem 2.4 in [3]) **Theorem A.2.** Let U(x) defined and positive on $[A, \infty)$ for some A > 0 given by

$$U(x) = \int_{A}^{x} u(y) \mathrm{d}y,$$

where u(y) is ultimately monotone. Then if $\rho \geq 0$ and L is a slowly varying function, then

$$U(x) = x^{\rho}L(x)$$
 as $x \to \infty$ implies $u(x) \sim \rho x^{\rho-1}L(x)$ as $x \to \infty$.

If $\rho > 0$ then u(x) is regularly varying.

A.2 Auxiliary Lemma for Theorem 2.13

Lemma A.3. For V and H as defined in Theorem 2.13, we have

$$\frac{H(1/x)}{2\theta V(x)} \rightarrow 1 \text{ as } x \rightarrow 0^+.$$

In particular if V(x) is slowly varying at 0, H(1/x) is slowly varying at ∞ .

Proof. Let r be the measure defined in (4.3). By Definition (4.8)

$$\mathbb{P}(\tau_1 > u) = 2\theta \int_0^1 x^u (1 - x) \mu(\mathrm{d}x) = \theta \int_0^1 (1 - y)^u y r(\mathrm{d}y).$$

We show first for any $\delta \in (0,1)$

$$\int_0^1 (1-y)^u r(\mathrm{d}y) = O(u^{\delta-3}) + \int_0^1 e^{-uy} y r(\mathrm{d}y). \tag{A.1}$$

To prove this, notice that for $u \ge 0$, $0 \le y \le 1$, for positive t, m and some C_m

$$|(1-y)^u - e^{-uy}| \le |1-y - e^{-y}|ue^{-(u-1)y}$$
 and $e^t \ge 1 + t^m/C_m$.

Then for any $\delta \in (0,1)$

$$\int_{0}^{1} e^{-uy} uy^{3} r(dy) \le C \int_{0}^{1} \frac{r(dy)}{u^{3-\delta} u^{1-\delta}} \le C u^{\delta-3} \int_{0}^{1} \frac{r(dy)}{u^{1-\delta}},$$

and after some rearrangement

$$\int_0^1 \frac{r(\mathrm{d}y)}{y^{1-\delta}} = C \int_{y=1}^\infty \frac{V(1/y)}{y^{1+\delta}} \mathrm{d}y < \infty,$$

since $V(u):=\int_u^1 r(\mathrm{d}y)/y$, is slowly varying as $u\to 0^+$. Since for $u\ge 0$

$$H(u) = \mathbb{E}\tau_1^2 I(\tau_1 < u) = 2 \int_0^u x \mathbb{P}(\tau_1 > x) dx - u^2 \mathbb{P}(\tau_1 > u),$$

by Fubini's theorem and (A.1) we derive

$$H(u) = O(u^{\delta - 1}) - \theta u^2 \int_0^1 e^{-uy} y r(\mathrm{d}y) + 2\theta \int_0^1 \frac{1 - e^{-yu}}{y} r(\mathrm{d}y) - 2\theta u \int_0^1 e^{-yu} r(\mathrm{d}y)$$

=: $O(u^{\delta - 1}) - \theta I_1(u) + 2\theta I_2(u) - 4I_3(u)$.

Also note that by integration by parts,

$$r(z) = \int_0^z r(ds) = -\int_0^z yV(dy) = \int_0^z [V(y) - V(z)]dy.$$

Then R can also be written as

$$r(z) = \int_{1/z}^{\infty} \frac{\mathrm{d}V(1/y)}{y},$$

and since V(1/y) is slowly varying as $y \to \infty$ we have by Theorem 1.6.5 in [3] that

$$\frac{r(z)}{zV(z)} \to 0$$
, as $z \to 0^+$.

Now let K be an arbitrary positive number. Since we will first take limits as $x\to 0^+$, we can assume that x is small enough so that Kx<1. Therefore splitting the integral $\int_0^1 = \int_0^{Kx} + \int_{Kx}^1$ and applying standard analysis, we derive

$$I_2(u) = \int_{y=0}^{Kx} \frac{1 - e^{-yu}}{y} r(\mathrm{d}y) + \int_{y=Kx}^{1} \frac{1 - e^{-yu}}{y} r(\mathrm{d}y)$$
$$= K \frac{r(Kx)}{Kx} + V(Kx)(1 + O(e^{-Kux})).$$

Then for u = 1/x, arbitrary K > 0, since for each fixed K

$$K \frac{r(Kx)}{KxV(x)} \to 0$$
 and $\frac{V(Kx)}{V(x)} \to 1$, as $x \to 0^+$

we have

$$\limsup_{x \to 0^+} \left| \frac{I_2(1/x)}{V(x)} - 1 \right| \le C e^{-K},$$

and since K > 0 is arbitrary $I_2(1/x)/V(x) \to 1$, as $x \to 0^+$ and u = 1/x.

Finally let again K > 0 be arbitrary. Then

$$\sup_{t>K} e^{-t}[1+t+t^2] = W_K, \quad W = \sup_{K>0} W_K$$

and notice that $W_K \to 0$ as $K \to \infty$. Then we have for fixed arbitrary K > 0, and x small enough so that Kx < 1,

$$I_1 + I_3 = \int_{y=0}^1 e^{-yu} (uy + (uy)^2) \frac{r(\mathrm{d}y)}{y} \le (1+W)u \int_{y=0}^{Kx} dr(y) + W_K \int_{y=Kx}^1 \frac{r(\mathrm{d}y)}{y}$$
$$= (1+W) \frac{r(Kx)}{x} + W_K V(Kx).$$

Therefore for u = 1/x, as above we have

$$\limsup_{x \to 0^+} \left| \frac{I_1(x^{-1}) + I_3(x^{-1})}{V(x)} \right| \le W_K,$$

and since K is arbitrary the claim follows from the fact that

$$\frac{I_1(x^{-1}) + I_3(x^{-1})}{V(x)} \to 0$$
, as $x \to 0^+$.

References

- [1] Billingsley, P.: Probability and measure. John Wiley & Sons, New York, 1995. MR-1324786
- [2] Billingsley, P.: Convergence of probability measures. John Wiley & Sons, New York, 1999. MR-1700749
- [3] Bingham, N.H., Goldie, C.M. and Teugels, J.L.: Regular Variation. *Cambridge University Press*, 1989. MR-1015093
- [4] Borodin, A.N. and Ibragimov, I.A.: Limit theorems for functionals of random walks. Trudy Mat. Inst. Steklov. 195, 1995. Transl. into English: Proc. Steklov Inst. Math. MR-1368394
- [5] Bradley, R.: Introduction to strong mixing conditions, Volume 3. Kendrick Press, 2007. MR-2325296
- [6] Brockwell, P. J. and Davis, R. A.: Time series, theory and methods. Springer, New York, 1991. MR-1093459
- [7] Cuny, C. and Lin, M.: Pointwise ergodic theorems with rate and application to the CLT for Markov chains. Ann. Inst. Henri Poincaré, Probab. Statist. 45, (2009), 710–733. MR-2548500
- [8] Deligiannidis, G. and Utev, S.: Variance of partial sums of stationary sequences. *Ann. Probab*, **41**, (2013), 3606–3616. MR-3127893
- [9] Deligiannidis, G. and Utev, S.: Asymptotic variance of the self-intersections of stable random walks using Darboux-Wiener theory. Sib. Math. J. 52, (2011), 639 650. MR-2883216
- [10] Derriennic, Y. and Lin, M.: The central limit thorem for Markov chains with normal transition operators started at a point. Probab. Theory Relat. Fields, 119, (2009), 508–528. MR-1826405
- [11] Doukhan, P., Massart, P. and Rio, E.: The functional central limit theorem for strongly mixing processes. *Ann. Inst. Henri Poincaré Probab. Statist.* **30**, (1994), 63–82. MR-1262892
- [12] Jewell, N.P. and Bloomfield, P.: Canonical correlations of past and future for time series: Definitions and theory, *Ann. Statist.* **11**, (1983), 837–847. MR-0707934
- [13] Feller, W.: An introduction to probability theory and its applications, Vol. II, 2nd ed., *John Willey & Sons*, New York, 1994. MR-0270403
- [14] Giraudo, D. and Volný, D.: A strictly stationary β -mixing process satisfying the central limit theorem but not the weak invariance principle. *Stoch. Proc. Appl.* **124**, (2014), 3769–3781. MR-3249354
- [15] Gordin, M. and Lifšic, B.: A remark about a Markov process with normal transition operator. *Third Vilnius Conf. Proba. Stat., Akad. Nauk Litovsk,* Vilnius. **1**, (1981), 147–148.
- [16] Geyer, C.J.: Practical Markov chain Monte Carlo. Stat. Sci. 7, (1992), 473-483.
- [17] Hardy, G.H. and Littlewood, J. E.: Solution of the Cesaro summability problem for power-series and Fourier series. *Mathematische Zeitschrift* **19**, (1924), 67–96. MR-1544641
- [18] Häggström, O. and Rosenthal, J.S.: On variance conditions for Markov chain CLTs. *Electron. Commun. Probab.* **12**, (2007), 454–464. MR-2365647
- [19] Holzmann, H.: The central limit theorem for stationary Markov processes with normal generator with applications to hypergroups. *Stochastics* 77, (2005), 371–380. MR-2165114
- [20] Holzmann, H.: Martingale approximations for continuous-time and discrete-time stationary Markov processes. *Stoch. Proc. Appl.* **115**, (2005), 1518–1529. MR-2158018
- [21] Ibragimov, I.A.: Some limit theorems for stationary processes. *Theory Probab. Appl.* **7**:4, (2005), 349–382. MR-0148125
- [22] Kipnis, C. and Varadhan, S.R.S.: Central limit theorem for additive functionals of reversible Markov processes and applications to simple exclusions. *Comm. Math. Phys* 104, (1986), 1–19. MR-0834478
- [23] Komorowski, T., Landim, C., and Olla, S.: Fluctuations in Markov processes: time symmetry and martingale approximation. *Springer*, Chicago, 2012. MR-2952852
- [24] Kontoyiannis, I. and Meyn, S.P.: Spectral theory and limit theorems for geometrically ergodic Markov processes. *Ann. App. Probab.* **13**, (2003), 304–362. MR-1952001
- [25] Lawler, G.: Conformally Invariant processes in the Plane, American Mathematical Society, Providence, 2005. MR-2129588

- [26] Longla, M., Peligrad M. and Peligrad C.: On the functional CLT for reversible Markov Chains with nonlinear growth of the variance. J. Appl. Probab. 49 (2012), 1091-1105. MR-3058990
- [27] Maxwell, M. and Woodroofe, M.: Central limit theorems for additive functionals of Markov chains. Ann. Probab. 28, (2000), 713-724. MR-1782272
- [28] Merlevède, F. and Peligrad, M.: Rosenthal-type inequalities for the maximum of partial sums of stationary processes and examples. Ann. Probab. 41, (2013), 914-960. MR-3077530
- [29] Olla, S.: Central limit theorems for tagged particles and for diffusions in random environments. Panoramas et Syntheses, 12, (2001) 75-100. MR-2226846
- [30] Peligrad, M. and Utev, S.: A new maximal inequality and invariance principle for stationary sequences. Ann. Probab. 33, (2005). 798-815. MR-2123210
- [31] Rio, E.: Théorie asymptotique des processus aléatoires faiblement dépendants. Mathématiques et Applications, Springer-Verlag, Berlin, 2000. MR-2117923
- [32] Rio, E.: Moment inequalities for sums of dependent random variables under projective conditions. J. Theoret. Probab. 22, (2009), 146-163. MR-2472010
- [33] Roberts, G.O. and Rosenthal, J.S.: General state space Markov chains and MCMC algorithms. Probability Surveys 1, (2004), 20-71. MR-2095565
- [34] Rudin, W.: Functional analysis. Second edition. McGraw-Hill, New York, 1991. MR-1157815
- [35] Seneta, E.: Regularly Varying Functions, Lecture Notes 508. Springer, Berlin, 1976. MR-0453936
- [36] Tierney, L.: A note on Metropolis-Hastings kernels for general state spaces. Ann. Appl. Probab. 8, (1998), 1-9. MR-1620401
- [37] Tóth, B.: Persistent random walks in random environment. Prob. Theory Rel. Fields 71, (1986), 615-625. MR-0833271
- [38] Toth, B.: Comment on a theorem of M. Maxwell and M. Woodroofe. Electron. Commun. Probab. 18, (2013), 1-4. MR-3033596
- [39] Zhao, O., Woodroofe, M. and Volný, D.: A central limit theorem for reversible processes with nonlinear growth of variance, J. Appl. Prob. 47, (2010), 1195-1202. MR-2752879

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