

Central limit theorems for sequential and random intermittent dynamical systems

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Central limit theorems for sequential and random intermittent dynamical systems.

Matthew Nicol * Andrew Török † Sandro Vaienti ‡ December 15, 2015

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We establish self-norming central limit theorems for non-stationary time series arising as observations on sequential maps possessing an indifferent fixed point. These

^{*}Department of Mathematics, University of Houston, Houston Texas, USA. e-mail: <nicol@math.uh.edu>.

[†]Department of Mathematics, University of Houston, Houston Texas, USA and Institute of Mathematics of the Romanian Academy, P.O. Box 1–764, RO-70700 Bucharest, Romania. e-mail: <torok@math.uh.edu>.

[‡]Aix Marseille Université, CNRS, CPT, UMR 7332, 13288 Marseille, France and Université de Toulon, CNRS, CPT, UMR 7332, 83957 La Garde, France. e-mail:<vaienti@cpt.univ-mrs.fr>.

transformations are obtained by perturbing the slope in the Pomeau-Manneville map. We also obtain quenched central limit theorems for random compositions of these maps.

1 Introduction

In a preceding series of two papers [13], [3], we considered a few statistical properties of nonstationary dynamical systems arising by the sequential composition of (possibly) different maps. The first article [13] dealt with the Almost Sure Invariance Principle (ASIP) for the non-stationary process given by the observation along the orbit obtained by concatenating maps chosen in a given set. We choose maps in one and more dimensions which were piecewise expanding, more precisely their transfer operator (Perron-Frobenius, "PF") with respect to the Lebesgue measure was quasi-compact structure on a suitable Banach space. The ASIP was then proved by applying a recent result by Cuny and Merlevede [7], whose first step was to approximate the original process with a reverse martingale difference plus an error. The latter was essentially bounded due to the presence of a spectral gap in the PF operator on a Banach space continuously injected in L^{∞} (from now on all the L^p spaces will be with respect to the ambient Lebesgue measure m and they will be denoted with L^p or $L^p(m)$.). Moreover, the same spectral property allowed us to show that for expanding maps chosen close enough, the variance σ_n^2 grows linearly, which permit to approximate the original process almost everywhere with a finite sum of i.i.d. Gaussian variables with the same variance.

The second paper [3] considered composition of Pomeau-Manneville like maps, obtained by perturbing the slope at the indifferent fixed point 0. We got polynomial decay of correlations for particular classes of centered observables, which could also be interpreted as the decay of the iterates of the PF operator on functions of zero (Lebesgue) average, and this fact is better known as loss of memory. In this situation the PF operator is not quasicompact and although the process given by the observation along a sequential orbit can be decomposed again as the sum of a reverse martingale difference plus an error, apriori the latter turns out to be bounded only in L^1 and this was an obstacle to obtain an almost sure result like the ASIP by only looking at the almost sure convergence of the reverse martingale difference. Instead one could hope to get a (distributional) Central Limit Theorem (CLT); in this regard a general approach to CLT for sequential dynamical systems has been proposed and developed in [6]. It basically applies to systems with a quasi-compact PF

operator and it is not immediately transposable to maps with do not admit a spectral gap. The main goal of our paper is to prove the CLT for the sequential composition of Pomeau-Manneville maps with varying slopes. A fundamental tool in obtaining such a result will be the polynomial loss of memory bound obtained in [3]; we are now going to recall it also because it will determine the regularity of the observables to which our CLT will apply; see Theorem 1.2.

We consider the family of Pomeau-Manneville maps

$$T_{\alpha}(x) = \begin{cases} x + 2^{\alpha} x^{1+\alpha}, & 0 \le x \le 1/2\\ 2x - 1, & 1/2 \le x \le 1 \end{cases}$$
 $0 < \alpha < 1.$ (1.1)

Actually in [3] we considered a slightly different family of this type, but pointed out that both versions could be worked out with the same techniques (see [1]), and lead to the same result; here we prefer to use the *classical* version (1.1). As in [18], we identify the unit interval [0,1] with the circle S^1 , so that the maps become continuous. If $0 < \beta_k < 1$ are given, denote by P_{β_k} or P_k the *Perron-Frobenius* operator associated with the map $T_k = T_{\beta_k}$ w.r.t. the measure m, where $0 < \beta_k \le \alpha$. For concatenations we use equivalently the notations

$$\mathcal{T}_m^{n-m+1} := T_{\beta_n} \circ T_{\beta_{n-1}} \circ \cdots \circ T_{\beta_m} = T_n \circ T_{n-1} \circ \cdots \circ T_m.$$

$$\mathcal{P}_m^{n-m+1} := P_{\beta_n} \circ P_{\beta_{n-1}} \circ \cdots \circ P_{\beta_m} = P_n \circ P_{n-1} \circ \cdots \circ P_m.$$

$$\mathcal{P}^n := \mathcal{P}_1^n \qquad \mathcal{T}^n := \mathcal{T}_1^n$$

where the exponent denotes the number of maps in the concatenation. We use for simplicity $\mathcal{T}^{\infty} := \cdots T_n \circ \cdots \circ T_1$ for a given sequence of transformations.

The Perron-Frobenius operator P_k associated to T_k satisfies the duality relation

$$\int_{M} P_{k} f g dm = \int_{M} f g \circ T_{k} dm, \text{ for all } f \in L^{1}, g \in L^{\infty}$$

and this is preserved under concatenation.

We next consider [18, 3] the cone C_2 of functions given by (here X(x) = x is the identity function):

$$\mathcal{C}_2 := \{ f \in C^0((0,1]) \cap L^1(m) \mid f \geq 0, \text{ } f \text{ decreasing, } X^{\alpha+1}f \text{ increasing, } f(x) \leq ax^{-\alpha} \ m(f) \}^1$$

¹By "decreasing" we mean "nonincreasing".

Remark 1.1 Some coefficients that appear later depend on the value a that defines the cone C_2 ; however, we will not write explicitly this dependence.

Fix $0 < \alpha < 1$; as proven in [3], provided a is large enough, the cone C_2 is preserved by all operators P_{β} , $0 < \beta \leq \alpha < 1$. The following polynomial decay result holds:

Theorem 1.2 ([3]) Suppose ψ, φ are in C_2 with equal expectation $\int \varphi dm = \int \psi dm$. Then for any $0 < \alpha < 1$ and for any sequence $T_{\beta_1}, \dots, T_{\beta_n}, n \ge 1$, of maps of Pomeau-Manneville type (1.1) with $0 < \beta_k \le \alpha < 1$, $k \in [1, n]$, we have

$$\int |P_{\beta_n} \circ \dots \circ P_{\beta_1}(\varphi) - P_{\beta_n} \circ \dots \circ P_{\beta_1}(\psi)| dm \le C_{\alpha}(\|\varphi\|_1 + \|\psi\|_1) n^{-\frac{1}{\alpha} + 1} (\log n)^{\frac{1}{\alpha}}, \quad (1.2)$$

where the constant C_{α} depends only on the map T_{α} , and $\|\cdot\|_1$ denotes the L^1 norm.

A similar rate of decay holds for observables φ and ψ that are C^1 on [0,1]; in this case the rate of decay has an upper bound given by

$$C_{\alpha} \mathcal{F}(\|\varphi\|_{C^{1}} + \|\psi\|_{C^{1}}) n^{-\frac{1}{\alpha}+1} (\log n)^{\frac{1}{\alpha}}$$

where the function $\mathcal{F}: \mathbb{R} \to \mathbb{R}$ is affine.

For the proof of the CLT Theorem 3.1 we need better decay than in L^1 . In this paper we improve the above result to decay in L^p , provided α is small enough.

Note that $\mathcal{P}^n \varphi \in \mathcal{C}_2$ if $\varphi \in \mathcal{C}_2$ and $m(\mathcal{P}^n \varphi) = m(\varphi)$, so

$$|[\mathcal{P}^n(\varphi) - \mathcal{P}^n(\psi)]|_{x}| \le |\mathcal{P}^n(\varphi)|_{x}| + |\mathcal{P}^n(\psi)|_{x}| \le am(\varphi)x^{-\alpha} + am(\psi)x^{-\alpha}$$

Proposition 1.3 Under the assumptions on Theorem 1.2, if $1 \le p < 1/\alpha$ then

$$||P_{\beta_n} \circ \dots \circ P_{\beta_1}(\varphi) - P_{\beta_n} \circ \dots \circ P_{\beta_1}(\psi)||_{L^p(m)} \le C_{\alpha,p}(||\varphi||_1 + ||\psi||_1)n^{1 - \frac{1}{p\alpha}} (\log n)^{\frac{1}{\alpha} \frac{1 - \alpha p}{p - \alpha p}}$$
(1.3)

where the constant $C_{\alpha,p}$ depends only on the map T_{α} and p.

As in Theorem 1.2, a similar L^p -decay result also holds for observables $\varphi, \psi \in C^1([0,1])$.

Proof For functions in the cone C_2 , Theorem 1.2 gives L^1 -decay; then Lemma 2.7 together with the preceding discussion implies L^p -decay for α small enough. Note that we use this Lemma with $K = 2a(\|\varphi\|_1 + \|\psi\|_1)$ and the L^1 -bound given by the Theorem, and then the coefficient in the L^p -bound is proportional to $(\|\varphi\|_1 + \|\psi\|_1)$ as well.

To prove the decay for C^1 observables, we use Lemma 2.4 (same approach as in the proof of Theorem 1.2).

Note that the convergence of the quantity (1.2) implies the decay of the non-stationary correlations with respect to m:

$$\left\| \int \psi \varphi \circ T_{\beta_n} \circ \cdots \circ T_{\beta_1} dm - \int \psi dm \int \varphi \circ T_{\beta_n} \circ \cdots \circ T_{\beta_1} dm \right\|$$

$$\leq \|\varphi\|_{\infty} \left\| P_{\beta_n} \circ \cdots \circ P_{\beta_1} (\psi) - P_{\beta_n} \circ \cdots \circ P_{\beta_1} \left(\mathbf{1} \left(\int \psi dm \right) \right) \right\|_{1}$$

provided φ is essentially bounded and $(\int \psi dm)\mathbf{1}$ is in the functional space where the convergence of (1.2) takes place. In particular, this holds for C^1 observables, by Theorem 1.2.

From now on we will take our observables as C^1 functions on the interval [0,1] and for any $\varphi \in C^1$, we will consider the following observation along a sequential orbit:

$$\varphi_k = [\varphi]_k := \varphi - \int \varphi(T_k \circ \cdots \circ T_1 x) dm.$$

As it is suggested by the preceding loss of memory result, centering the observable is the good way to define the process when it is not stationary, in order to consider limit theorems. Conze and Raugi [6] defined the sequence of transformations \mathcal{T}^{∞} to be *pointwise ergodic* whenever the law of large numbers is satisfied, namely

$$\lim_{n\to\infty} \frac{1}{n} \sum_{k=1}^{n} \left[\varphi(T_k \circ \cdots \circ T_1 x) - \int \varphi(T_k \circ \cdots \circ T_1 x) dm \right] = 0 \text{ for Lebesgue-a.e. } x.$$

We will prove in Theorem 2.10 that such a law of large numbers holds for our observations provided $0 < \alpha < 1$. It is therefore natural to ask about a non-stationary Central Limit Theorem for the sums

$$S_n := \sum_{k=1}^n [\varphi]_k \circ T_k \circ \cdots \circ T_1 \tag{1.4}$$

for a given sequence $\mathcal{T}^{\infty} := \cdots \circ T_n \circ \cdots \circ T_1$: this will be the content of the next sections.

To be more specific we will prove in Theorem 3.1 a non-stationary central limit theorem similar to that proved by Conze and Raugi [6] for (piecewise expanding) sequential systems:

$$\frac{S_n}{\sqrt{\operatorname{Var}(S_n)}} \to^d \mathcal{N}(0,1). \tag{1.5}$$

At this point, we would like to make a few comments about our result compared to that of Conze and Raugi. Theorem 5.1 in [6] shows that, when applied to the quantities defined above and for classes of maps enjoying a quasi-compact transfer operator:

- (1) If the norms $||S_n||_2$ are bounded, then the sequence $S_n, n \ge 1$ is bounded.
- (2) If $||S_n||_2 \to \infty$, then (1.5) holds.

We are not able to prove item (1) for the intermittent map following the same approach as in [6], since it uses the uniform boundedness of the sequence $\mathbf{H}_n \circ \mathcal{T}^k$, where the function \mathbf{H}_n is defined in (2.1) and is just the error in the martingale approximation as we discussed above. We can only prove that \mathbf{H}_n is bounded uniformly in n on each set of the form [a, 1), a > 0, and do not expect it to be bounded near 0 (look at the stationary case).

Instead, our central limit theorem will satisfy item (2) under the assumption that the variance $||S_n||_2$ grows at a certain rate and for some limitation on the range of values of α . It seems difficult to get such a result in full generality for the intermittent map considered here. Conze and Raugi proved the linear growth of the variance in their Theorem 5.3 under a certain number of assumptions, including the presence of a spectral gap for the transfer operator. We showed in our paper [13] that those assumptions apply to several classes of expanding maps even in higher dimensions.

However, for concatenations given by the same intermittent map T_{α} with $\alpha < 1/2$, the variance is linear in n, provided the observable is not a coboundary for T_{α} . In section 4 we prove that the linear growth of the variance still holds if we take maps T_{β_n} with β_n arbitrary but close to a fixed β , and an observable is not a coboundary for T_{β} ; therefore, the CLT holds. See Theorem 4.1. Our proof of Theorem 4.1 uses an estimate of interesting related work of Leppänen and Stenlund [16], which we learnt about after a first version of this paper was completed. Their result allowed us to give another example where variance grows linearly for a sequential dynamical system of intermittent type maps, and hence the non-stationary CLT holds. The focus of [16] is however more on the strong law of large numbers and convergence in probability rather than the CLT. They also consider quasi static systems, introduced in [17].

In section 5 we show that the variance grows linearly for almost all sequences when we compose intermittent maps chosen from a finite set and we take them according to a fixed probability distribution. This means that for almost all sequences (with respect to the induced Bernoulli measure) of maps, the central limit theorem holds (a *quenched*-like CLT). See Theorem 5.2.

Remark 1.4 For simplicity, in many of the following statements we will use as rate of

decay $n^{-\frac{1}{\alpha}+1}$, ignoring the log *n*-factor. This is correct if we take for α a slightly larger value (and is actually the correct rate of decay for the stationary case).

Notation 1.5 For any sequences of numbers $\{a_n\}$ and $\{b_n\}$, we will write $a_n \approx b_n$ if $c_1b_n \leq a_n \leq c_2b_n$ for some constants $c_2 \geq c_1 > 0$.

2 Cones and Martingales

In order to get the right martingale representation, we begin by recalling a few formulas concerning the transfer operator; the conditional expectation is considered with respect to the measure m, and \mathcal{B} denotes the Borel σ -algebra on [0,1]. We have:

$$\mathbb{E}[\varphi \mid \mathcal{T}^{-k}\mathcal{B}] = \frac{\mathcal{P}^k(\varphi)}{\mathcal{P}^k(\mathbf{1})} \circ \mathcal{T}^k$$

$$P(\varphi \circ T \cdot \psi) = \varphi \cdot P(\psi)$$

and therefore, for $0 \le \ell \le k$

$$\mathbb{E}[\varphi \circ \mathcal{T}^{\ell} \mid \mathcal{T}^{-k}\mathcal{B}] = \frac{\mathcal{P}_{\ell+1}^{k-\ell}(\varphi \cdot \mathcal{P}^{\ell}(\mathbf{1}))}{\mathcal{P}^{k}(\mathbf{1})} \circ \mathcal{T}^{k}.$$

Recall that for $L^2(m)$ -functions these conditional expectations are the orthogonal projections in $L^2(m)$.

We denote as above: $\varphi - m(\varphi \circ \mathcal{T}^j)$ by φ_j or $[\varphi]_j$. However, to simplify notation, it is convenient to assume that $\varphi_0 = [\varphi]_0 = 0$. Therefore we have for the centered sum (1.4): $S_n = \sum_{k=1}^n \varphi_k \circ \mathcal{T}^k = \sum_{k=0}^n \varphi_k \circ \mathcal{T}^k$.

Introduce

$$\mathbf{H}_n \circ \mathcal{T}^n := \mathbb{E}(S_{n-1} \mid \mathcal{T}^{-n}\mathcal{B}).$$

Hence $\mathbf{H}_1 = 0$, and the explicit formula for \mathbf{H}_n is

$$\mathbf{H}_{n} = \frac{1}{\mathcal{P}^{n} \mathbf{1}} \left[P_{n}(\varphi_{n-1} \mathcal{P}^{n-1} \mathbf{1}) + P_{n} P_{n-1}(\varphi_{n-2} \mathcal{P}^{n-2} \mathbf{1}) + \dots + P_{n} P_{n-1} \dots P_{1}(\varphi_{0} \mathcal{P}^{0} \mathbf{1}) \right].$$

$$(2.1)$$

It is not hard to check that setting

$$S_n = M_n + \mathbf{H}_{n+1} \circ \mathcal{T}^{n+1}$$

the sequence $\{M_n\}$ is a reverse martingale for the decreasing filtration $\{\mathcal{B}_n := \mathcal{T}^{-n}\mathcal{B}\}$:

$$\mathbb{E}(M_n \mid \mathcal{B}_{n+1}) = 0.$$

In particular,

$$M_n - M_{n-1} = \psi_n \circ \mathcal{T}^n$$
 with $\psi_n := \varphi_n + \mathbf{H}_n - \mathbf{H}_{n+1} \circ T_{n+1}$. (2.2)

We recall three lemmas from [14], stated in the current context:

Lemma 2.1 ([14, Lemma 2.6])

$$\sigma_n^2 := \mathbb{E}[(\sum_{i=1}^n \varphi_i \circ \mathcal{T}^i)^2] = \sum_{i=1}^n \mathbb{E}[\psi_i^2 \circ \mathcal{T}^i] - \int \mathbf{H}_1^2 + \int \mathbf{H}_{n+1}^2 \circ \mathcal{T}^{n+1}$$

(and $\mathbf{H}_1 = 0$).

To prove this Lemma we replace our \mathbf{H}_n with ω_n in [14].

Lemma 2.2 ([14, proof of Lemma 3.3]) Let $\mathbf{H}_{j}^{\varepsilon} = \mathbf{H}_{j} \mathbf{1}_{\{|\mathbf{H}_{j}| \leq \varepsilon \sigma_{n}\}}$, where for simplicity of notation we have left out the dependence on n. Then

$$\int \left(\sum_{j=1}^{n} \psi_{j} \circ \mathcal{T}^{j} \cdot \mathbf{H}_{j+1}^{\varepsilon} \circ \mathcal{T}^{j+1}\right)^{2} = \sum_{j=1}^{n} \int \left(\psi_{j} \circ \mathcal{T}^{j} \cdot \mathbf{H}_{j+1}^{\varepsilon} \circ \mathcal{T}^{j+1}\right)^{2}$$

The last formula in the proof of [14, Lemma 2.6] equivalently gives:

Lemma 2.3

$$\sigma_n^2 = \sum_{i=1}^n \mathbb{E}[\varphi_i^2 \circ \mathcal{T}^i] + 2\sum_{i=1}^n \mathbb{E}[(\mathbf{H}_i \varphi_i) \circ \mathcal{T}^i]$$

The following Lemma plays a crucial role all along this paper. In a slightly different form it was introduced and used in [18, Sect. 4], without a proof, and subsequently in [3]. We now give a detailed proof in a more general setting.

Lemma 2.4 Assume given a C^1 -function $\varphi : [0,1] \to \mathbb{R}$ and $h \in \mathcal{C}_2$. where the cone \mathcal{C}_2 is defined with a > 1.

Denote by X the function X(x) = x. If

$$\lambda \leq -|\varphi'|_{\infty}$$

$$\nu \geq -|\varphi + \lambda X|_{\infty}$$

$$\delta \geq \frac{a}{\alpha + 1} (|\varphi'|_{\infty} + |\lambda|) m(h)$$

$$\delta \geq \frac{a}{a - 1} |\varphi + \lambda X + \nu|_{\infty} m(h)$$

then

$$(\varphi + \lambda X + \nu)h + \delta \in \mathcal{C}_2.$$

Remark 2.5 It follows immediately that if $\varphi \in C^1([0,1])$ and $h \in \mathcal{C}_2$ then we can use Theorem 1.2 and Proposition 1.3 to obtain decay of $\mathcal{P}^{\ell}(\varphi h - m(\varphi h))$: consider $\Phi := (\varphi + \lambda X + \nu)h + \delta$, $\Psi := (\lambda X + \nu)h + \delta + m(\varphi h)$, with constants chosen according to Lemma 2.4 so that $\Phi, \Psi \in \mathcal{C}_2$ (by definition, $m(\Phi) = m(\Psi)$), and write

$$\mathcal{P}^{\ell}(\varphi \cdot h - m(\varphi \cdot h)) = \mathcal{P}^{\ell}(\Phi - \Psi).$$

Corollary 2.6 In particular, for a sequence $\omega_k \in C^1([0,1])$ with $\|\omega_k\|_{C^1} \leq K$ and $h_k \in C_2$ with $m(h_k) \leq M$ (e.g, $h_k := \mathcal{P}^k(\mathbf{1})$), one can choose constants λ , ν and δ so that

$$(\omega_k + \lambda X + \nu)h_k + \delta, (\lambda X + \nu)h_k + \delta + m(\omega_k h_k) \in \mathcal{C}_2$$
 for all $k \ge 1$

and therefore

$$||\mathcal{P}^n(\omega_k h_k - m(\omega_k h_k))||_1 \le C_{\alpha,K,M} n^{-\frac{1}{\alpha}+1} (\log n)^{\frac{1}{\alpha}} \quad \text{for all } n \ge 1, \ k \ge 1,$$

where the constant $C_{\alpha,K,M}$ has an explicit expression in terms of α, K and M. Decay in L^p now follows from Lemma 2.7: if $1 \le p < 1/\alpha$ then

$$||\mathcal{P}^n(\omega_k h_k - m(\omega_k h_k))||_p \le C_{\alpha,K,M,p} n^{-\frac{1}{p\alpha}+1}$$
 for all $n \ge 1, k \ge 1$

(ignoring the log-correction, see Remark 1.4) where the constant on the right hand side depends now upon p too.

Proof of Lemma 2.4 Denote $\Phi := (\varphi + \lambda X + \nu)h + \delta$. There are three conditions for Φ to be in C_2 .

 $\underline{\Phi}$ nonnegative and decreasing. If $\lambda \leq -\sup \varphi'$ and $\nu \geq -\inf(\varphi + \lambda X)$ then $\varphi + \lambda X + \nu$ is decreasing and nonnegative. Therefore Φ , is also decreasing (because $h \in \mathcal{C}_2$) and nonnegative provided $\delta \geq 0$.

 $\Phi X^{1+\alpha}$ increasing. For $0 < x < y \le 1$, need

$$\left[(\varphi(x) + \lambda x + \nu)h(x) + \delta \right] x^{1+\alpha} \le \left[(\varphi(y) + \lambda y + \nu)h(y) + \delta \right] y^{1+\alpha}$$

$$\iff \left[\varphi(x) + \lambda x + \nu \right] \le \left[\varphi(y) + \lambda y + \nu \right] \frac{h(y)}{h(x)} \frac{y^{\alpha+1}}{x^{\alpha+1}} + \delta \left[\frac{y^{\alpha+1}}{x^{\alpha+1}} - 1 \right] \frac{1}{h(x)}$$

Since $hX^{\alpha+1} \geq 0$ is increasing, $1 \leq \frac{h(y)}{h(x)} \frac{y^{\alpha+1}}{x^{\alpha+1}}$, so it suffices to have

$$\begin{split} \varphi(x) + \lambda x + \nu &\leq \left[\varphi(y) + \lambda y + \nu \right] + \delta \left[\frac{y^{\alpha + 1}}{x^{\alpha + 1}} - 1 \right] \frac{1}{h(x)} \\ \iff & \delta \geq - \left[\left(\varphi(y) + \lambda y + \nu \right) - \left(\varphi(x) + \lambda x + \nu \right) \right] \frac{h(x)}{\frac{y^{\alpha + 1}}{x^{\alpha + 1}} - 1}. \end{split}$$

By the mean value theorem and using that $\alpha \leq 1$, $y^{\alpha+1} - x^{\alpha+1} = (\alpha+1)\xi^{\alpha}(y-x) \geq (\alpha+1)x^{\alpha}(y-x) \geq (\alpha+1)x(y-x)$; therefore

$$0 \le \frac{h(x)}{\frac{y^{\alpha+1}}{x^{\alpha+1}} - 1} = \frac{h(x)x^{\alpha+1}}{y^{\alpha+1} - x^{\alpha+1}} \le \frac{h(x)x^{\alpha}}{(\alpha+1)(y-x)} \le \frac{am(h)}{(\alpha+1)(y-x)}.$$

Meanwhile,

$$-\left[\left(\varphi(y) + \lambda y + \nu\right) - \left(\varphi(x) + \lambda x + \nu\right)\right] \le \left(\left|\varphi'\right|_{\infty} + |\lambda|\right)(y - x).$$

Using these in the above lower bound for δ , we conclude that it suffices to have

$$\delta \geq \frac{a}{\alpha + 1} (|\varphi'|_{\infty} + |\lambda|) m(h)$$

 $\underline{\Phi X^{\alpha} \leq am(\Phi)}$. Using that $hX^{\alpha} \leq am(h)$,

$$[(\varphi + \lambda X + \nu)h + \delta]X^{\alpha} \le (\varphi + \lambda X + \nu)hX^{\alpha} + \delta \le \sup(\varphi + \lambda X + \nu)am(h) + \delta.$$

On the other hand, $am((\varphi + \lambda X + \nu)h + \delta) \ge a\inf(\varphi + \lambda X + \nu)m(h) + a\delta$, so it suffices to have

$$\sup(\varphi + \lambda X + \nu)am(h) + \delta \le a\inf(\varphi + \lambda X + \nu)m(h) + a\delta$$

$$\iff \delta \ge \frac{a}{a - 1} \left[\sup(\varphi + \lambda X + \nu) - \inf(\varphi + \lambda X + \nu)\right]m(h).$$

Note that, since the transfer operators are monotone,

$$|P_n \dots P_{k+1}[\varphi \mathcal{P}^k \mathbf{1}]|_x \le P_n \dots P_{k+1}[|\varphi|_\infty \mathcal{P}^k \mathbf{1}]|_x = |\varphi|_\infty P_n \dots P_{k+1}[\mathcal{P}^k \mathbf{1}]|_x$$

Since $|\varphi|_{\infty}P_n \dots P_{k+1}[\mathcal{P}^k \mathbf{1}]$ lies in the cone \mathcal{C}_2 this implies that

$$|P_n \dots P_{k+1}[\varphi \mathcal{P}^k \mathbf{1}]| \mid_x \le a|\varphi|_{\infty} x^{-\alpha}.$$

The following Lemma gives control over the L^p -norm of functions with such a bound.

Lemma 2.7 Suppose that $f \in L^1(m)$ and $|f(x)| \leq Kx^{-\alpha}$. Then, provided $p \geq 1$ and $\alpha p < 1$,

$$||f||_p \le C_{\alpha,p} ||f||_1^{\frac{1-\alpha p}{p-p\alpha}} K^{\frac{p-1}{p-p\alpha}}$$

In particular, if $|f(x)| \le Kx^{-\alpha}$ and $||f||_1 \le Mn^{1-\frac{1}{\alpha}}$, then

$$||f||_p \le C_{K,M,\alpha,p} n^{1-\frac{1}{p\alpha}} \text{ for } 1 \le p < 1/\alpha.$$

Therefore, for $1 \le p < 1/(2\alpha)$, there is $\delta > 0$ such that $||f||_p \le C_{K,M,\alpha,p} n^{-1-\delta}$.

Proof The case p=1 is obviously true, so we assume from now on that p>1. Denote $C_1:=||f||_1$. Compute, for $0 < x_* \le 1$, and $\alpha p < 1$: $\int_{x_*}^1 |f|^p dx \le \sup\{|f(x)|^{p-1} \mid x_* \le x \le 1\} \int_0^1 |f| dx \le K^{p-1} x_*^{-\alpha(p-1)} C_1$, and $\int_0^{x_*} |f|^p dx \le K^p \int_0^{x_*} x^{-\alpha p} dx = \frac{K^p}{1-\alpha p} x_*^{1-\alpha p}$. We want to minimize over x_* the quantity

$$G(x_*) := K^{p-1}C_1x_*^{-\alpha(p-1)} + K^p \frac{1}{1-\alpha p}x_*^{1-\alpha p} = Ax_*^{-\alpha(p-1)} + Bx_*^{1-\alpha p}.$$

It reaches its minimum value for $x_*^{\alpha-1} = \frac{B(1-\alpha p)}{A\alpha(p-1)}$, which gives for the minimum of $G^{1/p}$ the value

$$C_{\alpha,p}C_1^{\frac{1-\alpha p}{1-\alpha}\frac{1}{p}}K^{\frac{p-1}{p}\frac{1}{1-\alpha}}.$$

For the last statement notice that $\frac{1-p\alpha}{p\alpha} > 1 \iff 0 < \alpha p < 1/2$.

Corollary 2.8 We have:

1. $||\mathbf{H}_n||_q$ is uniformly bounded in n for $1 \le q < \frac{1}{2\alpha}$.

2. $||\mathbf{H}_n \circ \mathcal{T}^n||_r$ is uniformly bounded in n for $1 \le r < \frac{1}{2\alpha} - \frac{1}{2}$.

Proof Recall that \mathbf{H}_n is given in (2.1). By [3, Remark 1.3], $\mathcal{P}^n(\mathbf{1}) \geq D_\alpha > 0$ on (0,1]. We now apply Minkowski's inequality in the sum defining \mathbf{H}_n . Thanks to Lemma 2.7 each term of the form $P_n P_{n-1} \dots P_{n-\ell}(\varphi_{n-\ell-1}\mathcal{P}^{n-\ell-1}1)$, $\ell \in [0, n-1]$ will be bounded in L^p by $\frac{2}{D_\alpha} C_{\alpha,K,p} \ell^{1-\frac{1}{p\alpha}}$, where K is the C^1 norm of φ . The role of h_k in Lemma 2.6 is now played by $\mathcal{P}^{n-\ell-1}\mathbf{1}$ and therefore M=1. By summing over ℓ from 1 to infinity, we get a convergent series whenever $p\alpha < 1/2$. We now write $\int |\mathbf{H}_n \circ \mathcal{T}^n|^r dx = \int |\mathbf{H}_n|^r \mathcal{P}^n \mathbf{1} dx$. Since $\mathcal{P}^n \mathbf{1}$ belongs to $L^p(m)$ for $1 \leq p < \frac{1}{\alpha}$ by the definition of \mathcal{C}_2 and its invariance property, the function $|\mathbf{H}_n|^{r\frac{p}{p-1}}$ must be uniformly in $L^1(m)$ and therefore, by the previous item, $r\frac{p}{p-1} < \frac{1}{2\alpha}$. Thus we need $1 \leq r < \frac{p-1}{2p\alpha}$ for some $1 \leq p < \frac{1}{\alpha}$, which means $1 \leq r < \frac{1}{2\alpha} - \frac{1}{2}$.

As we said in the Introduction, we will also have a pointwise bound on the \mathbf{H}_n 's.

Lemma 2.9 For $0 < \alpha < 1/2$, there is a constant C depending on α and $K = ||\varphi||_{C^1}$, such that

$$|\mathbf{H}_n(x)| \le Cx^{-\alpha - 1}$$
 for all $x \in (0, 1], n \ge 1.$ (2.3)

Proof By using again formula (2.1) for \mathbf{H}_n (where $\varphi_0 = 0$) and the bound $\mathcal{P}^n(\mathbf{1}) \geq D_\alpha > 0$ we are left with the pointwise estimate of

$$P_n(\varphi_{n-1}\mathcal{P}^{n-1}\mathbf{1}) + P_nP_{n-1}(\varphi_{n-2}\mathcal{P}^{n-2}\mathbf{1}) + \dots + P_nP_{n-1}\dots P_1(\varphi_0\mathcal{P}^0\mathbf{1}).$$

By Corollary 2.6, for each $k \geq 1$ one can write $\varphi_k \mathcal{P}^k \mathbf{1} = (\varphi - m(\varphi \circ \mathcal{T}^k)) \mathcal{P}^k \mathbf{1} = A_k - B_k$ where $A_k, B_k \in \mathcal{C}_2$ with $m(A_k), m(B_k)$ uniformly bounded by some constant $C_{\alpha,K} < \infty$. Therefore, by the decay Theorem 1.2 (and ignoring the log-correction), there is a new constant C' depending only on α and K such that

$$\|\mathcal{P}_{k+1}^{n-k}(A_k - B_k)\|_1 \le C'(n-k)^{-\frac{1}{\alpha}+1}.$$
 (2.4)

We now recall the footnote to the proof of [18, Lemma 2.3]: if $f \in \mathcal{C}_2$ with $m(f) \leq M$ then

$$|x^{\alpha+1}f(x) - y^{\alpha+1}f(y)| \le a(1+\alpha)M|x-y| \text{ for } 0 < x, y \le 1.$$
 (2.5)

But a bound $|g(x) - g(y)| \le L|x - y|$ for the Lipschtz-seminorm $|g|_{\text{Lip}}$ implies

$$||g||_1 \ge C_L ||g||_{\infty}. \tag{2.6}$$

Combining the above observations and since $m(\mathcal{P}_{k+1}^{n-k}(f)) = m(f)$, we obtain that $|X^{\alpha+1}\mathcal{P}_{k+1}^{n-k}(A_k - B_k)|_{\text{Lip}} \leq |X^{\alpha+1}\mathcal{P}_{k+1}^{n-k}(A_k)|_{\text{Lip}} + |X^{\alpha+1}\mathcal{P}_{k+1}^{n-k}(B_k)|_{\text{Lip}} \leq L$ uniformly for $n \geq 1, 1 \leq k < n$, and then

$$||X^{\alpha+1}\mathcal{P}_{k+1}^{n-k}(A_k - B_k)||_{\infty} \le 1/C_L ||X^{\alpha+1}\mathcal{P}_{k+1}^{n-k}(A_k - B_k)||_1 \le C''(n-k)^{-\frac{1}{\alpha}+1}$$

for a new constant C'' depending only on α, K, L , which implies that

$$|\mathcal{P}_{k+1}^{n-k}(A_k - B_k)(x)| \le x^{-\alpha - 1}C''(n-k)^{-\frac{1}{\alpha} + 1}$$

and therefore, for $0 < \alpha < 1/2$,

$$\left| \sum_{k=1}^{n-1} \mathcal{P}_{k+1}^{n-k} (A_k - B_k)(x) \right| \le x^{-\alpha - 1} C'' \sum_{k=1}^{n-1} (n-k)^{-\frac{1}{\alpha} + 1} \le C x^{-\alpha - 1}$$

as desired.

We finish this Section by proving a type of Borel-Cantelli Lemma which is an unavoidable tool in proving non-stationary limit theorems.

Theorem 2.10 (Strong Borel-Cantelli) Suppose that for $j \geq 1$, $\psi_j \in C^1([0,1])$ with uniformly bounded C^1 -norms.

(a) If
$$0 < \alpha < 1/2$$
 then

$$\sum_{j=1}^{n} \psi_j(\mathcal{T}^j) - \sum_{j=1}^{n} m(\psi_j(\mathcal{T}^j)) = O(n^{1/2} (\log \log n)^{3/2}) \quad m\text{-a.e.}$$

and therefore, if $\liminf_{j} m(\psi_{j} \circ \mathcal{T}^{j}) > 0$ then

$$\frac{\sum_{j=1}^{n} \psi_j(\mathcal{T}^j x)}{\sum_{j=1}^{n} m(\psi_j \circ \mathcal{T}^j)} \to 1 \quad m\text{-}a.e. \ x.$$

(b) If $0 < \alpha < 1$ then

$$\frac{1}{n} \left[\sum_{j=1}^{n} \psi_j(\mathcal{T}^j x) - \sum_{j=1}^{n} m(\psi_j \circ \mathcal{T}^j) \right] \to 1 \quad m\text{-a.e. } x.$$

Proof To prove the first statement in part (a) we will use the Gál-Koksma Theorem 6.1 in the Appendix. By adding the same constant to all the ψ_j 's, we can assume without loss of generally that $\inf_j m(\psi_j \circ \mathcal{T}^j) > 0$. Thus, it suffices to give a linear upper bound for

 $\mathbb{E}[(\sum_{j=1}^n \psi_j(\mathcal{T}^j) - b_n)^2]$, where $b_n := \sum_{j=1}^n m(\psi_j(\mathcal{T}^j))$; note that the same estimate can be derived for sums over $m \le j \le n$. Expand

$$\mathbb{E}\left[\left(\sum_{j=1}^{n} \psi_{j} \circ \mathcal{T}^{j} - b_{n}\right)^{2}\right] = \sum_{j=1}^{n} \mathbb{E}\left[\psi_{j} \circ \mathcal{T}^{j} - m(\psi_{j} \circ \mathcal{T}^{j})\right]^{2} + 2\sum_{i=1}^{n} \sum_{j>i} \mathbb{E}\left[\left(\psi_{j} \circ \mathcal{T}^{j} - m(\psi_{j} \circ \mathcal{T}^{j})(\psi_{i} \circ \mathcal{T}^{i} - m(\psi_{i} \circ \mathcal{T}^{i})\right)\right]$$

and use the decay to estimate the mixed terms. Denote $\overline{\psi_j} = \psi_j - m(\psi_j \circ \mathcal{T}^j)$. Then, for j > i,

$$\begin{split} |\mathbb{E}[(\psi_{j}(\mathcal{T}^{j}) - m(\psi_{j}(\mathcal{T}^{j}))(\psi_{i}(\mathcal{T}^{i}) - m(\psi_{i}(\mathcal{T}^{i}))]| &= |\mathbb{E}[\overline{\psi}_{j} \circ \mathcal{T}^{j} \cdot \overline{\psi}_{i} \circ \mathcal{T}^{i}]| \\ &= |\mathbb{E}[(\overline{\psi}_{j} \circ \mathcal{T}_{i+1}^{j-i} \cdot \overline{\psi}_{i} \cdot \mathcal{P}^{i}(\mathbf{1})]| = |\mathbb{E}[(\overline{\psi}_{j} \cdot \mathcal{P}_{i+1}^{j-i}(\overline{\psi}_{i}\mathcal{P}^{i}(\mathbf{1}))]| \\ &\leq ||\overline{\psi}_{j}||_{\infty} ||\mathcal{P}_{i+1}^{j-i}(\overline{\psi}_{i}\mathcal{P}^{i}(\mathbf{1}))||_{1} \leq C(j-i)^{1-\frac{1}{\alpha}} \end{split}$$

where in the last inequality we used Corollary 2.6. Therefore

$$\mathbb{E}[(\sum_{j=1}^{n} \psi_{j}(\mathcal{T}^{j}) - b_{n})^{2}]$$

$$\leq 2 \sum_{i=1}^{n} |(\psi_{j}(\mathcal{T}^{i}) - m(\psi_{i}(\mathcal{T}^{i}))|_{\infty} m(\psi_{i}(\mathcal{T}^{i})) + 2C \sum_{i=1}^{n} \sum_{j>i} (j-i)^{1-\frac{1}{\alpha}} \leq nC',$$

where the constants C, C' are independent of j and n.. The conclusion now follows from the Gál-Koksma Theorem 6.1.

For (b), note that for $1/2 \le \alpha < 1$ the above computation still gives

$$\mathbb{E}\left[\left(\sum_{j=1}^{n} \psi_j(\mathcal{T}^j) - b_n\right)^2\right] \le Cn^{3 - \frac{1}{\alpha}}$$

which implies that

$$\sum_{j=1}^{n} \psi_{j}(\mathcal{T}^{j}) - b_{n} = O(n^{1-\eta}) \text{ a.s.}$$

for some $\eta > 0$, see the standard Lemma 2.11.

Lemma 2.11 Assume the random variables X_n have mean zero, and there are $M < \infty$, $\gamma < 2$ such that

$$||X_n||_{\infty} \le M$$
, $\operatorname{Var}\left(\sum_{k=1}^n X_k\right) \le Cn^{\gamma}$ for all n .

Then

$$\sum_{k=1}^{n} X_k = O(n^{\eta}) \text{ a.s. for } \eta > \frac{\gamma + 1}{3}.$$

Proof Denote $S_n := \sum_{k=1}^n X_k$. From Tchebycheff's inequality,

$$P(|S_n| > n^{1-\delta}) \le \frac{\operatorname{Var}(S_n)}{(n^{1-\delta})^2} \le Cn^{\gamma - 2\delta - 2}.$$

Pick $\delta > 0$ so that $\gamma - 2\delta - 2 < 0$ and $\omega > 0$ such that $\omega(2 - \gamma + 2\delta) > 1$. Then, for the subsequence $n_k := k^{\omega}$,

$$\sum_{k} P(|S_{n_k}| > n_k^{1-\delta}) < \infty$$

so, by Borel-Cantelli,

$$|S_{n_k}| = O(n_k^{1-\delta}) \text{ a.s.}$$
 (2.7)

Using (2.7), one has a.s.: if $n_k \le n < n_{k+1}$ for some k, then

$$|S_n| \le |S_{n_k}| + [n_{k+1} - n_k] \sup ||X_\ell||_{\infty} \le O(n_k^{1-\delta}) + Ck^{\omega - 1}M \le O(n^{1-\delta}) + C(n^{1/\omega})^{\omega - 1}M$$

therefore $|S_n| = O(n^{\eta})$ a.s. with

$$\eta = \max\left\{1 - \delta, \frac{\omega - 1}{\omega}\right\}.$$

Optimize over δ and ω to get the claimed lower bound on η .

3 Central Limit Theorem

We assume in this section that $0 < \alpha < 1/2$ (note that in the stationary case the CLT holds only in this range). With our approach we can only prove the non-stationary CLT for a lower upper bound on α , which will be stated later.

We define scaling constants $\sigma_n^2 = \mathbb{E}[(\sum_{j=1}^n \varphi_j \circ \mathcal{T}^j)^2]$. This sequence of constants play the role of non-stationary variance. As we pointed out in the Introduction, giving estimates on the growth and non-degeneracy of σ_n in this non-stationary setting is more difficult than in the usual stationary case.

Theorem 3.1 (CLT for C^1 functions) Let φ be a $C^1([0,1])$ function, and define S_n as in (1.4),

$$S_n := \sum_{k=1}^n \varphi_k \circ T_{\beta_k} \circ \cdots \circ T_{\beta_1}.$$

 $Assume\ that$

$$\sigma_n^2 := \operatorname{Var}(S_n) = \mathbb{E}[(\sum_{i=1}^n \varphi_i \circ \mathcal{T}^i)^2] \approx n^{\beta}.$$

Then, provided $\alpha < 1/8$ and $\beta > 2/3$ (see (3.6) for other cases),

$$\frac{S_n}{\sigma_n} \to^d \mathcal{N}(0,1).$$

Following the approach of Gordin we will express $S_n = \sum_{j=1}^n \varphi_j \circ \mathcal{T}^j$ as the sum of a (non-stationary) martingale difference array and a controllable error term and then use the following Theorem from Conze and Raugi [6, Theorem 5.8], which is a modification of a result of B. M. Brown [5] from martingale differences to reverse martingale differences.

Theorem 3.2 ([6, Theorem 5.8]) Let (X_i, \mathcal{F}_i) be a sequence of differences of square integrable reversed martingales, defined on a probability space $(\Omega, \mathcal{F}, \mathcal{P})$. For $n \geq 0$ let

$$S_n = X_0 + \ldots + X_{n-1}, \ \sigma_n^2 = \sum_{k=0}^{n-1} \mathbb{E}[X_k^2], \ V_n = \sum_{k=0}^{n-1} \mathbb{E}[X_k^2 | \mathcal{F}_{k+1}].$$

Assume the following two conditions hold:

- (i) the sequence of random variables $(\sigma_n^{-2}V_n)_{n\geq 1}$ converges in probability to 1.
- (ii) For each $\varepsilon > 0$, $\lim_{n \to \infty} \sigma_n^{-2} \sum_{k=0}^{n-1} \mathbb{E}[X_k^2 \mathbf{1}_{\{|X_k| > \varepsilon \sigma_n\}}] = 0$.

Then

$$\lim_{n \to \infty} \sup_{\alpha \in \mathbb{R}} \left| P \left[\frac{S_n}{\sigma_n} < a \right] - \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\alpha} e^{-\frac{x^2}{2}} dx \right| = 0.$$

Proof of Theorem 3.1

Let us take the quantity \mathbf{H}_n defined in (2.1) and then the function ψ_n given in (2.2)

$$\psi_n := \varphi_n + \mathbf{H}_n - \mathbf{H}_{n+1} \circ T_{n+1}.$$

We note that $\psi_n \circ \mathcal{T}^n$ is a reverse martingale difference scheme, uniformly bounded in $L^{r_1}(m)$, for some r_1 verifying the second item in Corollary 2.8; in particular we will take r_1 as the exponent for which $\mathbf{H}_{n+1} \circ \mathcal{T}^{n+1}$ is bounded in $L^{r_1}(m)$. That is, $1 \leq r_1 < \frac{1}{2\alpha} - \frac{1}{2}$.

We will verify conditions (i) and (ii) of Theorem 3.2. For condition (ii) we begin by noticing that the functions $\psi_n \circ \mathcal{T}^n$ have a uniformly bounded L^2 -norm if the same is true for $\mathbf{H}_{n+1} \circ \mathcal{T}_{n+1}$; this holds provided $2 < \frac{1}{2\alpha} - \frac{1}{2} \iff 0 < \alpha < \frac{1}{5}$. By Minkowski's inequality, $\|\psi_n \circ \mathcal{T}^n\|_{L^2(m)}$ will be bounded uniformly in n by some constant \hat{C} . Then we have by Hölder's and Tchebycheff's inequality

$$\sigma_n^{-2} \sum_{k=0}^{n-1} \mathbb{E}[\psi_k^2 \mathbf{1}_{\{|\psi_k| > \varepsilon \sigma_n\}}] \le \sigma_n^{-2} \hat{C} \sum_{k=0}^{n-1} m(|\psi_k| > \varepsilon \sigma_n)^{\frac{1}{2}} \le \sigma_n^{-2} \hat{C}^2 \frac{n}{\varepsilon \sigma_n}.$$

We note at this point that by prescribing a growth of the variance as $\sigma_n^2 \approx n^{\beta}$ we need $\beta > 2/3$.

The hard part lies in establishing (i). This is in contrast with the stationary setting where condition (i) is usually a straightforward consequence of the ergodic theorem.

Once we have established (i) and (ii) it follows that $\lim_{n\to\infty} \frac{1}{\sigma_n} \sum_{j=1}^n \psi_j \circ \mathcal{T}^j \to N(0,1)$ in distribution. Finally, since $[\sum_{j=1}^n \varphi_j \circ \mathcal{T}^j] - [\sum_{j=1}^n \psi_j \circ \mathcal{T}^j] = \mathbf{H}_{n+1} \circ \mathcal{T}^{n+1}$ is bounded in $L^r, r \geq 2$ (Corollary 2.8), $\lim_{n\to\infty} \frac{1}{\sigma_n} \sum_{j=1}^n \varphi_j \circ \mathcal{T}^j \to N(0,1)$ in distribution as well.

For (i), we first prove that

$$\frac{1}{\sigma_n^2} \sum_{j=1}^n \psi_j^2 \circ \mathcal{T}^j \to 1 \quad \text{in probability as } n \to \infty.$$

and then show that in our setting this implies (i) (see Theorem 3.5).

We follow [14, Lemma 3.3 and proof of Theorem 3.1 (II)], which uses an argument of Peligrad [19]. Since $\psi_j = \varphi_j + \mathbf{H}_j - \mathbf{H}_{j+1} \circ T_{n+1}$,

$$\psi_{j}^{2} = \varphi_{j}^{2} + 2\varphi_{j}\mathbf{H}_{j} + \mathbf{H}_{j}^{2} + \mathbf{H}_{j+1}^{2} \circ T_{n+1} - 2\mathbf{H}_{j+1} \circ T_{n+1}(\varphi_{j} + \mathbf{H}_{j})$$

$$= \varphi_{j}^{2} + 2\varphi_{j}\mathbf{H}_{j} + \mathbf{H}_{j}^{2} + \mathbf{H}_{j+1}^{2} \circ T_{n+1} - 2\mathbf{H}_{j+1} \circ T_{n+1}(\psi_{j} + \mathbf{H}_{j+1} \circ T_{n+1})$$

$$= \varphi_{j}^{2} + (\mathbf{H}_{j}^{2} - \mathbf{H}_{j+1}^{2} \circ T_{n+1}) - 2\psi_{j} \cdot \mathbf{H}_{j+1} \circ T_{n+1} + 2\varphi_{j}\mathbf{H}_{j}.$$

Therefore

$$\sum_{j=1}^{n} \psi_{j} \circ \mathcal{T}^{j} = \left(\mathbf{H}_{1}^{2} \circ \mathcal{T}_{1} - \mathbf{H}_{n+1}^{2} \circ \mathcal{T}_{n+1}\right) - \left[\sum_{j=1}^{n} \psi_{j} \circ \mathcal{T}^{j} \cdot \mathbf{H}_{j+1} \circ \mathcal{T}^{j+1}\right] + \left[\sum_{j=1}^{n} \varphi_{j}^{2} \circ \mathcal{T}^{j}\right] + 2\left[\sum_{j=1}^{n} (\varphi_{j} \cdot \mathbf{H}_{j}) \circ \mathcal{T}^{j}\right].$$

By the L^r uniform boundedness of $\mathbf{H}_n \circ \mathcal{T}^n$ (Corollary 2.8), $\frac{1}{\sigma_n^2} \mathbf{H}_{n+1}^2 \circ \mathcal{T}^{n+1} \to 0$ in probability.

Next we show that

$$\frac{1}{\sigma_n^2} \left[\sum_{j=1}^n \psi_j \circ \mathcal{T}^j \cdot \mathbf{H}_{j+1} \circ \mathcal{T}^{j+1} \right] \to 0 \quad \text{in probability.}$$
 (3.1)

Define

$$\mathbf{H}_{j}^{arepsilon}:=\mathbf{H}_{j}\mathbf{1}_{\{|\mathbf{H}_{j}|\leqarepsilon\sigma_{n}\}}$$

By Lemma 2.2,

$$U_n^2 := \int \left(\sum_{j=1}^n [\psi_j \circ \mathcal{T}^j \cdot \mathbf{H}_{j+1}^{\varepsilon} \circ \mathcal{T}^{j+1}] \right)^2 = \int \sum_{j=1}^n [\psi_j \circ \mathcal{T}^j \cdot \mathbf{H}_{j+1}^{\varepsilon} \circ \mathcal{T}^{j+1}]^2.$$

Hence, using Lemma 2.1 for the equal below,

$$U_n^2 \le \varepsilon^2 \sigma_n^2 \sum_{j=1}^n \int \psi_j^2 \circ \mathcal{T}^j$$

$$= \varepsilon^2 \sigma_n^2 \left[\int (\sum_{j=1}^n \varphi_j \circ \mathcal{T}^j)^2 + \int \mathbf{H}_1^2 \circ \mathcal{T}^1 - \int \mathbf{H}_{n+1}^2 \circ \mathcal{T}^{n+1} \right] \le \varepsilon^2 \sigma_n^4. \quad (3.2)$$

For any $a > \varepsilon$ we obtain, using Tchebycheff's inequality in the third and fourth lines below, the inequality (3.2), and our uniform L^r bound on $\mathbf{H}_j \circ \mathcal{T}^j$ (Corollary 2.8), given by the constant \hat{D}

$$m\left(\left|\frac{1}{\sigma_{n}^{2}}\sum_{j=1}^{n}\psi_{j}\circ\mathcal{T}^{j}\cdot\mathbf{H}_{j+1}\circ\mathcal{T}^{j+1}\right|>a\right)$$

$$\leq m\left(\max_{1\leq j\leq n}\left|\mathbf{H}_{j+1}\circ\mathcal{T}^{j+1}\right|>\varepsilon\sigma_{n}\right)+m\left(\left|\frac{1}{\sigma_{n}^{2}}\sum_{j=1}^{n}\psi_{j}\circ\mathcal{T}^{j}\cdot\mathbf{H}_{j+1}^{\varepsilon}\circ\mathcal{T}^{j+1}\right|>a\right)$$

$$\leq \sum_{j=1}^{n}m(|\mathbf{H}_{j+1}\circ\mathcal{T}^{j+1}|>\varepsilon\sigma_{n})+\frac{1}{a^{2}\sigma_{n}^{4}}U_{n}^{2}$$

$$\leq \frac{n}{(\varepsilon\sigma_{n})^{r}}\left(\max_{1\leq j\leq n}\int|\mathbf{H}_{j+1}\circ\mathcal{T}^{j+1}|^{r}\right)+\frac{\varepsilon^{2}}{a^{2}}\leq\frac{n\hat{D}}{(\varepsilon\sigma_{n})^{r}}+\frac{\varepsilon^{2}}{a^{2}}.$$

Take $a = \sqrt{\varepsilon}$; if we use that $\sigma_n^2 \approx n^{\beta}$, then $\beta > \frac{2}{r}$ with $1 \le r < \frac{1}{2\alpha} - \frac{1}{2}$, that is $\beta > \frac{4\alpha}{1-\alpha}$, allows us to obtain (3.1). We defer to the end of this proof the discussion about the possible choices for α, β .

Finally, we show that

$$\frac{1}{\sigma_n^2} \sum_{j=1}^n (\varphi_j^2 + 2\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j \to 1 \quad \text{in probability.}$$
 (3.3)

We know from our Strong Borel-Cantelli Theorem 2.10 that

$$\sum_{j=1}^{n} \varphi_j^2 \circ \mathcal{T}^j = \sum_{j=1}^{n} \mathbb{E}[\varphi_j^2 \circ \mathcal{T}^j] + o(n^{\frac{1}{2} + \varepsilon}) \quad m\text{-a.e.}$$
 (3.4)

We will show in Lemma 3.3 that

$$\frac{1}{\sigma_n^2} \left(\sum_{j=1}^n (\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j - \sum_{j=1}^n \mathbb{E}[(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j] \right) \to 0 \text{ in probability.}$$
 (3.5)

In view of Lemma 2.3, equations (3.3) and (3.5) impliy $\frac{1}{\sigma_n^2} [\sum_{j=1}^n \varphi^2 \circ \mathcal{T}^j + 2 \sum_{j=1}^n (\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j] \to 1$ in probability.

Lemma 3.3 For $< \alpha < 1/8$ and the variance growing as $\sigma_n^2 \approx n^{\beta}, \beta > 2/3$, we have:

$$\frac{1}{\sigma_n^2} \left(\sum_{j=1}^n (\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j - \sum_{j=1}^n \mathbb{E}[(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j] \right) \to 0 \text{ in probability.}$$

Proof Write $S_n = \sum_{j=1}^n (\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j$ and $E_n = \sum_{j=1}^n \mathbb{E}[(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j]$ and estimate

$$\mathbb{E}(|S_n - E_n| > \sigma_n^2 \varepsilon) = \mathbb{E}(|S_n - E_n|^2 > \sigma_n^4 \varepsilon^2)$$

$$\leq \frac{1}{\sigma_n^4 \varepsilon^2} \mathbb{E}(|S_n - E_n|^2).$$

When we estimate $\mathbb{E}(|S_n - E_n|^2)$ we have, as usual, the diagonal terms and a double summation of off-diagonal terms:

$$\mathbb{E}(|S_n - E_n|^2) = \sum_{j=1}^n \mathbb{E}([(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j - m[(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j)]^2)$$

$$+ 2\sum_{j=1}^n \sum_{i=1}^{j-1} \int [(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j - m((\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j)][(\varphi_i \mathbf{H}_i) \circ \mathcal{T}^i - m((\varphi_i \mathbf{H}_i) \circ \mathcal{T}^i)]dx.$$

The sum of diagonal terms is O(n) as $(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j \in L^2(m)$ with uniformly bounded norm if $\alpha < 1/5$.

We note that by prescribing a growth of the variance as $\sigma_n^2 \approx n^{\beta}$, the exponent β must verify $\beta > 1/2$.

We now consider

$$\sum_{j=1}^{n} \sum_{i=1}^{j-1} \int [(\varphi_{j} \mathbf{H}_{j}) \circ \mathcal{T}^{j} - m((\varphi_{j} \mathbf{H}_{j}) \circ \mathcal{T}^{j})] [(\varphi_{i} \mathbf{H}_{i}) \circ \mathcal{T}^{i} - m((\varphi_{i} \mathbf{H}_{i}) \circ \mathcal{T}^{i})] dx$$

$$= \sum_{j=1}^{n} \sum_{i=1}^{j-1} \int [\varphi_{j} \mathbf{H}_{j} - m((\varphi_{j} \mathbf{H}_{j}) \circ \mathcal{T}^{j})] \circ \mathcal{T}^{j} \cdot [\varphi_{i} \mathbf{H}_{i} - m((\varphi_{i} \mathbf{H}_{i}) \circ \mathcal{T}^{i})] \circ \mathcal{T}^{i} dx$$

$$= \sum_{j=1}^{n} \sum_{i=1}^{j-1} \int [\varphi_{j} \mathbf{H}_{j} - m((\varphi_{j} \mathbf{H}_{j}) \circ \mathcal{T}^{j})] \circ \mathcal{T}^{j-i}_{i+1} \cdot [\varphi_{i} \mathbf{H}_{i} - m((\varphi_{i} \mathbf{H}_{i}) \circ \mathcal{T}^{i})] \cdot \mathcal{P}^{i} \mathbf{1} dx$$

$$= \sum_{j=1}^{n} \sum_{i=1}^{j-1} \int [\varphi_{j} \mathbf{H}_{j} - m((\varphi_{j} \mathbf{H}_{j}) \circ \mathcal{T}^{j})] \cdot \mathcal{P}^{j-i}_{i+1} [\mathbf{H}_{i} \varphi_{i} \mathcal{P}^{i} \mathbf{1} - m((\varphi_{i} \mathbf{H}_{i}) \circ \mathcal{T}^{i}) \mathcal{P}^{i} \mathbf{1}] dx.$$

We will prove in Lemma 3.4 below that $||\mathcal{P}_{i+1}^{j-i}[\mathcal{P}^{i}\mathbf{1}\mathbf{H}_{i}\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\mathcal{T}^{i})]||_{2} \leq \frac{C^{*}i}{(j-i)^{\alpha^{*}}}$, where C^{*} is a constant depending only on α and the C^{1} norm of φ (and uniform in i and j). Here the numerator i comes about as $1 \leq i \leq j-1$ and $\alpha^{*} = \frac{1-2\alpha}{2\alpha}$ follows from the decay Theorem 1.2 and Lemma 2.7, provided $\alpha < \frac{1}{2}$. Note also that $||(\varphi_{j}\mathbf{H}_{j}) - m((\varphi_{j}\mathbf{H}_{j}) \circ \mathcal{T}^{j})||_{2}$ is uniformly bounded in j provided $\alpha < \frac{1}{4}$, see Corollary 2.8.

We have to show that each row summation satisfies

$$|\sum_{i=1}^{j-1} \int [(\varphi_j \mathbf{H}_j) - m((\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j)] \mathcal{P}_{i+1}^{j-i} [\mathcal{P}^i \mathbf{1} \mathbf{H}_i \varphi_i - \mathcal{P}^i \mathbf{1} m((\varphi_i \mathbf{H}_i) \circ \mathcal{T}^i)] \ dx| \le j^{\chi}$$

where $n^{1+\chi} = o(\sigma_n^4)$ otherwise the double summation contributes a term which is too large. So we divide the sum into two parts, with $0 < \delta < 1$

$$\sum_{i=j-j^{\delta}}^{j-1} \int [(\varphi_j \mathbf{H}_j) - m((\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j)] \mathcal{P}_{i+1}^{j-i} [\mathcal{P}^i \mathbf{1} \mathbf{H}_i \varphi_i - \mathcal{P}^i \mathbf{1} m((\varphi_i \mathbf{H}_i) \circ \mathcal{T}^i)] dx$$

$$+\sum_{i=1}^{j-j^o}\int [(\varphi_j\mathbf{H}_j)-m((\varphi_j\mathbf{H}_j)\circ\mathcal{T}^j)]\mathcal{P}_{i+1}^{j-i}[\mathcal{P}^i\mathbf{1}\mathbf{H}_i\varphi_i-\mathcal{P}^i\mathbf{1}m((\varphi_i\mathbf{H}_i)\circ\mathcal{T}^i)]\ dx.$$

The first sum we bound by C^*j^{δ} using L^2 bounds without decay. The second uses our decay estimate (see Lemma 3.4) and we get $\sum_{i=1}^{j-j^{\delta}} \frac{C^*i}{(j-i)^{\alpha^*}} \leq C^*j^{1-(\alpha^*-1)\delta} = C^*j^{1+\delta-\alpha^*\delta}$ provided

 $\alpha^* > 1 \ (\iff 0 < \alpha < 1/2).$ Then $|\sum_{i=1}^{j-1} \int [(\varphi_j \mathbf{H}_j) - m((\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j)] \mathcal{P}_{i+1}^{j-i} [\mathcal{P}^i \mathbf{1} \mathbf{H}_i \varphi_i - \mathcal{P}^i \mathbf{1} m((\varphi_i \mathbf{H}_i) \circ \mathcal{T}^i)] \ dx| \le C(j^{\delta} + j^{1+\delta-\alpha^*\delta})$ which is lowest for $\delta = 1/\alpha_*$. We obtain

$$\left|\sum_{j=1}^{n}\sum_{i=1}^{j-1}\int [(\varphi_{j}\mathbf{H}_{j})\circ\mathcal{T}^{j} - m((\varphi_{j}\mathbf{H}_{j})\circ\mathcal{T}^{j})][(\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i} - m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})]dx\right|$$

$$\leq C^{*} n^{1+1/\alpha_{*}} = C^{*} n^{1/(1-2\alpha)}$$

SO

$$E(|S_n - E_n|^2) \le Cn^{1/(1-2\alpha)}.$$

By dividing for σ_n^4 and asking again for a growth like $\sigma_n^2 \approx n^\beta$ we have now that $\beta > \frac{1}{2(1-2\alpha)}$. This estimate allows us to show that $\frac{1}{\sigma_n^2} \left(\sum_{j=1}^n (\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j - \sum_{j=1}^n E[(\varphi_j \mathbf{H}_j) \circ \mathcal{T}^j] \right) \to 0$ in probability.

We now collect the various inequalities involving β , which is the scaling of $\sigma_n^2 \approx n^{\beta}$, and α :

- for our proof of condition (ii) in Brown's Theorem 3.2 we need $\beta > \frac{2}{3}$ and $\alpha < \frac{1}{5}$;
- in Peligrad's argument we have $\beta > \frac{4\alpha}{1-\alpha}$;
- in Lemma 3.3, using that $\alpha < \frac{1}{5}$, we have $\beta > \frac{1}{2}$ and $\beta > \frac{1}{2(1-2\alpha)}$.

These give

$$\alpha < \frac{1}{5}, \quad \beta > \max\left\{\frac{2}{3}, \frac{4\alpha}{1-\alpha}, \frac{1}{2(1-2\alpha)}\right\} \tag{3.6}$$

which are all satisfied if $\alpha < \frac{1}{8}$, $\beta > \frac{2}{3}$, or $\alpha < \frac{1}{5}$, $\beta \ge 1$.

To conclude the proof we need the statement of Lemma 3.4, whose proof is in the Appendix, and of Theorem 3.5, which allows us to get the convergence in probability of the conditional expectations from condition (i) in Brown's Theorem.

Lemma 3.4 *For* $1 \le p < 1/\alpha$

$$\|\mathcal{P}_{k}^{n}\left(\left[\mathcal{P}^{i}\mathbf{1}\mathbf{H}_{i}\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})\right]\right)\|_{p}\leq i\ C_{\alpha,p}\ C_{\varphi}\ n^{-\frac{1}{p\alpha}+1}\left(\log n\right)^{\frac{1}{\alpha}\frac{1-\alpha p}{p-\alpha p}}$$

Theorem 3.5 The following inference holds:

$$\frac{1}{\sigma_n^2} \sum_{j=1}^n \psi_n^2 \circ \mathcal{T}^n \to^p \mathbf{1} \implies \frac{1}{\sigma_n^2} \sum_{j=1}^n \mathbb{E}[\psi_n^2 \circ \mathcal{T}^n | B_{n+1}] \to^p 1.$$

Proof To do this we will use Burckholder's inequality (Theorem 2.10 of [8]).

We will show that

$$\frac{1}{\sigma_n^2} \sum_{i=1}^n (\psi_n^2 \circ \mathcal{T}^n - \mathbb{E}[\psi_n^2 \circ \mathcal{T}^n | B_{n+1}]) \to 0 \quad \text{in probability.}$$

First define $V_n = \psi_n^2 \circ \mathcal{T}^n - \mathbb{E}[\psi_n^2 \circ \mathcal{T}^n | B_{n+1}]$ and note that $\mathbb{E}[V_n | B_{n+1}] = 0$. We define a martingale, reading from left to right,

$$S_1 = V_n, S_2 = V_n + V_{n-1}, V_n + V_{n-1} + V_{n-2} + \dots + V_1$$

with filtration

$$F_0 = B_{n+1}, F_1 = B_n, F_2 = B_{n-1}, ..., F_n = B_0 = B.$$

Then V_n is F_1 measurable as $\psi_n^2 \circ \mathcal{T}^n$ is B_n measurable, since $\mathbb{E}[\psi_n^2 \circ \mathcal{T}^n | B_{n+1}]$ is B_{n+1} measurable and $B_{n+1} \subset B_n$ $\mathbb{E}[\psi_n^2 \circ \mathcal{T}^n | B_{n+1}]$ is F_1 measurable. Similarly V_i is F_{n-i+1} measurable. This implies S_i is F_i measurable.

Note that $\mathbb{E}[V_{n-1}|F_1] = \mathbb{E}[V_{n-1}|B_n] = 0$ so

$$\mathbb{E}[S_{i+1}|F_i] = \mathbb{E}[V_{n-i}|F_i] + S_i = S_i.$$

Hence (S_i, F_i) is a martingale.

By Burckholder's inequality taking p=2 we have

$$\mathbb{E}|S_n|^2 \le C_1 \mathbb{E}(\sum_{i=1}^n V_i^2) \le C_2 \sigma_n^2$$

where C_2 is a universal constant.

Hence
$$P(|S_n| > \sigma_n^2 \varepsilon) = P(|S_n|^2 > \sigma_n^4 \varepsilon^2) \le \frac{C_2}{\varepsilon^2 \sigma_n^2}$$
 by Chebyshev.

4 Central Limit Theorem for nearby maps

Theorem 4.1 Given $\beta \in (0, 1/5)$ and $\varphi \in C^1([0,1])$ if φ is not a coboundary (up to a constant) for T_{β} there exists $\varepsilon > 0$ such that for all parameters $\beta_k \in (\beta - \varepsilon, \beta + \varepsilon)$ the variance grows linearly for any sequential system formed from concatenation of the maps T_{β_k} .

Therefore, by Theorem 3.1 and (3.6), the CLT holds.

Proof

Recall the quantities defined by a concatenation of different maps.

$$\mathbf{H}_{n} = \frac{1}{\mathcal{P}^{n} \mathbf{1}} \left[P_{n}(\varphi_{n-1} \mathcal{P}^{n-1} \mathbf{1}) + P_{n} P_{n-1}(\varphi_{n-2} \mathcal{P}^{n-2} \mathbf{1}) + \dots + P_{n} P_{n-1} \dots P_{1}(\varphi_{0} \mathcal{P}^{0} \mathbf{1}) \right]$$

and

$$\psi_n := \varphi_n + \mathbf{H}_n - \mathbf{H}_{n+1} \circ T_{n+1}.$$

First assume that the maps all coincide with T_{β} so that $P_{\beta}^{n}1 \to h_{\beta}$ (at a polynomial rate in L^{2}), $P_{n}P_{n-1}...P_{n-k} = P_{\beta}^{k}$, where h_{β} is the invariant density for T_{β} and P_{β} is the transfer operator for T_{β} with respect to Lebesgue measure. Furthermore $\varphi_{n} = \varphi - m(\varphi(T_{\beta}^{n})) \to \varphi - \int \varphi h_{\beta} dx$. Denote the \mathbf{H}_{n} corresponding to this situation by $\mathbf{H}_{\beta,n}$.

Note the terms $P_n P_{n-1} ... P_{n-j}(\varphi_{n-j-1} \mathcal{P}^{n-j-1} 1)$ decay at a polynomial rate in L^2 , $\|P_n P_{n-1} ... P_{n-j}(\varphi_{n-j-1} \mathcal{P}^{n-j-1} 1)\|_2 \leq \frac{C}{j^{\tau}}$ for some $\tau > 1$ for $\beta < 1/4$, by Proposition 1.3 and Lemma 2.4. Note that C and τ may be taken as uniform over all T_{β_k} if β_k is close to β .

Combining this with the fact that $P_{\beta}^{n}\mathbf{1} \to h_{\beta}$ in L^{2} (and hence $\frac{1}{P_{\beta}^{n}\mathbf{1}} \to \frac{1}{h_{\beta}}$ in L^{2} as both h_{β} and $P_{\beta}^{n}\mathbf{1}$ are bounded below by $\delta > 0$), we see that given $\varepsilon > 0$ there exists an N such that for all n > N, $\mathbf{H}_{\beta,n} = \frac{1}{h_{\beta}}[P_{\beta}(h_{\beta}\varphi - \int \varphi h_{\beta}dx) + P_{\beta}^{2}(h_{\beta}\varphi - \int \varphi h_{\beta}dx) + \dots + P_{\beta}^{N}(h_{\beta}\varphi - \int \varphi h_{\beta}dx)] + \gamma(\beta,n)$ where $\|\gamma(\beta,n)\|_{2} < \varepsilon$. We define $G_{\beta,N} = \frac{1}{h_{\beta}}[P_{\beta}(h_{\beta}\varphi - \int \varphi h_{\beta}dx) + \dots + P_{\beta}^{N}(h_{\beta}\varphi - \int \varphi h_{\beta}dx)]$ so that $\mathbf{H}_{\beta,n} = G_{\beta,N} + \gamma(\beta,n)$.

Now suppose φ is not a coboundary for T_{β} . Denote by \widetilde{P}_{β} the transfer operator for T_{β} with respect to the invariant measure $d\mu_{\beta} = h_{\beta}dx$. Then $\widetilde{P}_{\beta}^{n}(\varphi) = \frac{1}{h_{\beta}}P_{\beta}^{n}(h_{\beta}\varphi)$ where P_{β} is the transfer operator for T_{β} with respect to Lebesgue measure.

Hence $\frac{1}{h_{\beta}}[P_{\beta}(h_{\beta}\varphi - \int \varphi h_{\beta}dx) + P_{\beta}^{2}(h_{\beta}\varphi - \int \varphi h_{\beta}dx) + ... + P_{\beta}^{N}(h_{\beta}\varphi - \int \varphi h_{\beta}dx)] = \sum_{k=1}^{N} \widetilde{P}_{\beta}^{k}[\varphi - \int \varphi d\mu_{\beta}]$. If φ is not a coboundary then $\sum_{k=1}^{\infty} \widetilde{P}_{\beta}^{k}[\varphi - \int \varphi d\mu_{\beta}]$ converges to a coboundary \widetilde{H}_{β} so that

$$\varphi = \widetilde{\psi}_{\beta} + \widetilde{H}_{\beta} \circ T_{\beta} - \widetilde{H}_{\beta}$$

defines a martingale difference sequence $\{\widetilde{\psi}_{\beta} \circ T_{\beta}^n\}$, where $\widetilde{\psi}_{\beta} \neq 0$ in L^2 (as φ is not a coboundary for T_{β}). Suppose $\|\widetilde{\psi}_{\beta}\|_2 > \eta$.

Choose N large enough that for all n>N, $\|[H_{\beta,n}-H_{\beta,n+1}\circ T_{\beta}]-[\widetilde{H}_{\beta}-\widetilde{H}_{\beta}\circ T_{\beta}]\|_2<\frac{\eta}{20}$ and $\|\widetilde{H}_{\beta}-\sum_{k=1}^{N}\widetilde{P}_{\beta}^{k}[\varphi-\int \varphi d\mu_{\beta}]\|_2<\frac{\eta}{20}$. Then $\|\psi(\beta,n)\|_2>\frac{\eta}{2}$ for all n>N.

Now we consider a concatenation of maps T_{β_k} where β_k is close to β . The idea is to break \mathbf{H}_n into a sum of N terms uniformly close to $G(\beta, N)$ (no matter what the sequence of maps) and a small error.

Choose all β_k 's sufficiently close to β that when we form a concatenation of the maps T_{β_k} we have

$$\|G_{\beta,N} - \frac{1}{\mathcal{P}^{n} \mathbf{1}} \left[P_{n}(\varphi_{n-1} \mathcal{P}^{n-1} \mathbf{1}) + P_{n} P_{n-1}(\varphi_{n-2} \mathcal{P}^{n-2} \mathbf{1}) + \dots + P_{n} P_{n-1} \cdots P_{n-N}(\varphi_{n-N-1} \mathcal{P}^{n-N-1} \mathbf{1}) \right] \|_{2} < \frac{\eta}{20}.$$

We can do this as we have fixed N and the finite terms are continuous in L^2 as $\beta_k \to \beta$, see [16, Theorem 5.1] and Lemmas 2.4, 2.7.

Recall we also have $\|\gamma(\beta, n)\|_2 < \frac{\eta}{20}$ for all $n \ge N$.

Using the uniform contraction (τ and C are uniform for T_{β} where β is in a small neighborhood of β) we have

$$\|\mathbf{H}_{n} - \frac{1}{\mathcal{P}^{n}\mathbf{1}} [P_{n}(\varphi_{n-1}\mathcal{P}^{n-1}\mathbf{1}) + P_{n}P_{n-1}(\varphi_{n-2}\mathcal{P}^{n-2}\mathbf{1}) + \dots + P_{n}P_{n-1}\cdots P_{n-N}(\varphi_{n-N-1}\mathcal{P}^{n-N-1}\mathbf{1})]\|_{2} < \frac{\eta}{20}$$

for all n > N. Then $\|\psi_n\|_2 > \frac{\eta}{10}$ for all n > N and we have linear growth of variance for the concatenation of maps as $\sigma_n^2 = \sum_{k=1}^n E[\psi_n \circ \mathcal{T}^k]^2$.

5 Random compositions of intermittent maps

Suppose $S = \{T_{\alpha_1}, \ldots, T_{\alpha_\ell}\}$ is a finite number of intermittent type maps as in Section 1, with $\alpha_i < \frac{1}{4}$. We will take an iid selection of maps from S according to a probability vector $p = (p_1, \ldots, p_\ell)$ where the probability of choosing map T_{α_i} is p_i . This induces a Bernoulli measure ν on the shift space $\Omega := \{1, \ldots, l\}^{\mathbb{N}}$, where $(i_1, i_2, \ldots, i_n, \ldots)$ corresponds to the sequence of maps: first apply $T_{\alpha_{i_1}}$, then $T_{\alpha_{i_2}}$ and so on. Writing elements of $\omega \in \Omega$ as sequences $\omega := (\omega_0, \omega_1, \ldots, \omega_n, \ldots)$ the shift operator $S : \Omega \to \Omega$, $(S\omega)_i = \omega_{i+1}$ preserves the measure ν .

This random system also induces a Markov process on [0,1] with the transition probability function $P(x,A) = \sum_{i=1}^{\ell} p_{\alpha_i} 1_A(T_{\alpha_i}(x))$. A measure μ is invariant for the Markov process if $P^*\mu = \mu$. In this setting Bahsoun and Bose [4] have shown (among other results) that there is a unique absolutely continuous invariant measure μ and that if $\varphi : [0,1] \to \mathbb{R}$ is a Hölder function then φ satisfies an annealed CLT for this random dynamical system in

the sense that if $\int \varphi d \mu = 0$ then

$$(\nu \times \mu)\{(\omega, x) : \frac{1}{n} \sum_{i=1}^{n} \varphi(T_{(S^{j}\omega)_{0}} \dots T_{(\omega)_{0}} x) \in A\} \to \frac{1}{\sqrt{2\pi\sigma^{2}}} \int_{A} e^{-\frac{x^{2}}{2\sigma^{2}}} dx$$

for some $\sigma^2 \ge 0$. In fact the result of Bahsoun and Bose [4] also shows that this convergence is with respect to $(\nu \times m)$ where m is Lebesgue measure on [0, 1].

This follows from a well known result by Eagleson [9] which states the equivalence of the convergence in distribution for measures which are absolutely continuous one with respect to the other.

We will show that almost every realization of choices of concatenations of maps, i.e. with respect to the product measure ν , satisfies a self-norming CLT if:

(*) φ is not a coboundary for all maps i.e. there exists an i such that $\varphi \neq \psi \circ T_{\alpha_i} - \psi$ for any measurable (hence Hölder by standard Livsic theory) function ψ .

First we show that if we take a random composition of a finite number of intermittent type maps we obtain linear growth of the variance almost surely under assumption (*).

Lemma 5.1 If φ is not a coboundary for all maps, i.e. there exists an i such that $\varphi \neq \psi \circ T_{\alpha_i} - \psi$ for any measurable ψ , then for ν -almost every sequence of maps T_i

$$\sigma_n^2 := \int \left(\sum_{j=1}^n \varphi \circ \mathcal{T}^i - m(\varphi \circ \mathcal{T}^i) \right)^2 dx$$

grows at a linear rate in that $\sigma_n^2 \geq Cn$ for sufficiently large n for some C > 0.

Proof

Under our assumption φ is not a coboundary for one of the maps, say T_{α_1} .

We will construct a martingale decomposition using the transfer operator Q_{α_1} corresponding to the invariant measure μ_{α_1} for T_{α_1} . The invariant measure μ_{α_1} has a density h_{α_1} .

The coboundary function is defined by $H_{\alpha_1} = \sum_{j=1}^{\infty} Q_{\alpha_1}^j [\varphi - \int \varphi d\mu_{\alpha_1}]$ where Q_{α_1} is the adjoint operator of the Koopman operator $U\varphi = \varphi \circ T_{\alpha_1}$ with respect to the invariant measure $d\mu_{\alpha_1} = h_{\alpha_1} dx$ for T_{α_1} .

When we do the usual decomposition $\varphi = \psi_{\alpha_1} + H_{\alpha_1} - H_{\alpha_1} \circ T_{\alpha_1}$ then the martingale difference function ψ_{α_1} is bounded below from zero in L^2 . Suppose $\|\psi_{\alpha_1}\|_2 > \rho > 0$.

It is known that $Q_{\alpha_1}^n(\varphi) = \frac{1}{h_{\alpha_1}} P_{\alpha_1}^n(h_{\alpha_1}\varphi)$ where P_{α_1} is the adjoint of the Koopman operator of T_{α_1} with respect to Lebesgue measure. Furthermore $P_{\alpha_1}^n 1 \to h_{\alpha_1}$ (at a polynomial rate in L^2) and since $\Pi - j\mathbf{1}$ lies in the cone C_2 and $\int \Pi_j \mathbf{1} 1 dx = 1$,

$$P_{\alpha_1}^k[h_{\alpha_1}-\Pi_i 1]\to 0$$

in L^2 at a uniform polynomial rate, in fact $||P_{\alpha_1}^k[h_{\alpha_1} - \Pi_j 1]||_2 \le C \frac{1}{k^{1+\eta}}$ where C and η are uniform over $\Pi_j 1$.

Now we consider the quantities defined by a concatenation of different maps. We will use the notation from previous sections.

$$\mathbf{H}_{n} = \frac{1}{\mathcal{P}^{n} \mathbf{1}} \left[P_{n}(\varphi_{n-1} \mathcal{P}^{n-1} \mathbf{1}) + P_{n} P_{n-1}(\varphi_{n-2} \mathcal{P}^{n-2} \mathbf{1}) + \dots + P_{n} P_{n-1} \dots P_{1}(\varphi_{0} \mathcal{P}^{0} \mathbf{1}) \right]$$

and $\psi_n := \varphi_n + \mathbf{H}_n - \mathbf{H}_{n+1} \circ T_{n+1}$.

We will first consider what happens when we have a sequence of k maps T_{α_1} applied one after the other. We will suppose we have concatenated n maps and then apply k T_{α_1} maps in turn.

Then $\varphi_{n+k} = \varphi - \int \varphi(T_{\alpha_1}^k T_n ... T_1) dx = \varphi - \int \varphi P_{\alpha_1}^k \Pi_n 1 dx = \varphi - \int \varphi h_{\alpha_1} dx + \int \varphi P_{\alpha_1}^k [h_{\alpha_1} - \Pi_n 1] dx$ where $\|P_{\alpha_1}^k [h_{\alpha_1} - \Pi_n 1]\|_2 \leq \frac{C}{k^{1+\delta}}$.

We are considering here n fixed and k increasing.

Note the terms $P_nP_{n-1}...P_{n-j}(\varphi_{n-j-1}\mathcal{P}^{n-j-1}1)$ decay at a polynomial rate in L^2 , in fact $\|P_nP_{n-1}...P_{n-j}(\varphi_{n-j-1}\mathcal{P}^{n-j-1}1)\|_2 \leq \frac{C}{j^{1+\eta}}$. Note that C and η may be taken as uniform over all choices of T_{α_i} in the concatenation.

Combining this with the fact that $\mathcal{P}_{\alpha_1}^k\Pi_n 1 \to h_{\alpha_1}$ in L^2 (and hence $\frac{1}{\mathcal{P}_{\alpha_1}^k\Pi_n 1} \to \frac{1}{h_{\alpha_1}}$ in L^2 as both h_{α_1} and $\mathcal{P}_{\alpha_1}^k\Pi_n 1$ are bounded below by $\delta > 0$), we see that given $\rho > 0$ there exists an r such that for all m > n + rk, $\mathbf{H}_m = \frac{1}{h_{\alpha_1}}[P_{\alpha_1}(h_{\alpha_1}\varphi - \int \varphi h_{\alpha_1}dx) + P_{\alpha_1}^2(h_{\alpha_1}\varphi - \int \varphi h_{\alpha_1}dx)] + \gamma(m,\alpha_1)$ where $\|\gamma(m,\alpha_1)\|_2 < \frac{\rho}{20}$.

Now $\frac{1}{h_{\alpha_1}}[P(h_{\alpha_1}\varphi - \int \varphi h_{\alpha_1}dx) + P^2(h_{\alpha_1}\varphi - \int \varphi h_{\alpha_1}dx) + ... + P^k(h_{\alpha_1}\varphi - \int \varphi h_{\alpha_1}dx)] = \sum_{j=1}^k Q_{\alpha_1}^j[\varphi - \int \varphi d\mu_{\alpha_1}].$ The infinite sum $\sum_{j=1}^\infty Q_{\alpha_1}^j[\varphi - \int \varphi d\mu_{\alpha_1}]$ converges to H_{α_1} at a polynomial rate.

We choose k large enough that $||H_{\alpha_1} - \sum_{j=1}^i Q_{\alpha_1}^j [\varphi - \int \varphi d\mu_{\alpha_1}]||_2 \leq \frac{\rho}{20}$. Recall

$$\varphi = \psi_{\alpha_1} + H_{\alpha_1} \circ T_{\alpha_1} - H_{\alpha_1}$$

defines a martingale difference sequence $\{\psi_{\alpha_1} \circ T_{\alpha_1}^j\}$, where ψ_{α_1} is bounded away from zero in L^2 (as φ is not a coboundary for T_{α_1}). We assumed $\|\psi_{\alpha_1}\|_2 > \rho$.

We have shown that if we choose k and r large enough then $\|\mathbf{H}_m - H_{\alpha_1}\|_2 < \frac{\rho}{10}$ for all m > n + rk and hence as

$$\psi_m := \varphi_m + \mathbf{H}_m - \mathbf{H}_{m+1} \circ T_{m+1}$$

we see that $\|\psi_m - H_{\alpha_1}\|_2 \le \frac{\rho}{5}$ and hence $\|\psi_m\|_2 > \frac{\rho}{2}$.

This implies linear growth in the random composition setting as almost all choices of maps will have rk long sequences of the map T_{α_1} at a fixed frequency. In fact the only way we won't obtain linear growth almost surely is if the function φ is a coboundary for all the maps T_{α_i} .

The next theorem is an immediate consequence of the previous lemma and Theorem 3.1 (see (3.6) for the bound on α).

Theorem 5.2 If $\alpha_i < 1/5$ for all $1 \le i \le \ell$ and φ is not a coboundary for all maps then $\sigma_n^2 \ge Cn$ for some C > 0 and hence φ satisfies a CLT for ν almost every sequence of maps.

6 Appendices

6.1 Gál-Koksma Theorem.

We recall the following result of Gál and Koksma as formulated by W. Schmidt [20, 21] and stated by Sprindzuk [22]:

Theorem 6.1 Let $(\Omega, \mathcal{B}, \mu)$ be a probability space and let $f_k(\omega)$, (k = 1, 2, ...) be a sequence of non-negative μ measurable functions and g_k , h_k be sequences of real numbers such that $0 \le g_k \le h_k \le 1$, (k = 1, 2, ...,). Suppose there exists C > 0 such that

$$\int \left(\sum_{m < k \le n} (f_k(\omega) - g_k) \right)^2 d\mu \le C \sum_{m < k \le n} h_k$$

for arbitrary integers m < n. Then for any $\varepsilon > 0$

$$\sum_{1 \le k \le n} f_k(\omega) = \sum_{1 \le k \le n} g_k + O(\Theta^{1/2}(n) \log^{3/2 + \varepsilon} \Theta(n))$$

for μ -a.e. $\omega \in \Omega$, where $\Theta(n) = \sum_{1 \le k \le n} h_k$.

6.2 Proof of Lemma 3.4

Proof For simplicity of notation we discuss only the case k = 1; the general case is the same, since we use the n Perron-Frobenius maps in \mathcal{P}_k^n only for the decay given by Theorem 1.2.

The idea is to write $[\mathcal{P}^{i}\mathbf{1}\mathbf{H}_{i}\varphi_{i} - \mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})]$ as a difference of 2i functions in the cone of the same integral. By writing explicitly \mathbf{H}_{i} we get

$$[\mathcal{P}^{i}\mathbf{1}\mathbf{H}_{i}\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})]=\left[\sum_{k=1}^{i}\prod_{j=0}^{k-1}P_{i-j}(\varphi_{i-k}\mathcal{P}^{i-k}\mathbf{1})\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})\right]=\left[\sum_{k=1}^{i}\prod_{j=0}^{k-1}P_{i-j}(\varphi_{i-k}\mathcal{P}^{i-k}\mathbf{1})\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})\right]=\left[\sum_{k=1}^{i}\prod_{j=0}^{k-1}P_{i-j}(\varphi_{i-k}\mathcal{P}^{i-k}\mathbf{1})\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})\right]$$

$$\left[\sum_{k=1}^{i}\prod_{j=0}^{k-1}P_{i-j}(\varphi_{i-k}\mathcal{P}^{i-k}\mathbf{1})\varphi_{i}-\mathcal{P}^{i}\mathbf{1}\sum_{k=1}^{i}m((\varphi_{i}\frac{1}{\mathcal{P}^{i}\mathbf{1}}\prod_{j=0}^{k-1}P_{i-j}(\varphi_{i-k}\mathcal{P}^{i-1}\mathbf{1})\circ\mathcal{T}^{i})\right]=$$

$$\sum_{k=1}^{i} \left[\varphi_{i} \mathcal{P}_{i-k+1}^{k} (\varphi_{i-k} \mathcal{P}^{i-k} \mathbf{1}) - \mathcal{P}^{i} \mathbf{1} m ((\varphi_{i} \frac{1}{\mathcal{P}^{i} \mathbf{1}} \mathcal{P}_{i-k+1}^{k} (\varphi_{i-k} \mathcal{P}^{i-1} \mathbf{1}) \circ \mathcal{T}^{i}) \right].$$

Call $C_{k,i} := m((\varphi_i \frac{1}{\mathcal{P}_{i-k+1}} \mathcal{P}_{i-k+1}^k (\varphi_{i-k} \mathcal{P}^{i-1} 1) \circ \mathcal{T}^i);$ then consider the quantity

$$(*) := \varphi_i \mathcal{P}_{i-k+1}^k (\varphi_{i-k} \mathcal{P}^{i-k} \mathbf{1}) - \mathcal{P}^i \mathbf{1} C_{k,i}.$$

Since $\varphi_{i-k} \in C^1$ and $\mathcal{P}^{i-k}\mathbf{1} \in \mathcal{C}_2$ we can write by Lemma 2.4

$$\varphi_{i-k}\mathcal{P}^{i-k}\mathbf{1} = F_{i-k} - G_{i-k}$$

with $F_{i-k}, G_{i-k} \in \mathcal{C}_2$. By the invariance of the cone, the functions $h_{i-k}^{(1)} := \mathcal{P}_{i-k+1}^k F_{i-k}$; $h_{i-k}^{(2)} := \mathcal{P}_{i-k+1}^k G_{i-k}$ are still in the cone, and we rewrite (*) as

$$(*) = \varphi_i h_{i-k}^{(1)} - \varphi_i h_{i-k}^{(2)} - C_{i,k} \mathcal{P}^i \mathbf{1}.$$

Although the functions (in the cone), F_{i-k} , G_{i-k} are not of zero mean, we can still apply Lemma 2.4 and split the product of φ_i with them into the differences of two new functions belonging to the cone, namely

$$\varphi_i h_{i-k}^{(1)} = M_{i-k}^{(1)} - M_{i-k}^{(2)}; \ \varphi_i h_{i-k}^{(2)} = N_{i-k}^{(1)} - N_{i-k}^{(2)}$$

with $M_{i-k}^{(1,2)}, N_{i-k}^{(1,2)} \in \mathcal{C}_2$. We finally have

$$(*) = [M_{i-k}^{(1)} + N_{i-k}^{(2)}] - [M_{i-k}^{(2)} + N_{i-k}^{(1)} + C_{i,k}\mathcal{P}^{i}\mathbf{1}] := R_{i,k} - S_{i,k}$$

where the functions $R_{i,k}$, $S_{i,k}$ are in the cone and have the same expectation. Before continuing, let us summarize what we got

$$[\mathcal{P}^{i}\mathbf{1}\mathbf{H}_{i}\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})]=\sum_{k=1}^{i}(R_{i,k}-S_{i,k}).$$

By taking the power \mathcal{P}^n on both sides we have by our Theorem 1.2 on the loss of memory and Proposition 1.3

$$\|\mathcal{P}^{n}\left(\left[\mathcal{P}^{i}\mathbf{1}\mathbf{H}_{i}\varphi_{i}-\mathcal{P}^{i}\mathbf{1}m((\varphi_{i}\mathbf{H}_{i})\circ\mathcal{T}^{i})\right]\right)\|_{p}\leq \sum_{k=1}^{i}C_{\alpha,p}(\|R_{i,k}\|_{1}+\|S_{i,k}\|_{1})n^{-\frac{1}{p\alpha}+1}\left(\log n\right)^{\frac{1}{\alpha}\frac{1-\alpha p}{p-\alpha p}}.$$

From Lemma 2.4, one observes that if we have $\varphi \in C^1([0,1])$ and $H \in \mathcal{C}_2$ the splitting $\varphi H = A - B$, with $A, B \in \mathcal{C}_2$ is such that the functions A, B depend only on the C^1 norm of φ and the integrals $m(H), m(\varphi H)$. In our case since $\varphi_i(x) = \varphi(x) - m(\varphi \circ \mathcal{T}^i)$, we have that $\|\varphi_i\|_{C^1} \leq \|\varphi\|_{C^1}$; moreover, at each application of Lemma 2.4, the function H is either $\mathcal{P}^i \mathbf{1}$ or obtained by applying \mathcal{P}^ℓ to a function obtained in the previous step and which only depends upon $\|\varphi\|_{C^1}$; in conclusion the norms $\|R_{i,k}\|_1, \|S_{i,k}\|_1$ are bounded by a function C_{φ} which only depends on the choice of the observable φ . We finally get

$$\|\mathcal{P}^n\left(\mathcal{P}^i\mathbf{1}[\mathbf{H}_i\varphi_i - m((\varphi_i\mathbf{H}_i)\circ\mathcal{T}^i)]\right)\|_p \le i \ C_{\alpha,p} \ C_{\varphi} \ n^{-\frac{1}{p\alpha}+1} \left(\log n\right)^{\frac{1}{\alpha}\frac{1-\alpha p}{p-\alpha p}}.$$

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