



Probabilistic limit theorems via the operator perturbation method, under optimal moment assumptions

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PROBABILISTIC LIMIT THEOREMS VIA THE OPERATOR PERTURBATION METHOD, UNDER OPTIMAL MOMENT ASSUMPTIONS

FRANÇOISE PÈNE

CONTENTS

1. Introduction	2
2. Overview of the key ideas of the proofs	5
3. Quasi-compactness and perturbation	6
4. Taylor expansions for the resolvent and eigenprojectors	9
5. Expansions of Fourier eigenprojectors and eigenvalues in Markovian or dynamical contexts	12
5.1. General context and toy model	12
5.2. Square integrable observables	14
5.3. Non square integrable observable	19
6. Probabilistic limit theorems	22
6.1. Central and local limit theorems	22
6.2. Edgeworth expansions ($d = 1$)	26
7. Limit theorems for Markov random walks	28
7.1. General results	29
7.2. ρ -mixing Markov chains	30
7.3. V -geometrically ergodic Markov chains	33
7.4. Lipschitz iterative models	36
Appendix A. Markov Additive processes	40

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ABSTRACT. The Nagaev-Guivarc'h operator perturbation method is well known to provide various probabilistic limit theorems for Markov random walks. A natural conjecture is that this method should provide these limit theorems under the same moment assumptions as the optimal ones in the case of sums of independent and identically distributed random variables. In the past decades, assumptions have been weakened, without achieving fully this purpose (achieving it either with the help of an extra proof of the central limit theorem, or with an additional ε in the moment assumptions). The aim of this article is to give a positive answer to this conjecture via the Keller-Liverani theorem. We present here an approach allowing the establishment of limit theorems (including higher order ones) under optimal moment assumptions. Our method is based on Taylor expansions obtained via the perturbation operator method, combined with a new weak compactness argument without the use of any other extra tool (such as Martingale decomposition method, etc.).

1. INTRODUCTION

Let $(X_n)_{n \geq 0}$ be a Markov chain with values in Ω , with transition operator P and with stationary measure μ and $f : \Omega \times \Omega \times E \rightarrow \mathbb{R}$ be a measurable function. Let ν be the distribution of X_0 (i.e. the initial distribution of the Markov chain). We set \mathcal{P}_ν for the Markov distribution with transition operator P and initial probability measure ν . We are interested in the study of the Markov random walk $(S_n)_{n \geq 1}$ given by¹

$$S_n := \sum_{k=1}^n Y_k \quad \text{with } Y_k := f(X_{k-1}X_k, Z_k),$$

where Z_i are independent and identically distributed (i.i.d.) random variables independent of $(X_k)_k$ and with common distribution \mathbf{P} . We assume moreover throughout this article that Y_1 is centered with respect to $\mathcal{P}_\mu \otimes \mathbf{P}$. Our goal is to establish probabilistic limit theorems for $(S_n)_{n \geq 1}$ under moment assumptions known to be optimal in the case of sums of independent and identically distributed (i.i.d.) random variables. Recall that if $(Y_k)_{k \geq 1}$ were a sequence of centered i.i.d. random variables:

- If $Y_1 \in \mathbb{L}^2(\mu)$, then the usual central limit theorem (CLT) holds true : $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable \mathcal{W} with variance $\mathbb{E}[Y_1^2]$, with density $h_{\mathcal{W}}$.
- If $Y_1 \in \mathbb{L}^2(\mu)$ is \mathbb{Z} -valued and satisfies some non-lattice condition, then the usual local limit theorem (LLT) holds true: $\mathbb{P}(S_n = k) \sim h_{\mathcal{W}}(k/\sqrt{n})n^{-\frac{1}{2}}$, uniformly in $k \in \mathbb{Z}$.
- If $Y_1 \in \mathbb{L}^3(\mu)$ and satisfies some non-lattice condition, then there is a first order Edgeworth expansion: $\mathbb{P}\left(\frac{S_n}{\sqrt{n}} \leq x\right) = \mathbb{P}(\mathcal{W} \leq x) + \frac{\mathfrak{B}_1(x)}{n^{\frac{1}{2}}} h_{\mathcal{W}}(x) + o(n^{-\frac{1}{2}})$, uniformly in $x \in \mathbb{R}$, where \mathfrak{B}_1 is a polynomial function.
- If $Y_1 \in \mathbb{L}^{r+2}(\mu)$ is \mathbb{Z} -valued and satisfies some non-lattice condition, then there is an expansion of order r in the LLT: $\mathbb{P}(S_n = k) = h_{\mathcal{W}}(0)n^{-\frac{1}{2}} + \sum_{j=1}^{\lfloor r/2 \rfloor} \frac{A_j(k)}{n^{\frac{1}{2}+j}} + o(n^{-\frac{1+r}{2}})$, where A_j is a polynomial function.
- If $Y_1 \in \mathbb{L}^{r+2}(\mu)$ satisfies some non-lattice condition as well as some diophantine condition of the form $\mathbb{E}[e^{isY_1}] < e^{-\widehat{C}|s|^{-\alpha}}$ for all s large enough (see [8]), then there is an Edgeworth expansion of order r : $\mathbb{P}\left(\frac{S_n}{\sqrt{n}} \leq x\right) = \mathbb{P}(\mathcal{W} \leq x) + h_{\mathcal{W}}(x) \sum_{j=1}^r \frac{\mathfrak{B}_j(x)}{n^{\frac{j}{2}}} + o(n^{-\frac{r}{2}})$, uniformly in $x \in \mathbb{R}$, for some polynomial functions \mathfrak{B}_j .

¹We explain in appendix that this notion includes the discrete-time Markov additive processes considered in [24].

We will establish such results for Markov random walks. We will also investigate other results such as

- convergence to a stable distribution in the multi-dimensional setting,
- local limit theorem for observables with values in \mathbb{Z}^d , including the case of convergence to stable distributions,
- expansions in the LLT for non \mathbb{Z} -valued random variables.

We will state general results in the context of geometrically ergodic Markov chains and will illustrate all of them on a toy model of Knudsen gas. Let us recall that the Markov chain $(X_n)_{n \geq 0}$ (or equivalently its transition operator P) is said to be geometrically ergodic on some complex Banach space of functions \mathcal{B}_1 if its transition operator P satisfies

$$\exists \vartheta \in]0, 1[, \quad \|P^n - \mathbb{E}_\mu[\cdot]\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta^n). \quad (1)$$

As a consequence of our general results, we will prove limit theorems under optimal moment assumptions. We illustrate our results on classical families of Markov random walks (for ρ -mixing, V -geometrically ergodic Markov chains or Lipschitz iterative Markov chains) and obtain in particular the following result.

Theorem 1.1. *Let $m \geq 2$ and $\kappa > 0$. Assume one of the following conditions holds true:*

- *either P is ρ -mixing, $\nu = \mu$ and $Y_1 \in \mathbb{L}^m(\mathcal{P}_\mu \otimes \mathbf{P})$ centered;*
- *or there exists $\vartheta \in]0, 1[, C > 0$ and an unbounded continuous function $V : \Omega \rightarrow [1, +\infty[$ such that $\mathbb{E}_\mu[V] + \mathbb{E}_\nu[V^{\frac{\kappa}{m}}] < \infty$ and $\left\| \frac{P^n(\cdot) - \mathbb{E}_\mu[\cdot]}{V} \right\|_\infty \leq C\vartheta^n \|\cdot\|/V\|_\infty$. Assume $\sup_{(x,y) \in \Omega^2} \mathbb{E}[|f(x,y,Z_1)|^m]/(V(x) + V(y)) < \infty$, Y_1 is centered with respect to $\mathcal{P}_\mu \otimes \mathbf{P}$.*
- *or (Ω, d) is a non-compact metric space, $P(g) = \mathbb{E}[g(F(x, \theta))]$ with θ a random variable and² with $F(\cdot, \theta) : \Omega \rightarrow \Omega$ strictly contracting, $f : \Omega \rightarrow \mathbb{R}$ is Lipschitz continuous (we consider here $Y_n = f(X_n)$) and $\mathbb{E}[d(x_0, F(x_0, \theta))^{(r+1)}] + \mathbb{E}_\nu[d(x_0, \cdot)^\kappa] < \infty$ (for some fix $x_0 \in \Omega$).*

Then $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable with variance matrix $\sum_{n \in \mathbb{Z}} \text{Cov}_{\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}}(Y_1, Y_{|n|+1})$.³ Assume moreover that f is non-lattice (either in \mathbb{Z}^d if we precise that f takes its values in \mathbb{Z}^d or in \mathbb{R}^d otherwise). Then

- If f is \mathbb{Z}^d -valued, then $(S_n)_{n \geq 1}$ satisfies the local limit theorem (LLT).*
- if $d = 1$, $\kappa \geq m - 1$, and $m \geq 3$, then $(S_n)_{n \geq 1}$ satisfies a first order Edgeworth expansion.*
- if $d = 1$, $\kappa \geq m - 1$, $m \geq r + 2$ and if f is \mathbb{Z} -valued, then $(S_n)_{n \geq 1}$ satisfies a LLT with expansion of order r .*
- if $d = 1$, $\kappa \geq m - 1$ and $m \geq r + 2$, and if some diophantine condition is satisfied, then there is an Edgeworth expansion of order r and also an expansion of order r in the LLT.*

A simple case in which the diophantine condition holds true is if $\mathbb{E}[e^{isY_1}|X_0, X_1] < e^{-\widehat{C}|s|^{-\alpha}}$ for all s large enough and for some $\alpha > 0$ (with the additional assumption that $r < \alpha^{-1} + \frac{1}{2}$ for r -order Edgeworth expansion). Let us indicate that other examples in compact situations are given in [9].

This result will appear as an application of general results for Markov random walks, that are consequences of Taylor expansions for eigenprojectors and for the dominating eigenvalue of

²Actually, we state a much more general result (Theorem 7.9) under weaker assumptions on F and f .

³In this manuscript, given two d -dimensional square integrable random variables $A = (A_1, \dots, A_d)$ and $B = (B_1, \dots, B_d)$, we write $\text{Cov}(A, B)$ for the symmetric matrix $\left(\frac{\text{Cov}(A_i, B_j) + \text{Cov}(A_j, B_i)}{2} \right)_{i,j=1, \dots, d}$.

the operators P_t obtained from the transition operator P by Fourier perturbation (P and P_t acting on some complex Banach space \mathcal{B}_1). The use of operator perturbation techniques to prove probabilistic limit theorem is usually called the Nagaev-Guivarc'h method in reference to the seminal works by these two mathematicians [35, 36, 16] (see also [30] and [19]). This method was first implemented in the case of nice bounded observables so that $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_1)$ is smooth (continuous, \mathcal{C}^k , analytic) implying the smoothness of the eigenprojectors. The Keller-Liverani theorem [31] strengthen this approach making possible the study of the case of unbounded observables for which $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_1)$ is not continuous (see also [11] for a presentation of this method in french). The idea consists in considering two Banach spaces $\mathcal{B}_1 \subset \mathcal{B}_2$ (with continuous inclusion) and then in using the continuity of $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ to prove the continuity of the eigenprojectors as elements of $\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$.

In [22, 23], Hervé proved a local limit theorem and a Berry-Esséen estimate (i.e. an edgeworth expansion with error in $\mathcal{O}(n^{-\frac{1}{2}})$) for geometrically ergodic Markov chains (including V -geometrically ergodic) under the optimal assumption by proceeding in two steps: he first establishes the central limit theorem (CLT) using another method (using a martingale approximation, *à la* Gordin [13]) and then deduces from this result an expansion for the dominating eigenvalue, the continuity of the eigenprojectors being ensured by the Keller and Liverani theorem. This argument was reused in [26]. This method relies on the fact that we both have the continuity of the eigenprojectors and a proof of the CLT by another argument. Note also that it can only work for limit theorems using only the first non null derivative of the dominating eigenvalue λ_t of P_t , that is the dominating term of $\lambda_t - 1$ as t goes to 0 (remind that the geometrically ergodicity implies that $\lambda_0 = 1$ is the single dominating eigenvalue of P and that it is simple with only constant eigenvectors). As soon as we need more derivatives, we have to find another way.

The idea that continuity of the eigenprojectors and the first order term of $\lambda_t - 1$ as t goes to 0 are enough to prove convergence in distribution as well as local limit theorems has been implemented to prove convergence to stable distribution or gaussian distribution with non standard normalization in [1, 39] in the context of dynamical systems (chaotic billiards) and in [17] in the context affine random walks.

In [26], motivated by the establishment of further probabilistic limit theorems in Markovian context under the weakest possible moment assumptions, with Hervé, we extended the continuity statement of Keller and Liverani in a \mathcal{C}^r -smoothness result. This approach enabled us to prove some limit theorems under suboptimal assumptions, with an additional ε in the moment-type assumptions. In particular, we proved the first order Edgeworth expansion under the suboptimal moment assumption $m > 3$. The general \mathcal{C}^r -perturbation theorem of [26] was later used again in [24] and [25] in the context of Markov random walks and M -estimators, respectively.

In the present paper, using carefully a new weak compactity argument, we obtain the limit theorems under the optimal moment assumptions without requiring an extra probabilistic argument (such as martingale approximation). In Section 2 we present the key ideas with a focus on our new argument. We state in Sections 3 and 4 general results on quasi-compactness and Taylor expansions (in t) of the resolvent and so of eigenprojectors for general families $(P_t)_t$ of continuous linear operators. The last sections are devoted to the general context of Markov random walks. In Section 5, we establish Taylor expansions for the dominating eigenvalue. Probabilistic limit theorems are then inferred in Section 6 (CLT, LLT, higher order expansions in the Berry-Esséen theorem as well as in the local limit theorem) with the use of the general results of [10]. In Section 7, we state probabilistic limit theorems in a general context of Markov chains and apply it to examples, proving in particular Theorem 1.1. Actually, Section 7 contains more general results (especially for Lipschitz iterative models). We end this article with an appendix, in which

we explain how the Markov additive processes studied in e.g. [24] are in the scope of the present work.

2. OVERVIEW OF THE KEY IDEAS OF THE PROOFS

The key idea of the Nagaev-Guivarc'h method [35, 36, 16] consists in

- noticing that

$$\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}} [e^{itS_n}] = \mathbb{E}_\mu [P_t^n(\mathbf{1})],$$

and more generally that

$$\mathbb{E}_{\mathcal{P}_\nu \otimes \mathbf{P}^{\otimes \mathbb{N}}} [g(X_0) e^{itS_n} h(X_n)] = \mathbb{E}_\nu [g P_t^n(h)],$$

with

$$P_t(h)(x) = \int_E P(e^{itf(x, \cdot, \omega)} h(\cdot))(x) d\mathbf{P}(\omega) = \mathbb{E} [e^{itY_1} h(X_1) | X_0 = x],$$

- using the fact that the geometric ergodicity of the Markov chain (see (1)), i.e. the quasi-compactness with single simple dominating eigenvalue 1 of P on some Banach space \mathcal{B}_1 :

$$\exists \vartheta \in]0, 1[, \quad \|P^n - \mathbb{E}_\mu[\cdot]\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta^n)$$

will imply the one of P_t , for small $|t|$ with a uniform bound:

$$\exists \vartheta_1 \in]0, 1[, \quad \sup_{|t| < b} \|P_t^n - \lambda_t^n \Pi_t\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta_1^n), \quad (2)$$

with $\lambda_t \in \mathbb{C}$ the dominating eigenvalue of P_t and with $\Pi_t \in \mathcal{L}(\mathcal{B}_1)$ the corresponding eigenprojector,

- in proving the smoothness of $t \mapsto \lambda_t$ and of $t \mapsto \Pi_t$,
- inferring the probabilistic limit theorems using characteristic functions as in the case of sums of i.i.d. random variables.

In this whole paper, we will work with different Banach spaces satisfying some continuous embedding property that we introduce now. For two Banach spaces $(\mathcal{B}_j, \|\cdot\|_{(j)})$, $j \in \{1, 2\}$, the notation $\mathcal{B}_1 \hookrightarrow \mathcal{B}_2$ will mean that $\mathcal{B}_1 \subset \mathcal{B}_2$ and $\|\cdot\|_{(2)} \leq \|\cdot\|_{(1)}$.

In [31], Keller and Liverani proved that when $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_1)$ is not continuous but only $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ is continuous with $\mathcal{B}_1 \hookrightarrow \mathcal{B}_2$, it may still be possible to implement this method to get (2) and the continuity of $t \mapsto \Pi_t \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$.

This idea has been extended in [26] to prove the \mathcal{C}^r -smoothness of $t \mapsto \Pi_t \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ and of $t \mapsto \lambda_t \in \mathbb{C}$. To this end, assuming $Y_1 \in \mathbb{L}^{r+1}$ and exploiting actually only the fact that $Y_1 \in \mathbb{L}^{r+\varepsilon}$, we worked with a double chain of Banach spaces:

$$\mathcal{B}_0 \hookrightarrow \tilde{\mathcal{B}}_0 \hookrightarrow \mathcal{B}_1 \hookrightarrow \tilde{\mathcal{B}}_1 \hookrightarrow \dots \hookrightarrow \mathcal{B}_r \hookrightarrow \tilde{\mathcal{B}}_r \hookrightarrow \mathbb{L}^1$$

such that $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_j, \tilde{\mathcal{B}}_j)$ is continuous and such that $t \mapsto P_t \in \mathcal{L}(\tilde{\mathcal{B}}_j, \mathcal{B}_{j+m})$ is \mathcal{C}^m with P_t acting quasi-compactly on each \mathcal{B}_j and $\tilde{\mathcal{B}}_j$. It is proved in [26] that the resolvent and so the eigenprojectors are \mathcal{C}^r as elements of $\mathcal{L}(\mathcal{B}_j, \tilde{\mathcal{B}}_{j+m})$ and that $t \mapsto \lambda_t$ is also \mathcal{C}^r .

This way of proceeding allowed us to prove limit theorems under suboptimal hypotheses (typically $Y_1 \in \mathbb{L}^{r+\varepsilon}$ when the optimal condition in the i.i.d. case was $Y_1 \in \mathbb{L}^r$). This was already a great improvement, but was not totally satisfactory because of the additional ε in the moment assumptions.

We present here an approach that allows to obtain the optimal moment assumptions. Before entering deaplier in the presentation of the operator method used here, we explain the key ideas

making this adaptation possible (in particular a new key weak compactness argument in \mathbb{L}^p combined with tailored adaptations of the chain of Banach spaces). Assuming $Y_1 \in \mathbb{L}^{r+1}$, we work with a single chain of Banach spaces:

$$\mathcal{B}_0 \hookrightarrow \mathcal{B}_1 \hookrightarrow \mathcal{B}_2 \hookrightarrow \dots \hookrightarrow \mathcal{B}_{r+1} \hookrightarrow \mathbb{L}^1$$

such that the operators P_t are quasi-compact on $\mathcal{B}_1, \dots, \mathcal{B}_r$ and such that $t \mapsto P_t$ admits a Taylor expansion with error in $\mathcal{O}(t^m)$ in $\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_r)$, and with an error in $o(t^m)$ in $\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_{r+1})$. Compared to [26], we establish (see Proposition 4.1, Theorem 5.5 and Corollary 5.6):

- Taylor expansions for eigenprojectors
 - take the first space \mathcal{B}_0 to be the space of constant functions which is preserved just by P , we do not assume that P_t acts on \mathcal{B}_0 (gain of space at the beginning of the chain),
 - replace \mathcal{C}^r -smoothness of $t \mapsto P_t$ and $t \mapsto \Pi_t$ by Taylor expansions with error in $\mathcal{O}(t^r)$ in $\mathcal{L}(\mathcal{B}_0, \mathcal{B}_r)$, and with an error in $o(t^r)$ in $\mathcal{L}(\mathcal{B}_0, \mathcal{B}_{r+1})$ (gain of space all along the chain, gain of space at the end of the chain for the estimate in $\mathcal{O}(t^r)$),⁴
 - choose the spaces \mathcal{B}_j so that $\mathcal{B}_j \hookrightarrow \mathbb{L}^{\frac{r+1}{j}}$,
- Taylor expansion for the dominating eigenvalue by a key weak compactness argument
 - The $o(t^{r+1})$ -Taylor expansion of the dominating eigenvalue will follow from an $o(t^r)$ -Taylor expansion of $\mathbb{E}_\mu[\Pi_t(\mathbf{1})]$ and from $o(t^{r+1})$ Taylor expansions of $\mathbb{E}_\mu[(e^{itf} - 1)]$ and of $\mathbb{E}_\mu[(e^{itf} - 1)(\Pi_t - \Pi_0)(\mathbf{1})]$ (Fact 3.2). The main issue is to prove this last expansion. But, to this end, only r -order Taylor expansions of $e^{itf} - 1$ and $(\Pi_t - \Pi_0)(\mathbf{1})$ will be needed.
 - $o(t^{r+1})$ -Taylor expansion of $\mathbb{E}_\mu[(e^{itf} - 1)(\Pi_t - \Pi_0)(\mathbf{1})]$: make an $o(t^r)$ -Taylor expansion of $(e^{itf} - 1)$ and study individually each term $\mathbb{E}_\mu[\frac{(itf)^k}{k!}(\Pi_t - \Pi_0)(\mathbf{1})]$, for $k = 1, \dots, r$.
 - $o(t^{r+1})$ -Taylor expansion of $\mathbb{E}_\mu[\frac{(itf)^k}{k!}(\Pi_t - \Pi_0)(\mathbf{1})]$: first use an $\mathcal{O}(t^{r+1-k})$ -Taylor expansion of $(\Pi_t - \Pi_0)(\mathbf{1}) \in \mathcal{B}_{r+1-k}$ to prove weak compactness in the reflexive space $\mathbb{L}^{\frac{r+1}{r+1-k}}$, second use $o(t^{r+1-k})$ -Taylor expansions in $\mathcal{B}_{r+1} \hookrightarrow \mathbb{L}^1(\mu)$ to identify the weak limit. Conclude.
 - gather the terms and conclude the desired $o(t^{r+1})$ -Taylor expansion of $\mathbb{E}_\mu[(e^{itf} - 1)(\Pi_t - \Pi_0)(\mathbf{1})]$.
 - The $o(t^{r+1})$ -Taylor expansion of $\mathbb{E}_\mu[(e^{itf} - 1)]$ and the $o(t^r)$ -Taylor expansion of $\mathbb{E}_\mu[\Pi_t(\mathbf{1})]$ follow directly from the corresponding Taylor expansions of respectively $(e^{itf} - 1)$ in $\mathbb{L}^1(\mu)$ and $\Pi_t(\mathbf{1})$ in $\mathcal{B}_{r+1} \hookrightarrow \mathbb{L}^1(\mu)$.

3. QUASI-COMPACTNESS AND PERTURBATION

Let \mathcal{B}_1 be complex Banach space. We write $\mathcal{L}(\mathcal{B}_1)$ for the set of continuous linear operators on \mathcal{B}_1 . For any continuous linear operator $P \in \mathcal{B}_1$, we write $\rho(P)$ for its spectral radius and $\rho_{ess}(P)$ for its essential spectral radius. Recall that the operator P is said to be quasi-compact if $\rho_{ess}(P) < \rho(P)$.

Theorem 3.1 (Browder [4]). *Let $P \in \mathcal{L}(\mathcal{B}_1)$ be quasi-compact. Let $r \in]\rho_{ess}(P), \rho(P)[$. Then the spectrum of P outside $B(0, r)$ consists in a finite number of eigenvalues $\lambda_1, \dots, \lambda_s$ (isolated*

⁴Let us indicate for completeness that, in our examples, it has been proved in [26] that $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_0, \mathcal{B}_{r+1})$ is \mathcal{C}^r . Indeed, in these situations, there is a continuum of Banach spaces and so the space between \mathcal{B}_r and \mathcal{B}_{r+1} (to get $o(t^r)$) can be spread along the chain to add spaces $\tilde{\mathcal{B}}_j$, up to slightly moving the spaces $\mathcal{B}_0, \dots, \mathcal{B}_r$. Nevertheless we do not need this smoothness but just the Taylor expansion with error in $o(t^r)$.

in the spectrum of P) and there exist positive integers m_1, \dots, m_s such that

$$\mathcal{B}_1 = \bigoplus_{j=0}^s \mathcal{E}_j, \quad \text{with } P(\mathcal{E}_j) \subset \mathcal{E}_j \quad \text{and} \quad \|P|_{\mathcal{E}_0}^n\| = \mathcal{O}(r^n),$$

with, for all $j = 1, \dots, s$, $\mathcal{E}_j := \ker(P - \lambda_j \text{Id})^{m_j}$ and $\dim(\mathcal{E}_j) < \infty$. For every $j = 0, \dots, s$, there exists a continuous linear projection $\Pi_{[j]} : \mathcal{B}_1 \rightarrow \mathcal{E}_j$ such that

$$\sum_{j=0}^s \Pi_{[j]} = \text{Id}, \quad P\Pi_{[j]} = \Pi_{[j]}P, \quad \Pi_{[j]}\Pi_{[\ell]} = \delta_{j,\ell}\Pi_{[j]}$$

and

$$\forall n \geq 0, \quad P^n \Pi_{[j]} = \frac{1}{2i\pi} \int_{\Gamma_j} z^n (z \text{Id} - P)^{-1} dz, \quad (3)$$

with Γ_0 an oriented circle $\mathcal{C}(0, r_0)$ containing no λ_j and with $r_0 < r$, and with, for $j = 1, \dots, s$, Γ_j an oriented circle $\mathcal{C}(\lambda_j, r_j) \subset \mathbb{C}$ such that λ_j is the only spectral value of P contained in the closed disk $\mathcal{D}(\lambda_j, r_j]$.

If moreover $m_j = 1$ for every $j = 1, \dots, s$, then

$$P^n = \sum_{j=1}^s \lambda_j^n \Pi_{[j]} + P^n \Pi_{[0]}, \quad \text{with } \|P^n \Pi_{[0]}\| = \mathcal{O}(r^n).$$

We consider now a quasi-compact operator $P \in \mathcal{L}(B)$ with simple peripheral spectrum and a family of quasi-compact operators $(P_t)_{|t| < \delta}$ such that $P_0 = P$ and admitting the same type of decomposition as P :

$$P_t^n = \sum_{j=1}^s \lambda_{j,t}^n \Pi_{[j],t} + N_t^n, \quad \text{with } \|N_t^n\| = \mathcal{O}(r^n),$$

with $\lambda_{j,t}$ contained in the open disk $\mathcal{D}(\lambda_j, r_j[$. We will use the Keller-Liverani perturbation theorem recalled at the end of this section to prove that the family of operators we are considering satisfies this property.

Due to Theorem 3.1, the regularity in t of the eigenelements $\Pi_{[j],t}$ of P_t will follow from the regularity in t of the resolvent $(z \text{Id} - P_t)^{-1}$ uniformly on $z \in \bigcup_{j=0}^s \Gamma_j$. Such results will be stated in Proposition 4.1 thanks to Theorem 3.3 (Keller-Liverani perturbation theorem).

As explained in Section 2, we will deduce from this an higher order Taylor expansion for the dominating eigenvalues due to the following key formula.

Fact 3.2 (see [1], or [21]). *In this context (assuming $m_j = 1$), if $v_j \in \ker(P - \lambda_j \text{Id})$ and $\varphi_j \in \ker(P^* - \lambda_j \text{Id})$ are such that $\varphi_j(v_j) = 1$ and that $t \mapsto \varphi_j \circ \Pi_{[j],t}(v_j)$ is continuous, then $P_t(\Pi_{[j],t}(v_j)) = \lambda_{j,t}(\Pi_{[j],t}(v_j))$ and, for t small enough, $\varphi_j(\Pi_{[j],t}(v_j)) \neq 0$, and so*

$$\lambda_{j,t} - \lambda_j = \frac{\varphi_j((P_t - P)(\Pi_{[j],t}(v_j)))}{\varphi_j(\Pi_{[j],t}(v_j))},$$

where we used the fact that $\varphi_j \circ P = \lambda_j \varphi_j$. Since $\Pi_{[j]} = \varphi_j(\cdot)v_j$, it follows also that

$$\begin{aligned} \lambda_{j,t} - \lambda_j &= \varphi_j((P_t - P)v_j) + \frac{\varphi_j((P_t - P)(\Pi_{[j],t}(v_j) - \Pi_{[j]}(\Pi_{[j],t}(v_j))))}{\varphi_j(\Pi_{[j],t}(v_j))} \\ &= \varphi_j((P_t - P)v_j) + \frac{\varphi_j((P_t - P)(\text{Id} - \Pi_{[j]})(\Pi_{[j],t} - \Pi_{[j]})(v_j))}{\varphi_j(\Pi_{[j],t}(v_j))}. \end{aligned}$$

In Theorem 3.3 below, an auxiliary space \mathcal{B}_2 is used to study the spectral properties of a family of continuous linear operators acting on \mathcal{B}_1 .

Theorem 3.3 (Keller-Liverani perturbation theorem, see [31] and [11]). *Let V be a neighbourhood of 0 in \mathbb{R}^d . Let $(P_t)_{t \in V}$ be a family of continuous linear operators on a Banach space $(\mathcal{B}_1, \|\cdot\|_{(1)})$. Let $(\mathcal{B}_2, \|\cdot\|_{(2)})$ be a Banach space such that $\mathcal{B}_1 \hookrightarrow \mathcal{B}_2$. Assume that there exist positive real numbers $C > 0$, $R > \rho(P_0)$ and $r \in]\rho_{\text{ess}}(P_0), R[$ such that $(P_t)_{t \in V}$ satisfies the following uniform Doeblin-Fortet type inequality⁵*

$$\forall t \in V, \forall f \in \mathcal{B}_1, \forall n \geq 0 \quad \|P_t^n f\|_{(1)} \leq Cr^n \|f\|_{(1)} + R^n \|f\|_{(2)}. \quad (4)$$

Assume moreover that $P_0 \in \mathcal{L}(\mathcal{B}_2)$. Let $\varepsilon \in]0, R - r/2[$. Assume $P_0 \in \mathcal{L}(\mathcal{B}_1)$ has eigenvalues $\lambda_{[1]}, \dots, \lambda_{[m]}$ of modulus strictly larger than $r + 2\varepsilon$ (here the eigenvalues are repeated with their multiplicities), and with no other eigenvalues of modulus larger than or equal to $r + \varepsilon$. Assume moreover that $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ is continuous at 0. Then there exists a neighbourhood U of 0 contained in V such that, for any $t \in U$, $P_t \in \mathcal{L}(\mathcal{B}_1)$ is quasi-compact with $\rho_{\text{ess}}(P_t) < r + \varepsilon$ and with eigenvalues $\lambda_{[1],t}, \dots, \lambda_{[m],t}$ of modulus strictly larger than $r + \varepsilon$ and such that, furthermore,

$$\sup_{t \in U, z \in \mathbb{C} : r + \varepsilon < |z| < R + \varepsilon, \inf_j |z - \lambda_{[j]}| > \varepsilon} \|(z \text{Id} - P_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_1)} < +\infty,$$

$$\lim_{t \rightarrow 0} \sup_{z \in \mathbb{C} : r + \varepsilon < |z| < R + \varepsilon, \inf_j |z - \lambda_{[j]}| > \varepsilon} \|(z \text{Id} - P_t)^{-1} - (z \text{Id} - P_0)^{-1}\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} = 0,$$

and, for all $j=1, \dots, m$,

$$\lim_{t \rightarrow 0} \lambda_{[j],t} = \lambda_j,$$

$$\forall t \in U, \quad \dim \sum_{i: \lambda_i = \lambda_j} \bigcup_{k \geq 0} \ker(P_t - \lambda_{i,t} \text{Id})^k = \dim \bigcup_{k \geq 0} \ker(P_t - \lambda_{[j]} \text{Id})^k,$$

$$\lim_{t \rightarrow 0} \|\Pi_{[0]} - \Pi_{[0],t}\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} = 0 \quad \text{and} \quad \lim_{t \rightarrow 0} \left\| \Pi_{[j]} - \sum_{i: \lambda_i = \lambda_j} \Pi_{[i],t} \right\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} = 0,$$

where $\Pi_{[i],t}$ are the projectors associated to $P_t \in \mathcal{L}(\mathcal{B}_1)$ and $\lambda_{i,t}$ as considered in Theorem 3.1.

This theorem ensures that

$$P_t^n = \left(\sum_{j=1}^m P_t^n \Pi_{[j],t} \right) + N_t^n,$$

with $N_t := P_t \circ \Pi_{[0],t}$ satisfies $\sup_t \|N_t^n\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}((r + \varepsilon)^n)$. It will be crucial to notice that in the particular case where all the characteristic spaces $\bigcup_{k \geq 0} \ker(P_t - \lambda_{i,t} \text{Id})^k$ consist only of eigenvectors, then this decomposition can be simplified in

$$P_t^n = \left(\sum_{j=1}^m \lambda_{[j],t}^n \Pi_{[j],t} \right) + N_t^n.$$

⁵We call it Doeblin-Fortet type inequality in reference to [6].

4. TAYLOR EXPANSIONS FOR THE RESOLVANT AND EIGENPROJECTORS

The continuity in $t \in \mathbb{R}^d$ of the eigenprojectors stated in Theorem 3.3 appears as a consequence of the continuity in t of the Resolvent $R_{z,t} := (z \text{Id} - P_t)^{-1}$ of P_t , due to Formula (3) (taken with $n = 0$) providing an expression of $\Pi_{[j],t}$ as an integral of $(z \text{Id} - P_t)^{-1}$. In the present section, we establish higher order Taylor expansions for the Resolvent $t \mapsto R_{z,t} = (z \text{Id} - P_t)^{-1}$, that will imply immediately the corresponding Taylor expansions for $t \mapsto \Pi_{[j],t}$ thanks to (3). This section is devoted to the proof of the next result providing Taylor expansions of the resolvent. This result contains estimates of orders 0 and 1 (useful to establish the convergence in distribution to stable distributions for non square integrable observables) but also higher order Taylor expansions for the resolvent. Since our result holds true in multi-dimension, we need to introduce some different notions related to the multilinear forms appearing in multi-dimensional Taylor expansions.

We write $\mathcal{B}_1^d = (\mathcal{B}_1)^d$ for the set of d -dimensional vectors with entries in \mathcal{B}_1 , and more generally $\mathcal{B}_\ell^{d \otimes k}$ for the set of k -linear maps on $(\mathbb{R}^d)^k$ with values in \mathcal{B}_ℓ , and we identify its elements \mathcal{H} with $\mathcal{H} = (h_{i_1, \dots, i_k})_{i_1, \dots, i_k=1, \dots, d}$. We write $t^{\otimes m}$ for $(t_{i_1} \dots t_{i_m})_{i_1, \dots, i_m=1, \dots, d}$ and \cdot for the scalar product. Finally, for any $\mathcal{H}^{(1)} \in \mathcal{L}(\mathcal{B}_i, \mathcal{B}_j^{d \otimes m})$ and $\mathcal{H}^{(2)} \in \mathcal{L}(\mathcal{B}_j, \mathcal{B}_\ell^{d \otimes k})$, we write

$$\mathcal{H}^{(2)} \otimes \mathcal{H}^{(1)} := \left(\frac{1}{(m+k)!} \sum_{\sigma \in \mathfrak{S}_{m+k}} \mathcal{H}_{i_{\sigma(1)}, \dots, i_{\sigma(m)}}^{(2)} H_{i_{\sigma(m+1)}, \dots, i_{\sigma(m+k)}}^{(1)} \right)_{i_1, \dots, i_{m+k}=1, \dots, d},$$

where we write as usual \mathfrak{S}_{m+k} for the set of permutations of $\{1, \dots, m+k\}$. Note that, with these notations, $\mathcal{H}^{(2)} \otimes \mathcal{H}^{(1)} \in \mathcal{L}(\mathcal{B}_i, \mathcal{B}_\ell^{d \otimes m+k})$ and that $(\mathcal{H}^{(2)} \otimes \mathcal{H}^{(1)}) \cdot t^{\otimes(k+m)} = (\mathcal{H}^{(2)} \cdot t^{\otimes k})(\mathcal{H}^{(1)} \cdot t^{\otimes m})$. In dimension 1 (when $d = 1$), \otimes as well as \cdot both correspond to the usual product, $t^{\otimes m}$ simply to t^m and $\mathcal{B}_\ell^{d \otimes k}$ to \mathcal{B}_ℓ .

Proposition 4.1. *Let $\delta > 0$, $\Gamma \subset \mathbb{C}$ and r be a nonnegative integer. Let $(\mathcal{B}_j, \|\cdot\|_{(j)})$, $j \in \{0, \dots, r+1\}$ be a chain of $(r+2)$ Banach spaces, increasing in the sense that $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$ for all $j = 0, \dots, r$. Let $(P_t)_t$ be a family of linear operators acting continuously on $\mathcal{B}_1, \dots, \mathcal{B}_{r+1}$ (for all $t \in \mathbb{R}^d$ such that $|t| < \delta$). Assume $P = P_0$ acts continuously on \mathcal{B}_0 and that*

$$K_0 := \sup_{j=1, \dots, r+1} \sup_{|t| < \delta} \sup_{z \in \Gamma} \left(\|(z \text{Id} - P_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_j)} + \|(z \text{Id} - P)^{-1}\|_{\mathcal{L}(\mathcal{B}_0)} \right) < \infty. \quad (5)$$

Let $(P_0^{(k)})_{k=0, \dots, r}$ be a family of operators such that, for all $k = 0, \dots, r$, $P_0^{(k)} \in \bigcap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k}^{d \otimes k})$ and such that $P_0^{(0)} = P_0$. Then

(A) for all $j = 0, \dots, r$,

$$\left\| (z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} \right\|_{\mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+1})} \leq K_0^2 \|P_t - P\|_{\mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+1})}. \quad (6)$$

(B) If $r = 1$ then, setting $P'_0 = P_0^{(1)}$,

$$\begin{aligned} & \left\| (z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} - (z \text{Id} - P)^{-1} (P'_0 \cdot t) (z \text{Id} - P)^{-1} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_2)} \\ & \leq K_0^2 \|P_t - P - P'_0 \cdot t\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_2)} + |t| K_0^3 \|P_t - P\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} \|P'_0\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_1^d)}. \end{aligned} \quad (7)$$

(C) If ⁶

$$\forall j = 0, \dots, r, \quad \left\| P_t - \sum_{k=0}^{r-j} \frac{P_0^{(k)}}{k!} \cdot t^{\otimes k} \right\|_{\mathcal{L}(\mathcal{B}_j, \mathcal{B}_r)} \leq K_1 |t|^{r-j},$$

⁶Note that \mathcal{B}_{r+1} does not play any role in this result and thus we can take $\mathcal{B}_{r+1} = \mathcal{B}_r$.

then there exists a constant \widetilde{K}_r which is given by a polynomial expression in K_0, K_1 such that

$$\left\| (z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} - \sum_{j=1}^r R_{z,0}^{(j)} t^{\otimes j} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_r)} \leq \widetilde{K}_r |t|^r,$$

for all $|t| < \delta$ and all $z \in \Gamma$, with

$$R_{z,0}^{(j)} := (z \text{Id} - P)^{-1} \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1 : k_1 + \dots + k_\ell = j} \mathcal{A}_{k_\ell} \otimes \dots \otimes \mathcal{A}_{k_1},$$

and $\mathcal{A}_\ell := \frac{P_0^{(\ell)}}{\ell!} (z \text{Id} - P)^{-1} \in \bigcap_{j=0}^{r-\ell} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+\ell}^{d^{\otimes \ell}})$.

(D) If for all $j = 0, \dots, r$, $\left\| P_t - \sum_{k=0}^{r-j} \frac{P_0^{(k)}}{k!} t^{\otimes k} \right\|_{\mathcal{L}(\mathcal{B}_j, \mathcal{B}_{r+1})} = o(|t|^{r-j})$, then, on $\mathcal{L}(\mathcal{B}_0, \mathcal{B}_{r+1})$,

$$(z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} - \sum_{j=1}^r R_{z,0}^{(j)} t^{\otimes j} = o(|t|^r),$$

uniformly in $|t| < \delta$ and $z \in \Gamma$, with the same notations as in the previous item.

It may be worthwhile to note that

$$R_{z,0}^{(j)} t^{\otimes j} = (z \text{Id} - P)^{-1} \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1 : k_1 + \dots + k_\ell = j} (\mathcal{A}_{k_\ell} t^{\otimes k_\ell}) \dots (\mathcal{A}_{k_1} t^{\otimes k_1}).$$

Since, in our examples, P_t will have the form $P(e^{it.f.})$, the Taylor expansions of $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+m})$ will be proved using Taylor expansions in t of $t \mapsto (e^{it.f.}) \in \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+m})$ and the operators $P_0^{(\ell)}$ will have the form $P_0^{(\ell)} = P((if)^{\otimes \ell})$.

Remark 4.2. In practice we will prove the first part of Assumption (5) (about the control of $\|(z \text{Id} - P_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_j)}$) by applying Theorem 3.3 with Banach spaces $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$ for all $j = 1, \dots, r+1$ and thus we will need a $(r+3)$ -th Banach space \mathcal{B}_{r+2} . For the second part of Assumption (5) (control of $\|(z \text{Id} - P)^{-1}\|_{\mathcal{L}(\mathcal{B}_0)}$), a useful idea (used several times in applications) will be to take the eigenspace associated to the dominating eigenvalue (in our applications, it will be the space of constant functions).

We could have stated the previous lemma in a much more general way by replacing $P_0^{(k)} t^{\otimes k}$ by $a_{k,t}$ and $\mathcal{A}_k t^{\otimes k}$ by $\frac{a_{k,t}}{k!} (z \text{Id} - P)^{-1} \in \bigcap_{j=0}^{r-\ell} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+\ell})$. We have not chosen this presentation since we do not have application in mind, except maybe the case of convergence to stable distribution, but for which in practice Item (A) is enough (see Proposition 5.11 and Example 5.12).

Proof of Proposition 4.1. We will use the following key identity

$$(z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} = (z \text{Id} - P_t)^{-1} (P_t - P) (z \text{Id} - P)^{-1}.$$

Observe that Item (A) is a direct consequence of this identity. Analogously

$$\begin{aligned} & (z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} - (z \text{Id} - P)^{-1} (P'_0 t) (z \text{Id} - P)^{-1} \\ &= (z \text{Id} - P_t)^{-1} (P_t - P) (z \text{Id} - P)^{-1} - (z \text{Id} - P)^{-1} (P'_0 t) (z \text{Id} - P)^{-1} \\ &= (z \text{Id} - P_t)^{-1} (P_t - P - P'_0 t) (z \text{Id} - P)^{-1} + \left[(z \text{Id} - P_t)^{-1} - (z \text{Id} - P)^{-1} \right] (P'_0 t) (z \text{Id} - P)^{-1} \end{aligned}$$

The first term has norm less than $K_0^2 \|P_t - P - P'_0 t\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_2)}$ in $\mathcal{L}(\mathcal{B}_0, \mathcal{B}_2)$ and the second one can be rewritten

$$(z \text{Id} - P_t)^{-1} (P_t - P) (z \text{Id} - P)^{-1} (P'_0 t) (z \text{Id} - P)^{-1},$$

which ends the proof of (B).

To establish (C), we prove by induction on $m = 1, \dots, r$ that

$$\left\| R_{z,t} - R_{z,0} - \sum_{j=1}^{m-1} R_{z,0}^{(j)} \cdot t^{\otimes j} \right\|_{\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_r)} \leq \widetilde{K}_m |t|^m, \quad (8)$$

for all $|t| < \delta$ and $z \in \Gamma$. Due to Item (A) applied with $j = r - 1$, (8) holds true for $m = 1$ with $\widetilde{K}_1 = K_1 K_0^2$. for all $|t| < \delta$ and $z \in \Gamma$.

Let $N = 2, \dots, r$. Assume (8) holds true for all $m = 0, \dots, N - 1$, and let us prove it holds also true for $m = N$. Observe that

$$\begin{aligned} R_{z,t} - R_{z,0} &= R_{z,t}(P_t - P)R_{z,0} \\ &= R_{z,t} \sum_{k=1}^{N-1} \frac{P_0^{(k)}}{k!} \cdot t^{\otimes k} R_{z,0} + \mathcal{O}(|t|^N) \\ &= \sum_{k=1}^{N-1} (R_{z,t} - R_{z,0}) \mathcal{A}_k \cdot t^{\otimes k} + R_{z,0} \sum_{k=1}^{N-1} \mathcal{A}_k \cdot t^{\otimes k} + \mathcal{O}(|t|^N), \end{aligned} \quad (9)$$

in $\mathcal{L}(\mathcal{B}_{r-N}, \mathcal{B}_r)$, with $\mathcal{O}(|t|^N)$ bounded by $K_0^2 K_1$ (uniformly in t, z). Recall that $\mathcal{A}_k \in \mathcal{L}(\mathcal{B}_{r-N}, \mathcal{B}_{r-N+k}^{d^{\otimes k}})$. It follows from the inductive hypothesis that, for any $k = 1, \dots, N - 1$,

$$R_{z,t} - R_{z,0} = \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1 : k_1 + \dots + k_\ell \leq N-k-1} R_{z,0}(\mathcal{A}_{k_\ell} \cdot t^{\otimes k_\ell}) \cdots (\mathcal{A}_{k_1} \cdot t^{\otimes k_1}) + \mathcal{O}(|t|^{N-k})$$

in $\mathcal{L}(\mathcal{B}_{r-N+k}, \mathcal{B}_r)$ uniformly in t, z . Using this formula in the first sum in the right hand side of (9) ends the induction and so the proof of Item (C).

It remains finally to prove Item (D). To this end, it is enough to prove by induction on $m = 0, \dots, r$ that, on $\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_{r+1})$,

$$R_{z,t} - R_{z,0} - \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1 : k_1 + \dots + k_\ell \leq m} (z \text{Id} - P)^{-1} (\mathcal{A}_{k_\ell} \cdot t^{\otimes k_\ell}) \cdots (\mathcal{A}_{k_1} \cdot t^{\otimes k_1}) = o(|t|^m) \quad (10)$$

uniformly in $|t| < \delta$ and $z \in \Gamma$. This is true for $m = 0$ since

$$R_{z,t} - R_{z,0} = R_{z,t}(P_t - P)R_{z,0} = o(1),$$

on $\mathcal{L}(\mathcal{B}_r, \mathcal{B}_{r+1})$ uniformly in t, z . Fix $N = 2, \dots, r$. Assume (10) for all $m = 0, \dots, N - 1$ and let us prove it holds also true for $m = N$. Observe that

$$\begin{aligned} R_{z,t} - R_{z,0} &= R_{z,t}(P_t - P)R_{z,0} \\ &= R_{z,t} \sum_{k=1}^N \frac{P_0^{(k)}}{k!} \cdot t^{\otimes k} R_{z,0} + o(|t|^N) \\ &= \sum_{k=1}^N (R_{z,t} - R_{z,0}) \mathcal{A}_k \cdot t^{\otimes k} + R_{z,0} \sum_{k=1}^N \mathcal{A}_k \cdot t^{\otimes k} + o(|t|^N), \end{aligned}$$

in $\mathcal{L}(\mathcal{B}_{r-N}, \mathcal{B}_{r+1})$ uniformly in t, z . Recall that $\mathcal{A}_k \in \mathcal{L}(\mathcal{B}_{r-N}, \mathcal{B}_{r-N+k}^{d^{\otimes k}})$. It follows from the inductive hypothesis that

$$R_{z,t} - R_{z,0} = \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1 : k_1 + \dots + k_\ell \leq N-k} R_{z,0}(\mathcal{A}_{k_\ell} \cdot t^{\otimes k_\ell}) \cdots (\mathcal{A}_{k_1} \cdot t^{\otimes k_1}) + o(|t|^{N-k})$$

in $\mathcal{L}(\mathcal{B}_{r-N+k}, \mathcal{B}_{r+1})$ uniformly in t, z . This ends the proof of Item (D) and so of Proposition 4.1. \square

To make easier the comparison with previous works, let us recall the result of [26, Appendix A] about \mathcal{C}^r -smoothness.

Proposition 4.3 ([26]). *Assume there exists a double chain of Banach spaces*

$$\mathcal{B}_0 \hookrightarrow \tilde{\mathcal{B}}_0 \hookrightarrow \mathcal{B}_1 \hookrightarrow \tilde{\mathcal{B}}_1 \hookrightarrow \dots \hookrightarrow \mathcal{B}_r \hookrightarrow \tilde{\mathcal{B}}_r.$$

Assume $(P_t)_{t \in U}$ is a family of linear operators acting continuously on all these Banach spaces (with U an open subset of \mathbb{R}^d) such that

$$\sup_{j=0,\dots,r} \sup_{t \in U} \sup_{z \in \Gamma} \left(\|(z \text{Id} - P_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_j)} + \|(z \text{Id} - P_t)^{-1}\|_{\mathcal{L}(\tilde{\mathcal{B}}_j)} \right) < \infty,$$

and such that $t \mapsto P_t \in \bigcap_{j=0}^r \mathcal{L}(\mathcal{B}_j, \tilde{\mathcal{B}}_j)$ is continuous on $U \subset \mathbb{R}^d$ and such that $t \mapsto P_t \in \bigcap_{j=0}^{r-m} \mathcal{L}(\tilde{\mathcal{B}}_j, \mathcal{B}_{j+m})$ is \mathcal{C}^m . Then $t \mapsto (z \text{Id} - P_t)^{-1} \in \mathcal{L}(\mathcal{B}_0, \tilde{\mathcal{B}}_r)$ is \mathcal{C}^r , with derivatives at 0 given by the Taylor expansion established in Proposition 4.1.

Proof. This result follows directly [26, Proposition A] applied with $I = \bigcup_{j=1}^m \{\mathcal{B}_j, \tilde{\mathcal{B}}_j\}$, $T_0(\mathcal{B}_j) = \tilde{\mathcal{B}}_j$, $T_1(\tilde{\mathcal{B}}_j) = \mathcal{B}_{j+1}$ (up to identify I with a subset of \mathbb{R}). \square

5. EXPANSIONS OF FOURIER EIGENPROJECTORS AND EIGENVALUES IN MARKOVIAN OR DYNAMICAL CONTEXTS

5.1. General context and toy model. In this section, we will see how the general results of Section 4 can be implemented to study dynamical or markovian random walks $(S_n)_{n \geq 1}$ defined as follows.

Hypothesis 5.1. *Let $(\Omega, \mathcal{F}, \mu)$ and $(E, \mathcal{T}, \mathbf{P})$ be two probability spaces.*

- (I) *either $X_n = T^n$ where $T : \Omega \rightarrow \Omega$ is a μ -preserving transformation, with transfer operator P and $f : \Omega \times E \rightarrow \mathbb{R}^d$ is a measurable $\mu \otimes \mathbf{P}$ -centered function. We consider ν a probability measure on Ω absolutely continuous with density h with respect to μ . To unify notations with the markovian setting, we also set $\mathcal{P}_\mu = \mu$ and $f(x, y, \omega) := f(y, \omega)$.*
- (II) *or $(X_n)_{n \geq 0}$ is a Markov chain (identified with the canonical Markov chain) with values in Ω and with stationary measure μ and $f : \Omega \times \Omega \times E \rightarrow \mathbb{R}^d$ is a measurable function. Let ν be the distribution of X_0 (i.e. the initial distribution of the Markov chain). We set \mathcal{P}_ν for the Markov distribution with transition operator P and initial probability measure ν . We assume that $((x_k)_k, \omega) \mapsto f(x_0, x_1, \omega)$ is $\mathcal{P}_\mu \otimes \mathbf{P}$ -centered.*

We set

$$P_t(h)(x) = \int_E P \left(e^{it \cdot f(x, \cdot, \omega)} h(\cdot) \right) (x) d\mathbf{P}(\omega),$$

and $S_n := \sum_{k=1}^n Y_k$ with $Y_k := f(X_{k-1}, X_k, Z_k)$ where Z_i are i.i.d. random variables independent of $(X_k)_{k \geq 0}$ and with common distribution \mathbf{P} .⁷

In the dynamical setting (I), identifying $H(x, \omega)$ with $H(x)$:

$$\mathbb{E}_\mu [g P_t^n(hG)] = \mathbb{E}_{\mu \otimes \mathbf{P}^{\otimes n}} \left[(g \circ T^n) e^{it \cdot S_n} \cdot hG \right] = \mathbb{E}_{\nu \otimes \mathbf{P}^{\otimes n}} \left[(g \circ T^n) e^{it \cdot S_n} G \right]. \quad (11)$$

⁷Let us indicate for completeness that Markov random walks are usually defined as a Markov chain $(X_n, \tilde{S}_n)_{n \geq 0}$ satisfying

$$\mathbb{E} [h(X_n, \tilde{S}_n) | X_{n-1} = x, \tilde{S}_{n-1} = y] = \mathbb{E} [h(X_n, y + \tilde{S}_n) | X_{n-1} = x, \tilde{S}_{n-1} = 0].$$

and that we prove in appendix that these two definitions are equivalent.

In the Markovian setting (II), $P_t(h)(x) = \mathbb{E} \left[e^{it \cdot Y_1} h(X_1) \middle| X_0 = x \right]$, and so

$$\mathbb{E}_\nu [g P_t^n(h)] = \mathbb{E}_{\mathcal{P}_\nu \otimes \mathbf{P}^{\otimes n}} \left[g(X_0) e^{it \cdot S_n} h(X_n) \right]. \quad (12)$$

Before considering applications seen in Theorem 1.1, we will state results in the general dynamical or markovian context and will illustrate them on the following toy model of Knudsen gas considered in [3], which is one of the simplest example of uniformly geometrically ergodic Markov chains (i.e. satisfying (1) with $\mathcal{B}_1 = L^\infty$ the set of uniformly bounded complex-valued functions on Ω).

Example 5.2 (a Toy model of Knudsen gas). *Let r be a nonnegative integer. Let $(\Omega, \mathcal{F}, \mu)$ be a probability space and $\alpha \in]0, 1[$. Consider the Markov chain $(X_n)_n$ with transition operator given by*

$$P(h) = \alpha h + (1 - \alpha) \mathbb{E}_\mu[h], \quad \text{i.e.} \quad P(h - \mathbb{E}_\mu[h]) = \alpha(h - \mathbb{E}_\mu[h]).$$

This Markov chain describes the evolution of a process which, at each step, remains the same with probability α and changes to an independent copy of distribution μ with probability $1 - \alpha$. We consider also $Y_n = f(X_n)$, with $f : \Omega \rightarrow \mathbb{R}^d$ centered and admitting moments of order $(r+1)$ with respect both to the invariant distribution μ and to the initial distribution ν . Then P_t is given by

$$P_t(h)(x) = \alpha e^{itf(x)} h(x) + (1 - \alpha) \mathbb{E}_\mu[e^{itf} h].$$

More generally, we will consider the context described in the following:

Remark 5.3. *Assume Hypothesis 5.1 and that P is geometrically ergodic on some Banach space $\mathcal{B}_1 \hookrightarrow \mathbb{L}^1(\mu)$ containing the constant functions, so that $\|P^n - \mathbb{E}_\mu\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta^n)$ (for some $\vartheta \in]0, 1[$). Let $\vartheta_1 \in]\vartheta, 1[$. Assume moreover that the assumptions of Theorem 3.3 are satisfied for this choice of (P_t, \mathcal{B}_1) and for some \mathcal{B}_2 . Then there exists $\delta_0 > 0$ such that, in $\mathcal{L}(\mathcal{B}_1)$,*

$$\forall |t| < \delta_0, \quad P_t^n = \lambda_t^n \Pi_t + N_t^n, \quad \text{with} \quad \sup_{|t| < \delta_0} \|N_t^n\|_{\mathcal{B}_1} = \mathcal{O}(\vartheta_1^n), \quad (13)$$

with $t \mapsto \Pi_t \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ and $t \mapsto \lambda_t \in \mathbb{C}$ continuous and $\Pi_0 = \mathbb{E}_\mu[\cdot] \mathbf{1}$ and $\lambda_0 = 1$. Moreover Theorems 3.1 and 3.3 ensure that Π_t and N_t are given by the following formulas

$$\Pi_t := \frac{1}{2i\pi} \int_{\Gamma_1} (z \text{Id} - P_t)^{-1} dz, \quad N_t^n := \frac{1}{2i\pi} \int_{\Gamma_0} z^n (z \text{Id} - P_t)^{-1} dz, \quad (14)$$

with Γ_1 the oriented circle $\mathcal{C}(1, \delta)$ and Γ_0 the oriented circle $\mathcal{C}(0, a)$, with $\vartheta_1 < a < a + \delta < 1$ and that

$$\sup_{|t| < \delta_0, z \in \Gamma_1 \cup \Gamma_0} \|(z \text{Id} - P_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_1)} < \infty,$$

with $t \mapsto R_{z,t} = (z \text{Id} - P_t)^{-1} \in \mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ continuous on $\{t \in \mathbb{R}^d : |t| < \delta_0\}$ (uniformly in $z \in \Gamma_1 \cup \Gamma_0$).

Example 5.4 (Knudsen gas). *We consider again the Knudsen gas introduced in Example 5.2. Note that*

$$\forall p \in [1, +\infty], \quad \|P^n(h) - \mathbb{E}_\mu[h]\|_{\mathbb{L}^p(\mu)} = \alpha^n \|h - \mathbb{E}_\mu[h]\|_{\mathbb{L}^p(\mu)},$$

and if $\nu \neq \mu$, it is worthwhile to notice that we also have

$$\forall \gamma \in [0, +\infty[, \quad \left\| \frac{P^n(h) - \mathbb{E}_\mu[h]}{(1 + |f|)^\gamma} \right\|_\infty = \alpha^n \left\| \frac{h - \mathbb{E}_\mu[h]}{(1 + |f|)^\gamma} \right\|_\infty \leq (1 + \mathbb{E}_\mu[(1 + |f|)^\gamma]) \alpha^n \left\| \frac{h}{(1 + |f|)^\gamma} \right\|_\infty.$$

Thus, since $|P_t^n(h)| \leq P^n(|h|)$, Theorem 3.3 applies with

- $\mathcal{B}_i := \mathbb{L}^{p_i}(\mu)$ for all p_1, p_2 such that $1 \leq p_2 < p_1 \leq +\infty$,

- and also (useful when $\nu \neq \mu$) with $\mathcal{B}_i = (1 + |f|)^{\gamma_i} L^\infty$ for all γ_1, γ_2 such that $0 \leq \gamma_1 < \gamma_2 < \infty$ and $\mathbb{E}[|f|^{\gamma_2}] < \infty$, where we write again L^∞ for the set of bounded measurable complex valued functions on Ω and \mathcal{B}_i being endowed with the norm $\left\| \frac{\cdot}{(1+|f|)^{\gamma_i}} \right\|_\infty$.

In the general context of Remark 5.3, Item (C) (resp. Item (D)) of Proposition 4.1 applied to this context will provide the following expansions for Π and N^n :

$$\left\| \Pi_t - \sum_{k=0}^{r-1} \frac{\Pi_0^{(k)}}{k!} \cdot t^{\otimes k} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_r)} = \mathcal{O}(|t|^r) \quad \text{and} \quad \left\| N_t^n - \sum_{k=0}^{r-1} \frac{(N^n)_0^{(k)}}{k!} \cdot t^{\otimes k} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_r)} \leq \mathcal{O}(a^n) \mathcal{O}(|t|^r),$$

(resp. with $(\mathcal{O}(t^r), \mathcal{B}_r)$ being replaced by $(o(t^r), \mathcal{B}_{r+1})$) with $\Pi_0^{(k)}, (N^n)_0^{(k)} \in \bigcap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k}^{d^{\otimes k}})$ given by

$$\frac{\Pi_0^{(k)}}{k!} := \frac{1}{2i\pi} \int_{\Gamma_1} \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1: k_1 + \dots + k_\ell = k} (z \text{Id} - P)^{-1} \mathcal{A}_{k_\ell} \otimes \dots \otimes \mathcal{A}_{k_1} dz \quad (15)$$

and

$$\frac{(N^n)_0^{(k)}}{k!} := \frac{1}{2i\pi} \int_{\Gamma_0} z^n \sum_{\ell \geq 1, k_1, \dots, k_\ell \geq 1: k_1 + \dots + k_\ell = k} (z \text{Id} - P)^{-1} \mathcal{A}_{k_\ell} \otimes \dots \otimes \mathcal{A}_{k_1} dz, \quad (16)$$

with $\mathcal{A}_\ell = \frac{P^{(\ell)}}{\ell!} (z \text{Id} - P)^{-1} \in \bigcap_{j=0}^{r-\ell} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+\ell}^{d^{\otimes \ell}})$.

5.2. Square integrable observables.

Theorem 5.5 (Key result for probabilistic limit theorems). *Let $\delta_0 > 0$. Assume Hypothesis 5.1. Let r be a nonnegative integer and $Y_1 \in \mathbb{L}^{r+1}(\mathcal{P}_\mu \otimes \mathbf{P})$. Let $(\mathcal{B}_j, \|\cdot\|_{(j)})$, $j = 0, \dots, r+1$ be a chain of $(r+2)$ Banach spaces such that, for all $j = 1, \dots, r+1$, $\mathcal{B}_{j-1} \hookrightarrow \mathcal{B}_j$, $\mathcal{B}_j \hookrightarrow \mathbb{L}^{\frac{r+1}{j}}(\mu)$. Assume that P_t (for $|t| < \delta_0$) acts continuously on $\mathcal{B}_1, \dots, \mathcal{B}_{r+1}$ and that P_0 acts continuously on \mathcal{B}_0 . Assume moreover that $h, \mathbf{1} \in \mathcal{B}_0$ and that $g: \Omega \rightarrow \mathbb{R}$ is such that $\mathbb{E}_\nu[g \cdot]$ defines a continuous linear form on \mathcal{B}_r . We also assume that (13) with (14) hold true on \mathcal{B}_1 and that*

$$\sup_{j=1, \dots, r+1} \sup_{|t| < \delta_0} \sup_{z \in \Gamma_0 \cup \Gamma_1} \left(\|(z \text{Id} - P_t)^{-1}\|_{\mathcal{B}_j} + \|(z \text{Id} - P)^{-1}\|_{\mathcal{B}_0} \right) < \infty, \quad (17)$$

and that, for all $m = 0, \dots, r$, $P_t - \sum_{k=0}^m \frac{P_0^{(k)}}{k!} \cdot t^{\otimes k}$ is both in $\mathcal{O}(t^m)$ in $\bigcap_{j=0}^{r-m} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+m})$ and in $\mathcal{O}(t^m)$ in $\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_{r+1})$, with

$$P_0^{(k)}(h)(x) = i^k \mathcal{Q}_k(h) := \int_E P \left((if(x, \cdot, \omega))^{\otimes k} h(\cdot) \right) (x) d\mathbf{P}(\omega) \in \bigcap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k}^{d^{\otimes k}}).$$

Then

$$\mathbb{E}_\nu [g P_t^n(h)] = \lambda_t^n \left(\sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\nu [g \Pi_0^{(\ell)}(h)] + \mathcal{O}(t^r) \right) + \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\nu [g (N^n)_0^{(\ell)}(h)] + \mathcal{O}(a^n |t|^r), \quad (18)$$

with $\Pi_0^{(\ell)}$ and $(N^n)_0^{(\ell)}$ given by (15) given by (16).

If moreover $\mathbb{E}_\nu[g \cdot]$ defines a continuous linear form on \mathcal{B}_{r+1} , then

$$\mathbb{E}_\nu [g P_t^n(h)] = \lambda_t^n \left(\sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\nu [g \Pi_0^{(\ell)}(h)] + o(t^r) \right) + \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\nu [g (N^n)_0^{(\ell)}(h)] + \mathcal{O}(a^n o(|t|^r)). \quad (19)$$

Moreover

$$\lambda_t - 1 = \sum_{k=2}^{r+1} \frac{\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}}[(it \cdot Y_1)^k]}{k!} + \frac{\sum_{k=1}^r \sum_{\ell=1}^{r+1-k} \frac{t^{\otimes(k+\ell)}}{k! \ell!} \mathbb{E}_\mu [P_0^{(k)}(\Pi_0^{(\ell)}(\mathbf{1}) - \mathbb{E}_\mu[\Pi_0^{(\ell)}(\mathbf{1})])]}{1 + \sum_{\ell=1}^{r-1} \frac{t^{\otimes \ell}}{\ell!} \mathbb{E}_\mu[\Pi_0^{(\ell)}(\mathbf{1})]} + o(|t|^{r+1}).$$

Note that, in the Markovian context,

$$t^{\otimes(k+\ell)} \mathbb{E}_\mu[P_0^{(k)}(\Pi_0^{(\ell)}(g) - \mathbb{E}_\mu[\Pi_0^{(\ell)}(g)])] = \text{Cov}_{\mathcal{P}_\mu \otimes \mathbf{P}} \left((it.Y_1)^k, \Pi_0^{(\ell)}(\mathbf{1})(X_1).t^{\otimes \ell} \right).$$

In particular, if $f(x, y, z) = f(y) \in \mathbb{R}$, then

$$\mathbb{E}_\mu[P_0^{(k)}(\Pi_0^{(\ell)}(g) - \mathbb{E}_\mu[\Pi_0^{(\ell)}(g)])] = i^k \text{Cov}_\mu \left(f^{\otimes k}, \Pi_0^{(\ell)}(\mathbf{1}) \right).$$

Before proving Theorem 5.5, we state a corollary and apply it to our Knudsen gas model.

Corollary 5.6. *Let $\delta_0 > 0$, $\vartheta \in]0, 1[$, $R > 0$ and r be a nonnegative integer. Assume Hypothesis 5.1 with $Y_1 \in \mathbb{L}^{r+1}(\mathcal{P}_\mu \otimes \mathbf{P})$. Let $(\mathcal{B}_j, \|\cdot\|_{(j)})_{j=0, \dots, r+2}$ be a chain of $(r+3)$ Banach spaces such that:*

- for all $j = 0, \dots, r+1$, $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$, and for all $j = 1, \dots, r+1$, $\mathcal{B}_j \hookrightarrow \mathbb{L}^{\frac{r+1}{j}}(\mu)$,
- for all $j = 1, \dots, r+1$, P is geometrically ergodic on \mathcal{B}_j : $\|P^n - \mathbb{E}_\mu\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta^n)$,
- $\mathbf{1} \in \mathcal{B}_0$, $P \in \mathcal{L}(\mathcal{B}_0)$ and $\sup_{z \in \Gamma_0 \cup \Gamma_1} \|(z \text{Id} - P)^{-1}\|_{\mathcal{B}_0} < \infty$,
- for all $m = 0, \dots, r$, $P_t - \sum_{k=0}^m \frac{P_0^{(k)}}{k!} \cdot t^{\otimes k}$ is both in $\mathcal{O}(t^m)$ in $\bigcap_{j=0}^{r-m} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+m})$ and in $\mathcal{O}(t^m)$ in $\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_{r+1})$, with

$$P_0^{(k)}(h)(x) = i^k \mathcal{Q}_k(h) := \int_E P \left((if(x, \cdot, \omega))^{\otimes k} h(\cdot) \right) (x) d\mathbf{P}(\omega) \in \bigcap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k}^{\otimes k}),$$

- $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_{r+1}, \mathcal{B}_{r+2})$ is continuous at 0, $P \in \mathcal{L}(\mathcal{B}_{r+2})$,
- for all $|t| < \delta_0$ and all $j = 1, \dots, r+1$,

$$\forall f \in \mathcal{B}_j, \forall n \geq 0 \quad \|P_t^n f\|_{(j)} \leq \vartheta^n \|f\|_{(j)} + R^n \|f\|_{(j+1)}. \quad (20)$$

Assume that $h \in \mathcal{B}_0$ and that $g : \Omega \rightarrow \mathbb{R}$ is such that $\mathbb{E}_\nu[g \cdot]$ defines a continuous linear form on \mathcal{B}_r . Then the assumptions of Theorem 5.5 are satisfied (except maybe the fact that $\mathbb{E}_\nu[g \cdot]$ is a linear form on \mathcal{B}_{r+1}) and

$$\mathbb{E}_\nu[g P_t^n(h)] = \lambda_t^n \left(\sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\nu[g \Pi_0^{(\ell)}(h)] + \mathcal{O}(t^r) \right) + \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\nu[g (N^n)_0^{(\ell)}(h)] + \mathcal{O}(a^n |t|^r), \quad (21)$$

with $\Pi_0^{(\ell)}$ and $(N^n)_0^{(\ell)}$ given by (15) given by (16) and

$$\lambda_t - 1 = \sum_{k=2}^{r+1} \frac{\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}}[(it.Y_1)^k]}{k!} + \frac{\sum_{k=1}^r \sum_{\ell=1}^{r+1-k} \frac{t^{\otimes(k+\ell)}}{k!\ell!} \mathbb{E}_\mu[P_0^{(k)}(\Pi_0^{(\ell)}(g) - \mathbb{E}_\mu[\Pi_0^{(\ell)}(g)])]}{1 + \sum_{\ell=1}^{r-1} \frac{t^{\otimes \ell}}{\ell!} \mathbb{E}_\mu[\Pi_0^{(\ell)}(\mathbf{1})]} + o(|t|^{r+1}).$$

Proof. Observe that, given $\vartheta_1 \in]\vartheta, 1[$, up to reduce the value of δ_0 , Theorem 3.3 applied with $\mathcal{B}_j, \mathcal{B}_{j+1}$ for $j = 1, \dots, r+1$ ensures that (13), (14) and the first part of (17). We conclude by Theorem 5.5. \square

Example 5.7 (Knudsen gas). *Let $Q \in]r, r+1[$. Since $Y_1 \in \mathbb{L}^{r+1}(\mathcal{P}_\mu \otimes \mathbf{P})$ and since $|P_t(H)| \leq P(|H|)$, Example 5.2 satisfies the assumptions of Corollary 5.6 for $g \in \mathcal{V}_1$ for (18) and for $g \in \mathcal{V}_{r+1-Q}$ for (19) with $\mathcal{B}_j = \mathcal{V}_j$ for $j = 0, \dots, r$, $\mathcal{B}_{r+1} = \mathcal{V}_Q$ and $\mathcal{B}_{r+2} = \mathcal{V}_{r+1}$, with*

- either, if $\nu = \mu$, $\mathcal{V}_j := \mathbb{L}^{\frac{r+1}{j}}(\mu)$ (with convention $\frac{r+1}{0} = \infty$),
- or, in the general case, $\mathcal{V}_j = (1 + |f|)^j L^\infty$.

This follows from Example 5.4 and from the two following facts

$$\left\| P_t - \sum_{k=0}^m \frac{P((it.f)^k \cdot)}{k!} \right\|_{\mathcal{L}(\mathcal{V}_j, \mathcal{V}_{j+m})} \leq \|P\|_{\mathcal{V}_{j+m}} \left\| e^{it.f} - \sum_{k=0}^m \frac{(it.f)^k}{k!} \right\|_{\mathcal{V}_m} \leq \|P\|_{\mathcal{V}_{j+m}} \left\| \frac{|t f|^m}{m!} \right\|_{\mathcal{V}_m} = \mathcal{O}(|t|^m),$$

and analogously

$$\begin{aligned} \left\| P_t - \sum_{k=0}^{r-j} \frac{P((it.f)^k)}{k!} \right\|_{\mathcal{L}(\mathcal{V}_j, \mathcal{V}_Q)} &\leq \|P\|_{\mathcal{V}_Q} \left\| e^{itf} - \sum_{k=0}^{r-j} \frac{(itf)^k}{k!} \right\|_{\mathcal{V}_{Q-j}} \\ &\leq \|P\|_{\mathcal{V}_Q} \| |tf|^{Q-j} \|_{\mathcal{V}_{Q-j}} = \mathcal{O}(t^{Q-j}). \end{aligned}$$

Proof of Theorem 5.5. Items (C) and (D) of Proposition 4.1 apply, so that

$$\left\| N_t^n - \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot (N^n)_0^{(\ell)} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_r)} = \mathcal{O}(a^n |t|^r), \quad \left\| \Pi_t - \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \Pi_0^{(\ell)} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_r)} = \mathcal{O}(|t|^r) \quad (22)$$

and

$$\left\| N_t^n - \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot (N^n)_0^{(\ell)} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_{r+1})} = \mathcal{O}(a^n o(|t|^r)), \quad \left\| \Pi_t - \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \Pi_0^{(\ell)} \right\|_{\mathcal{L}(\mathcal{B}_0, \mathcal{B}_{r+1})} = o(|t|^r), \quad (23)$$

with $\Pi_0^{(\ell)}, (N^n)_0^{(\ell)} \in \bigcap_{j=0}^{r-\ell} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+\ell}^{d^{\otimes \ell}})$ given by (15) and (16). This ends the proof of (18) and (19).

Due to Fact 3.2 with $v_0 = \mathbf{1}$ and $\varphi_0 = \mathbb{E}_\mu[\cdot]$,

$$\lambda_t - 1 = \mathbb{E}_\mu[(P_t - P)(\mathbf{1})] + \frac{\mathbb{E}_\mu[(P_t - P)(\text{Id} - \Pi_0)(\Pi_t - \Pi_0)(\mathbf{1})]}{\mathbb{E}_\mu[\Pi_t(\mathbf{1})]}.$$

The first expectation of the above right hand side is in $\sum_{k=1}^{r+1} \frac{\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}}[(it.Y_1)^k]}{k!} + o(|t|^{r+1})$ since Y_1 is $(r+1)$ times integrable and $\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}}[Y_1] = 0$. Moreover, since $\mathbf{1} \in \mathcal{B}_0$ and $\mathcal{B}_{r+1} \hookrightarrow \mathbb{L}^1(\mu)$, we also know that

$$\mathbb{E}_\mu[\Pi_t(\mathbf{1})] = \sum_{\ell=0}^r \frac{t^{\otimes \ell}}{\ell!} \cdot \mathbb{E}_\mu[\Pi_0^{(\ell)}(\mathbf{1})] + o(t^r).$$

It remains to study

$$\mathbb{E}_\mu[(P_t - P)(\Pi_t - \Pi_0)(\mathbf{1})].$$

Due to the dominated convergence theorem, since Y_1 is $(r+1)$ times integrable,

$$\left\| e^{it.Y_1} - \sum_{k=0}^r \frac{(it.Y_1)^k}{k!} \right\|_{\mathbb{L}^{\frac{r+1}{r}}(\mathcal{P}_\mu \times \mathbf{P})} = o(t^r).$$

Thus, since $\mathcal{B}_1 \hookrightarrow \mathbb{L}^{r+1}(\mathcal{P}_\mu \otimes \mathbf{P})$, in case (I):

$$\begin{aligned} \mathbb{E}_\mu[(P_t - P)(\Pi_t - \Pi_0)(\mathbf{1})] &= \mathbb{E}_{\mu \otimes \mathbf{P}}[(e^{it.f} - 1)(\Pi_t - \Pi_0)(\mathbf{1})] \\ &= \sum_{k=1}^r \frac{\mathbb{E}_{\mu \otimes \mathbf{P}}[(it.f)^k(\Pi_t - \Pi_0)(\mathbf{1})]}{k!} + o(|t|^r \|(\Pi_t - \Pi_0)(\mathbf{1})\|_{(1)}) . \end{aligned} \quad (24)$$

and, in case (II):

$$\begin{aligned} \mathbb{E}_\mu[(P_t - P)(\Pi_t - \Pi_0)(\mathbf{1})] &= \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}}[(e^{it.Y_1} - 1)(\Pi_t - \Pi_0)(\mathbf{1})(X_1)] \\ &= \sum_{k=1}^r \frac{\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}}[(it.Y_1)^k(\Pi_t - \Pi_0)(\mathbf{1})(X_1)]}{k!} + o(|t|^r \|(\Pi_t - \Pi_0)(\mathbf{1})\|_{(1)}) . \end{aligned} \quad (25)$$

Fix $k = 1, \dots, r$. Note that $Y_1^k \in \mathbb{L}_{\frac{r+1}{k}}(\mathcal{P}_\mu \times \mathbf{P})$. Due to Proposition 4.1, we know that the quantity $\Psi_{r+1-k,t}(\mathbf{1}) := \Pi_t(\mathbf{1}) - \sum_{\ell=0}^{r+1-k} \frac{t^{\otimes \ell}}{\ell!} \cdot \Pi_0^{(\ell)}(\mathbf{1})$ is in $\mathcal{O}(t^{r+1-k})$ in $\mathcal{B}_{r+1-k} \hookrightarrow \mathbb{L}_{\frac{r+1}{r+1-k}}(\mu)$ and is in $o(t^{r+1-k})$ in $\mathcal{B}_{r+1} \hookrightarrow \mathbb{L}^1(\mu)$. We will deduce that

$$\mathbb{E}_{\mu \otimes \mathbf{P}} [f^{\otimes k} \cdot \Psi_{r+1-k,t}(\mathbf{1})] = o(|t|^{r+1-k}), \quad \text{resp. } \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}} [Y_1^{\otimes k} \cdot \Psi_{r+1-k,t}(\mathbf{1})(X_1)] = o(|t|^{r+1-k}). \quad (26)$$

Indeed, since $(\Psi_{r+1-k,t}(\mathbf{1})/|t|^{r+1-k})_t$ is bounded in $\mathbb{L}_{\frac{r+1}{r+1-k}}$, it is contained in a relative compact set for the weak topology. Let h be one of its weak limits as $t \rightarrow 0$. Recall that $(\Psi_{r+1-k,t}(\mathbf{1})/|t|^{r+1-k})_t$ converges to 0 in $L^1(\mu)$. This implies that, for any bounded measurable function H ,

$$\begin{aligned} \mathbb{E}_{\mu \otimes \mathbf{P}} [H \cdot h] &= \lim_{t \rightarrow 0} |t|^{-(r+1-k)} \mathbb{E}_{\mu \otimes \mathbf{P}} [H \cdot \Psi_{r+1-k,t}(\mathbf{1})] = 0, \\ \text{resp. } \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}} [H \cdot h] &= \lim_{t \rightarrow 0} |t|^{-(r+1-k)} \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}} [H^{\otimes k} \cdot \Psi_{r+1-k,t}(\mathbf{1})(X_1)] = 0. \end{aligned}$$

We conclude that $h = 0$, and so that $(\Psi_{r+1-k,t}(\mathbf{1})/|t|^{r+1-k})_t$ converges weakly in $\mathbb{L}_{\frac{r+1}{r+1-k}}$ to 0. This ends the proof of (26) since $Y_1^k \in \mathbb{L}_{\frac{r+1}{k}}(\mathcal{P}_\mu \otimes \mathbf{P})$. We conclude by combining this with (24) and (25) respectively (and using the fact that $\|(\Pi_t - \Pi_0)(\mathbf{1})\|_{(1)} = \mathcal{O}(t)$ since $\mathbf{1} \in \mathcal{B}_0$). \square

We now study the consequences on the smoothness of λ .

Proposition 5.8. *Assume Assumptions of Theorem 5.5 and $r \geq 1$. Then $\lambda_t = 1 - \frac{a}{2} t^{\otimes 2} + o(|t|^2)$, with*

$$a \cdot t^{\otimes 2} = \mathbb{E}_\mu[(t \cdot Y_1)^2] + 2 \sum_{n \geq 0} \mathbb{E}_\mu[(t \cdot \mathcal{Q}_1)(P^n(t \cdot \mathcal{Q}_1(\mathbf{1})))] = \sum_{n \in \mathbb{Z}} \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}}[(t \cdot Y_1)(t \cdot Y_{|n|+1})].$$

Proof of Proposition 5.8. We write Π'_0 for $\Pi_0^{(1)}$. We study the term of order t^2 of $\lambda_t - 1$. Observe that $\mathbb{E}_\mu[P'_0(\mathbf{1})] = 0$ since Y_1 is centered. Due to (14) and to Item (B) of Proposition 4.1,

$$\begin{aligned} \mathbb{E}_\mu [\Pi'_0(\mathbf{1})] &= \frac{1}{2i\pi} \int_{\Gamma_1} \mathbb{E}_\mu [(z \text{Id} - P)^{-1} P'_0(z \text{Id} - P)^{-1}(\mathbf{1})] dz \\ &= \frac{1}{2i\pi} \int_{\Gamma_1} (z - 1)^{-2} \mathbb{E}_\mu [P'_0(\mathbf{1})] dz = 0, \end{aligned}$$

since $\mathbb{E}_\mu[P(h)] = \mathbb{E}_\mu[h]$ and $P(\mathbf{1}) = \mathbf{1}$. Thus $\mathbb{E}_\mu[\Pi'_0(\mathbf{1})] = 0$. Therefore, it follows from Theorem 5.5 that

$$\begin{aligned} \lambda_t - 1 &= -\frac{\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}}[(t \cdot Y_1)^2]}{2} + t^{\otimes 2} \cdot \mathbb{E}_\mu [P'_0(\Pi'_0(\mathbf{1}))] + o(|t|^2) \\ &= -\frac{\mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}}[(t \cdot Y_1)^2]}{2} + i \mathbb{E}_\mu [(t \cdot \mathcal{Q}_1)(t \cdot \Pi'_0(\mathbf{1}))] + o(|t|^2). \end{aligned}$$

It follows from (15) combined with $(z \text{Id} - P)^{-1}(\mathbf{1}) = (z - 1)^{-1} \mathbf{1}$ that

$$\begin{aligned} \Pi'_0(\mathbf{1}) &= \frac{1}{2i\pi} \int_{\Gamma_1} (z \text{Id} - P)^{-1} P'_0(z \text{Id} - P)^{-1}(\mathbf{1}) dz \\ &= \frac{i}{2i\pi} \int_{\Gamma_1} (z \text{Id} - P)^{-1} \mathcal{Q}_1(\mathbf{1})(z - 1)^{-1} dz \\ &= \frac{i}{2i\pi} \int_{\Gamma_1} (z - 1)^{-1} \sum_{n \geq 0} z^{-n-1} P^n(\mathcal{Q}_1(\mathbf{1})) dz \\ &= i \sum_{n \geq 0} P^n(\mathcal{Q}_1(\mathbf{1})), \end{aligned} \quad (27)$$

where we used the fact that $\|z^{-n}P^n(\mathcal{Q}_1(\mathbf{1}))\|_{\mathcal{L}(\mathcal{B}_1)} \leq (1-\delta)^{-n}a^n\|\mathcal{Q}_1(\mathbf{1})\|_{(1)}$ (since $\mathbb{E}_\mu[\mathcal{Q}_1(\mathbf{1})] = 0$ and recalling that $a < 1 - \delta$). This ends the proof of the expression of \mathbf{a} . \square

Example 5.9 (Knudsen gas, normal case). *Due to Example 5.7, Example 5.2 with $r \geq 1$ satisfies the assumptions of Proposition 5.8 and so $\lambda_t - 1 = -\frac{\mathbf{a}}{2}.t^{\otimes 2} + o(|t|^2)$, with*

$$\mathbf{a}.t^{\otimes 2} = \sum_{n \in \mathbb{Z}} \mathbb{E}_{\mathcal{P}_\mu} \left[(t.f(X_1)).(t.f(X_{|n|+1})) \right] = \sum_{n \in \mathbb{Z}} \alpha^{|n|} \mathbb{E}_\mu[(t.f)^2] = \frac{1+\alpha}{1-\alpha} \mathbb{E}_\mu[(t.f)^2].$$

If moreover $r \geq 2$, then the next proposition (Proposition 5.10) also applies, with

$$\mathbf{b} = \mathbb{E}[f^{\otimes 3}] \left(1 + 6 \sum_{n,m \geq 1} \alpha^{n+m} + 3 \sum_{n \geq 1} 2\alpha^n \right) = \frac{\alpha^2 + 4\alpha + 1}{(1-\alpha)^2} \mathbb{E}[f^{\otimes 3}].$$

Let us compute now the term of order 3 in the Taylor expansion of $t \mapsto \lambda_t$.

Proposition 5.10. *Assume the Assumptions of Theorem 5.5 with $r \geq 2$. Then $\lambda_t = 1 - \frac{\mathbf{a}}{2}.t^{\otimes 2} - \frac{i\mathbf{b}}{6}.t^{\otimes 3} + o(|t|^3)$, with \mathbf{a} as in Proposition 5.8 and with*

$$\mathbf{b}.t^{\otimes 3} = \sum_{n,m \geq 0} a_{0,n,n+m} \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}} [(t.Y_1)(t.Y_{1+n})(t.Y_{1+n+m})],$$

with $a_{0,n,n+m} = \#\{(p,q,r) \in \mathbb{Z}^3 : \{p,q,r\} = \{0,n,n+m\}\}$.

Proof. It follows from Theorem 5.5 that

$$\lambda_t - 1 = -\frac{\mathbf{a}.t^{\otimes 2}}{2} - \frac{i}{6} \mathbb{E}_{\mathcal{P}_\mu \otimes \mathbf{P}} [(t.Y_1)^3] + \frac{it^{\otimes 3}}{2} \cdot \mathbb{E}_\mu [\mathcal{Q}_1(\Pi_0''(\mathbf{1}))] - \frac{t^{\otimes 3}}{2} \cdot \mathbb{E}_\mu [\mathcal{Q}_2(\Pi_0'(\mathbf{1}))] + o(|t|^3).$$

Recall that, due to (27), $\Pi_0'(\mathbf{1}) = i \sum_{n \geq 0} P^n(\mathcal{Q}_1(\mathbf{1}))$. It follows moreover from (15) that

$$\begin{aligned} \Pi_0''(\mathbf{1}) &:= \frac{1}{2i\pi} \int_{\Gamma_1} (z \text{Id} - P)^{-1} (P_0'' + 2P_0'(z \text{Id} - P)^{-1} P_0')(z-1)^{-1}(\mathbf{1}) dz \\ &= \frac{1}{2i\pi} \int_{\Gamma_1} (z-1)^{-1} (z \text{Id} - P)^{-1} (-\mathcal{Q}_2 + 2i\mathcal{Q}_1(z \text{Id} - P)^{-1} i\mathcal{Q}_1)(\mathbf{1}) dz \\ &= -\frac{1}{2i\pi} \int_{\Gamma_1} I_1(z) + I_2(z) dz, \end{aligned}$$

with

$$\begin{aligned} I_1(z) &:= (z-1)^{-2} \left(\mathbb{E}_\mu[\mathcal{Q}_2(\mathbf{1})] + 2\mathbb{E}_\mu \left[\mathcal{Q}_1 \left(\sum_{m \geq 0} z^{-m-1} P^m(\mathcal{Q}_1(\mathbf{1})) \right) \right] \right) \\ I_2(z) &:= (z-1)^{-1} \sum_{n \geq 0} z^{-n-1} P^n \left(\mathcal{Q}_2(\mathbf{1}) - \mathbb{E}_\mu[\mathcal{Q}_2(\mathbf{1})] \right. \\ &\quad \left. + 2\mathcal{Q}_1 \left(\sum_{m \geq 0} z^{-m-1} P^m(\mathcal{Q}_1(\mathbf{1})) \right) - 2\mathbb{E}_\mu \left[\mathcal{Q}_1 \left(\sum_{m \geq 0} z^{-m-1} P^m(\mathcal{Q}_1(\mathbf{1})) \right) \right] \right). \end{aligned}$$

It follows that

$$\begin{aligned} -\Pi_0''(\mathbf{1}) &:= -2 \sum_{m \geq 0} (m+1) \mathbb{E}_\mu [\mathcal{Q}_1(P^m(\mathcal{Q}_1(\mathbf{1})))] + \sum_{n \geq 0} (P^n(\mathcal{Q}_2(\mathbf{1})) - \mathbb{E}_\mu[\mathcal{Q}_2(\mathbf{1})]) \\ &\quad + 2 \sum_{n \geq 0} \left(P^n \left(\mathcal{Q}_1 \sum_{m \geq 0} (P^m \mathcal{Q}_1(\mathbf{1})) - \sum_{m \geq 0} \mathbb{E}_\mu [\mathcal{Q}_1(P^m(\mathcal{Q}_1(\mathbf{1})))] \right) \right). \end{aligned}$$

This ends the proof of the lemma. \square

5.3. Non square integrable observable. In view of establishing results of convergence to a stable random variable, we consider now a less smooth situation. If we assume that the distribution of Y_1 is in the standard domain of attraction of an α_0 -stable distribution with $\alpha_0 \in]1, 2[$ (so that $\mathbb{P}(|Y| > s) \sim |s|^{-\alpha_0}$ as $s \rightarrow +\infty$), then we expect that $\lambda_t - 1 \sim -c|t|^{\alpha_0}$. But, unlike in Theorem 5.5, we cannot use an argument of weak convergence to conclude, since we do not have convergence of $\frac{P_t - P_0 - tP'_0}{t^{\alpha_0}}$ and thus we cannot hope the convergence of $\frac{\Pi_t - \Pi_0 - t\Pi'_0}{t^{\alpha_0}}$. The next general statement can be seen as a first step to convergence to stable random variables. We will apply it immediatly on our easy Knudsen gas model.

Proposition 5.11 ($d = 1$). *Let $\delta > 0$. Assume Hypothesis 5.1 with $d = 1$ and that there exist two Banach spaces $\mathcal{B}_1, \mathcal{B}_2$ such that $\mathbf{1} \in \mathcal{B}_1 \hookrightarrow \mathcal{B}_2 \hookrightarrow \mathbb{L}^1(\mu)$, and such that P_t (for any $|t| < \delta$) acts continuously on \mathcal{B}_1 . Assume that P is geometrically ergodic on \mathcal{B}_1 and \mathcal{B}_2 : $\|P^n - \mathbb{E}_\mu[\cdot]\| = \tilde{N}_0^n \|\cdot\|_{\mathcal{L}(\mathcal{B}_j)} = \mathcal{O}(\vartheta^n)$ and that Theorem 3.3 holds true for $(P_t)_t$ with this couple of spaces $(\mathcal{B}_1, \mathcal{B}_2)$. Assume moreover that $P_t - P_0$ is in $\mathcal{O}(t)$ in $\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ and in $o(1)$ in $\mathcal{L}(\mathcal{B}_2, \mathbb{L}^1(\mu))$. Let $\vartheta_1 > \vartheta$, $\gamma \in]0, 1[$. We assume that $\|\tilde{N}_0^n\|_{\mathcal{L}(\mathcal{B}_2)} + \sup_{|t| < \delta} \|\tilde{N}_t^n\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta_1^n)$, that $\|\tilde{N}_t - \tilde{N}_0\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} = \mathcal{O}(t)$, where $\tilde{N}_t := P_t - \Psi_t(\cdot)\mathbf{1}$ with $\Psi_0(\mathbf{1}) = \gamma + o(t)$ in \mathbb{C} . Then*

$$\begin{aligned} \lambda_t - 1 &= \gamma \mathbb{E}_\mu \left[(P_t - P)(\text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) \right] + o(|t|^2) \\ &= \gamma \sum_{n \geq 0} \mathbb{E}_\mu \left[(P_t - P) \tilde{N}_t^n(\mathbf{1}) \right] + o(|t|^2). \end{aligned}$$

Example 5.12 (Knudsen gas, stable case). *Let $\alpha_0 \in]1, 2]$ and $p \in]1, \alpha_0[$ be such that $\|f\|_{\mathbb{L}^p(\mu)} < \infty$. Consider the Knudsen gas introduced in Example 5.2 with $r = 0$ and $d = 1$. Assume that the characteristic function φ_f of f with respect to μ satisfies*

$$\varphi_f(t) - 1 = \mathbb{E}_\mu[e^{itf}] - 1 = -|t|^{\alpha_0}(1 - i\beta \text{sgn}(t))L_0(|t|^{-1}), \quad (28)$$

with $|\beta| < \tan(\alpha_0\pi/2)$, $c > 0$ and with L_0 slowly varying at infinity. Then

$$\begin{aligned} \lambda_t - 1 &\sim (\varphi_f(t) - 1)(1 - \alpha)^2 \sum_{n \geq 0} \alpha^n (n + 1)^{\alpha_0} = (\varphi_f(t) - 1)(1 - \alpha) \mathbb{E}[\mathcal{G}^{\alpha_0}] \\ &\sim (1 - \alpha) \left(\mathbb{E} \left[e^{it\mathcal{G}f} \right] - 1 \right) \sim \left(\mathbb{E} \left[e^{i(1-\alpha)^{\frac{1}{\alpha_0}} t \mathcal{G}f} \right] - 1 \right). \end{aligned} \quad (29)$$

as $t \rightarrow 0$, where \mathcal{G} is a geometric random variable with parameter $1 - \alpha$ (i.e. $\mathbb{P}(\mathcal{G} = n) = (1 - \alpha)\alpha^n$ for all positive integer n) independent of f (up to enlarge the probability space).

The toy Knudsen gas model considered here can also be studied by induction, considering the successive "renewal times" T_j at which X_{T_j} is chosen independently of $X_{T_{j-1}}$. The increments $T_j - T_{j-1}$ of these times are independent with the same distribution as \mathcal{G} and the sums $Y_{T_{j-1}} + \dots + Y_{(T_j)-1} = (T_j - T_{j-1})Y_{T_{j-1}}$ between these times are independent with the same distribution as $\mathcal{G}f$. Furthermore the number N_n of j 's such that $t_j \leq n$ is a sum of n independent Bernoulli random variables with parameter $(1 - \alpha)$, and so $N_n \sim n(1 - \alpha)$ almost surely. Thus it will not be surprising that $\sum_{k=1}^n Y_k$ will "behave asymptotically" as $\sum_{j=1}^{N_n} Z_j$ and as $\sum_{j=1}^{(1-\alpha)n} Z_j$, where Z_j are independent random variables with the same distribution as $\mathcal{G}f$. This is coherent with (29) (via the arguments of Section 6.1).

Proof of Example 5.12. We apply Proposition 5.11 with $\mathcal{B}_1 = \mathbb{L}^\infty(\mu)$ and $\mathcal{B}_2 = \mathbb{L}^p(\mu)$. We have already seen in Example 5.4 that P is geometrically ergodic on \mathcal{B}_1 and on \mathcal{B}_2 and that P_t acts continuously on both these spaces. Moreover

$$\|P_t - P\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} \leq \|e^{itf} - 1\|_{\mathbb{L}^p(\mu)} = \|tf\|_{\mathbb{L}^p(\mu)} = \mathcal{O}(|t|),$$

since $f \in \mathbb{L}^p(\mu)$ and, setting $q \in]1, +\infty[$ such that $\frac{1}{p} + \frac{1}{q} = 1$ and using the dominated convergence theorem, we also obtain that

$$\|P_t - P\|_{\mathcal{L}(\mathcal{B}_2, \mathbb{L}^1(\mu))} \leq \|e^{itf} - 1\|_{\mathbb{L}^q(\mu)} = o(1),$$

Here we take $\tilde{N}_t(h) := \alpha e^{itf} h$ (taking $\gamma = 1 - \alpha$ and $\Psi_t = (1 - \alpha)\mathbb{E}_\mu[P_t(\cdot)]$). Observe that $\tilde{N}_t - \tilde{N}_0$ is in $\mathcal{O}(t)$ in $\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ and in $o(1)$ in $\mathcal{L}(\mathcal{B}_2, \mathbb{L}^1(\mu))$. It follows from Proposition 5.11 that

$$\begin{aligned} \lambda_t - 1 &= (1 - \alpha)\mathbb{E}_\mu \left[(P_t - P)(\text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) \right] + o(|t|^2) \\ &= (1 - \alpha)\mathbb{E}_\mu \left[(e^{itf} - 1)(\text{Id} - \alpha e^{itf})^{-1}(\mathbf{1}) \right] + o(|t|^2) \\ &= (1 - \alpha)\mathbb{E}_\mu \left[\frac{e^{itf} - 1}{1 - \alpha e^{itf}} \right] + o(|t|^2). \end{aligned} \quad (30)$$

But

$$\begin{aligned} \mathbb{E}_\mu \left[\frac{e^{itf} - 1}{1 - \alpha e^{itf}} \right] &= \sum_{n \geq 0} \mathbb{E}_\mu \left[(e^{itf} - 1)(\alpha e^{itf})^n \right] \\ &= \sum_{n \geq 0} \alpha^n \mathbb{E}_\mu \left[e^{it(n+1)f} - e^{itnf} \right]. \end{aligned} \quad (31)$$

Recall that Karamata proved that there exists $u_0 > 0$ and two functions c, ε_0 such that $\lim_{s \rightarrow +\infty} c(s) > 0$ and $\lim_{s \rightarrow +\infty} \varepsilon_0(s) = 0$ such that $L_0(u) = c(u)e^{\int_{u_0}^u \frac{\varepsilon_0(s)}{s} ds}$. Set $N_t = \lfloor u_0/\sqrt{t} \rfloor$ and let us control (31) as follows

$$\begin{aligned} \sum_{n \geq 0} \alpha^n \mathbb{E} \left[e^{i(n+1)tf} - e^{intf} \right] &= \mathbb{E} \left[\frac{(e^{itf} - 1)(\alpha e^{itf})^{N_t+1}}{1 - \alpha e^{itf}} \right] + \sum_{n=0}^{N_t} \alpha^n \mathbb{E} \left[e^{i(n+1)tf} - e^{intf} \right] \\ &= \mathcal{O} \left(\frac{\alpha^{u_0/\sqrt{t}}}{1 - \alpha} \|e^{itf} - 1\|_{\mathbb{L}^1} \right) + \varphi_f(t) - 1 - \sum_{n=1}^{N_t} \alpha^n ((\varphi_f((n+1)t) - 1) - (\varphi_f(nt) - 1)) \\ &\sim (\varphi_f(t) - 1) \left[1 + \sum_{n=1}^{N_t} \alpha^n ((n+1)^{\alpha_0} - |n|^{\alpha_0}) \right] \\ &\sim (\varphi_f(t) - 1) \left[1 + (\alpha^{N_t+1}(N_t+1)^{\alpha_0} - \alpha) - \sum_{n=1}^{N_t} (n+1)^{\alpha_0} (\alpha^{n+1} - \alpha^n) \right] \\ &\sim (\varphi_f(t) - 1)(1 - \alpha) \left[\sum_{n \geq 0} (n+1)^{\alpha_0} \alpha^n \right], \end{aligned} \quad (32)$$

due to the dominated convergence theorem since, for all $n \geq 1$, $L_0(|nt|^{-1}) \sim L_0(|t|^{-1})$ as $t \rightarrow 0$ and, for all $n = 1, \dots, N_t$,

$$\frac{L_0(|nt|^{-1})}{L_0(|t|^{-1})} = \frac{c(|nt|^{-1})}{c(|t|^{-1})} e^{-\int_{|nt|^{-1}}^{|t|^{-1}} \frac{\varepsilon_0(s)}{s} ds} = \mathcal{O} \left(n^{\sup_{|s| > |\sqrt{t}u_0|^{-1}} |\varepsilon_0(s)|} \right).$$

Combining (30), (31), and (32), we obtain the first equivalent. The others follow since $\mathbb{E}[\mathcal{G}^{\alpha_0}] = (1 - \alpha) \sum_{n \geq 0} \alpha^n (n+1)^{\alpha_0}$, since $\mathbb{E}[e^{it\mathcal{G}f}] = \mathbb{E}[\varphi_f(\mathcal{G}t)]$ and since $\varphi_f(ut) - 1 \sim u^{\alpha_0}(\varphi_f(t) - 1)$ (using the Lebesgue dominated convergence theorem and again the Karamata representation of slowly varying functions). \square

Proof of Proposition 5.11. Observe that $\tilde{N}_0(\mathbf{1}) = (1 - \gamma)\mathbf{1}$. It follows from Theorem 3.3 applied to $(P_t)_t$ that, for $\vartheta_2 \in]\vartheta_1, 1[$, up to reduce δ if necessary, Formulas (13) and (14) hold true (with $\delta_0 = \delta$) on \mathcal{B}_1 and that

$$\Pi_t = \frac{1}{2i\pi} \int_{\Gamma_1} (z \text{Id} - P_t)^{-1} dz \quad (33)$$

is continuous from $] -\delta, \delta[$ to $\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ and assuming moreover that $\Gamma_1 \subset \{z \in \mathbb{C} : |z| > \vartheta_2\}$. Furthermore, $\sup_{|t| < \delta} \|\tilde{N}_t^n\|_{\mathcal{L}(\mathcal{B}_1)} = \mathcal{O}(\vartheta_1^n)$ ensures that $(z \text{Id} - \tilde{N}_t)^{-1} = \sum_{n \geq 0} z^{-n-1} \tilde{N}_t^n \in \mathcal{L}(\mathcal{B}_1)$ is uniformly bounded in $|z| \geq \vartheta_2$ (and thus in $z \in \Gamma_1$) and that $(z \text{Id} - \tilde{N}_0)^{-1} = \sum_{n \geq 0} z^{-n-1} \tilde{N}_0^n \in \mathcal{L}(\mathcal{B}_2)$ uniformly in $|z| \geq \vartheta_2$. Due to Fact 3.2, since $\mathcal{B}_2 \hookrightarrow \mathbb{L}^1(\mu)$,

$$\lambda_t - 1 = \frac{\mathbb{E}_\mu[(P_t - P)(\Pi_t(\mathbf{1}))]}{\mathbb{E}_\mu[\Pi_t(\mathbf{1})]} = \mathbb{E}_\mu[(P_t - P)\Pi_t(\mathbf{1})] (1 + \mathcal{O}(t)) , \quad \text{as } t \rightarrow 0 .$$

Let us observe that

$$(z \text{Id} - P_t)^{-1}(h) = (z \text{Id} - \tilde{N}_t)^{-1} \left(h + \frac{\Psi_t(z \text{Id} - \tilde{N}_t)^{-1}(h)}{1 - \Psi_t((z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}))} \mathbf{1} \right) . \quad (34)$$

Indeed $g := (z \text{Id} - P_t)^{-1}(h)$ satisfies

$$(z \text{Id} - \tilde{N}_t)(g) = (z \text{Id} - P_t)(g) + \Psi_t(g)w_t = h + c_t(h)\mathbf{1} ,$$

with

$$c_t(h) = \Psi_t(g) = \Psi_t((z \text{Id} - \tilde{N}_t)^{-1}(h + c_t(h)\mathbf{1})) ,$$

from which we infer that $c_t(h) = \frac{\Psi_t((z \text{Id} - \tilde{N}_t)^{-1}(h))}{1 - \Psi_t((z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}))}$ and so (34). Thus, applying (34) to $h = \mathbf{1}$, we obtain that

$$(z \text{Id} - P_t)^{-1}(\mathbf{1}) = (z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1})b_t(z), \quad \text{with } b_t(z) := \frac{1}{1 - \Psi_t((z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}))} . \quad (35)$$

Note that we can recover the fact that $(z \text{Id} - P)^{-1}(\mathbf{1}) = (z - 1)^{-1}$. Moreover, we will prove that

$$\sup_{z \in \Gamma_1} |b_t(z) - b_0(z)| = \mathcal{O}(|t|) . \quad (36)$$

To this end, we first notice that $\Psi_t = \mathbb{E}_\mu[(P_t - \tilde{N}_t)(\cdot)]$ and so

$$\sup_{z \in \Gamma_1} (\Psi_t - \Psi_0)(z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) = \mathcal{O}(t) \quad (37)$$

since $\sup_{z \in \Gamma_1} \|(z \text{Id} - \tilde{N}_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_1)} < \infty$, since $\mathcal{B}_2 \hookrightarrow \mathbb{L}^1(\mu)$ and since

$$\|P_t - P\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} + \|\tilde{N}_t - \tilde{N}_0\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} = \mathcal{O}(t) . \quad (38)$$

Second, it follows from (38) and

$$\sup_{|t| < \delta} \sup_{z \in \Gamma_1} \left(\|(z \text{Id} - \tilde{N}_0)^{-1}\|_{\mathcal{L}(\mathcal{B}_2)} + \|(z \text{Id} - \tilde{N}_t)^{-1}\|_{\mathcal{L}(\mathcal{B}_1)} \right) < \infty ,$$

that

$$\left((z \text{Id} - \tilde{N}_t)^{-1} - (z \text{Id} - \tilde{N}_0)^{-1} \right) = (z \text{Id} - \tilde{N}_0)^{-1}(\tilde{N}_t - \tilde{N}_0)(z \text{Id} - \tilde{N}_t)^{-1} = \mathcal{O}(t) \quad (39)$$

in $\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)$ uniformly in $z \in \Gamma_1$ and so that

$$\psi_0 \left(\left((z \text{Id} - \tilde{N}_t)^{-1} - (z \text{Id} - \tilde{N}_0)^{-1} \right) (\mathbf{1}) \right) = \mathcal{O}(t) \quad (40)$$

since $\mathcal{B}_2 \hookrightarrow \mathbb{L}^1(\mu)$. Combining (37) and (40), we have proved (36).

It follows also from (39), that $(z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) = (z - 1 + \gamma)^{-1} \mathbf{1} + \mathcal{O}(t)$ in \mathcal{B}_2 . Since moreover $P_t - P = o(1)$ in $\mathcal{L}(\mathcal{B}_2, \mathbb{L}^1(\mu))$ and $\mathbb{E}_\mu[(P_t - P)(\mathbf{1})] = \varphi_f(t) - 1 = o(t)$ (since Y_1 is centered),

$$\sup_{z \in \Gamma_1} \left| \mathbb{E}_\mu \left[(P_t - P)(z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) \right] \right| = \sup_{z \in \Gamma_1} \left| \mathbb{E}_\mu \left[(P_t - P)(z - 1 + \gamma)^{-1} \mathbf{1} \right] \right| + o(t) = o(t).$$

Combining this with (33), (35) and (36), we infer that

$$\begin{aligned} \mathbb{E}_\mu [(P_t - P)\Pi_t(\mathbf{1})] &= \frac{1}{2i\pi} \int_{\Gamma_1} \mathbb{E}_\mu \left[(P_t - P)(z \text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) \right] b_0(z) dz \\ &\quad + \mathcal{O} \left(o(t) \sup_{z \in \Gamma_1} |b_t(z) - b_0(z)| \right) \\ &= \gamma \mathbb{E}_\mu \left[(P_t - P)(\text{Id} - \tilde{N}_t)^{-1}(\mathbf{1}) \right] + o(t^2), \quad \text{as } t \rightarrow 0, \end{aligned}$$

where we used the fact that $b_0(z) = \frac{z-1+\gamma}{z-1}$. \square

6. PROBABILISTIC LIMIT THEOREMS

Let $\delta_0 > 0$. Let $(S_n)_{n \geq 1}$ be a sequence of \mathbb{X} -valued random variables with $\mathbb{X} = \mathbb{R}^d$ or \mathbb{Z}^d defined on a probability space $(\mathcal{M}, \mathbb{P})$ such that

$$\forall n \in \mathbb{N}, \quad \forall t \in \mathbb{R}^d, |t| < \delta_0, \quad \mathbb{E}[e^{it \cdot S_n}] = \lambda_t^n \Phi_t + M_{t,n}. \quad (41)$$

We set $\mathbb{X}^* := \mathbb{R}^d$ if $\mathbb{X} = \mathbb{R}^d$ and $\mathbb{X}^* := [-\pi, \pi]^d$ if $\mathbb{X} = \mathbb{Z}^d$.

Remark 6.1. In Remark 5.3, we have seen general situations in which $\mathbb{E}[e^{it \cdot S_n}] = \mathbb{E}_\nu[P_t^n(h_0)]$ (see (11) and (12)) for some h_0 and some family of operators $(P_t)_t$ such that (41) holds true with λ and $\Phi_t = \mathbb{E}_\nu[\Pi_t(h_0)]$ continuous in t and with $\sup_{|t| < \delta_0} |M_{t,n}| = \mathcal{O} \left(\sup_{|t| < \delta_0} \|N_t^n(h_0)\|_{\mathcal{B}_2} \right)$ decaying exponentially fast in n . Recall moreover that further Taylor expansions have been studied in Theorem 5.5.

The goal of this section is to establish probabilistic limit theorems for $(S_n)_{n \geq 1}$. More precisely, we will study situations in which $(S_n)_{n \geq 1}$ satisfies the same kind of limit theorems as sums of independent identically distributed random variables with characteristic function behaving at 0 as $t \mapsto \lambda_t$.

6.1. Central and local limit theorems. Let \mathcal{W} be a \mathbb{R}^d -valued random variable.

Theorem 6.2 (Central Limit Theorem (CLT)). *Let $(A_n)_{n \geq 1}$ be a sequence of (normalizing) $d \times d$ matrices converging to 0. Assume that $\lim_{t \rightarrow 0} \Phi_t = 1$, that $\lim_{n \rightarrow +\infty} \sup_{|t| < \delta_0} |M_{t,n}| = 0$ and that $\lim_{n \rightarrow +\infty} \lambda_{A_n^* t}^n = \mathbb{E}[e^{it \cdot \mathcal{W}}]$ for all $t \in \mathbb{R}^d$ (writing A_n^* for the transpose matrix of A_n). Then $(A_n S_n)_{n \geq 1}$ converges in distribution to \mathcal{W} .*

Proof of Theorem 6.2. We prove the convergence of the characteristic functions. We fix $t \in \mathbb{R}^d$ and write

$$\mathbb{E} \left[e^{it \cdot (A_n S_n)} \right] - \mathbb{E}[e^{it \cdot \mathcal{W}}] = \lambda_{A_n^* t}^n \Phi_{A_n^* t} + M_{A_n^* t, n} - \mathbb{E}[e^{it \cdot \mathcal{W}}] = o(1),$$

since $\lim_{n \rightarrow +\infty} \Phi_{A_n^* t} = 1$, $\lim_{n \rightarrow +\infty} \lambda_{A_n^* t}^n = \mathbb{E}[e^{it \cdot \mathcal{W}}]$ and since

$$\lim_{n \rightarrow +\infty} M_{A_n^* t, n} \leq \lim_{n \rightarrow +\infty} \sup_{|u| < \delta_0} |M_{u, n}| = 0.$$

Thus $(A_n S_n)_{n \geq 1}$ converges in distribution to \mathcal{W} . \square

The condition $\lim_{n \rightarrow +\infty} \lambda_{A_n^* t}^n = \mathbb{E}[e^{it \cdot \mathcal{W}}]$ means that λ_t behaves at 0 as the characteristic function of a distribution belonging to the domain of attraction of the stable distribution of \mathcal{W} . In particular, if $\lambda_t - 1 \sim -a|\Sigma t|_\alpha^\alpha$ as $t \rightarrow 0$ with $|s|_\alpha^\alpha = \sum_{i=1}^d |s_i|^\alpha$ and with B an invertible matrix, then, setting $A_n = n^{-\frac{1}{\alpha}} \text{Id}$ and considering \mathcal{W} such that $\mathbb{E}[e^{it \cdot \mathcal{W}}] = e^{-a|Bt|_\alpha^\alpha}$, the following estimate holds true for any $t \in \mathbb{R}^d$

$$\begin{aligned} \left| \lambda_{A_n^* t}^n - \mathbb{E}[e^{it \cdot \mathcal{W}}] \right| &= \left| \lambda_{t/n^{\frac{1}{\alpha}}}^n - e^{-an|Bt/n^{\frac{1}{\alpha}}|_\alpha^\alpha} \right| \\ &\leq n \left| \lambda_{t/n^{\frac{1}{\alpha}}} - e^{-a|Bt/n^{\frac{1}{\alpha}}|_\alpha^\alpha} \right| \leq n o \left(|Bt/n^{\frac{1}{\alpha}}|_\alpha^\alpha \right) = o(1), \end{aligned}$$

as n goes to infinity.

Remark 6.3. We consider again the context of Remark 5.3. We recall that the fact that (41) holds with $\lim_{t \rightarrow 0} \Phi_t = 1$ and $\lim_{n \rightarrow +\infty} \sup_{|t| < \delta_0} |M_{t,n}| = 0$ follows from Remark 6.1. Furthermore,

- (a) Under the assumptions of Theorem 5.5 for $r \geq 1$ (and so of Proposition 5.8), Theorem 6.2 holds true with $A_n = \text{Id} / \sqrt{n}$ and \mathcal{W} a centered Gaussian random variable with variance

$$\Sigma := \sum_{n \in \mathbb{Z}} \text{Cov}_{\mathcal{P}_\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}} (Y_1, Y_{|n|+1}).$$

Indeed, Proposition 5.8 ensures⁸ that $\lambda_t - 1 = \frac{1}{2} |\Sigma^{\frac{1}{2}} t|_2^2 + o(|t|^2)$, whence we infer that, for every fixed $t \in \mathbb{R}^d$,

$$\begin{aligned} \left| \lambda_{t/\sqrt{n}}^n - \mathbb{E}[e^{it \cdot \mathcal{W}}] \right| &= \left| \lambda_{t/n^{\frac{1}{\alpha}}}^n - e^{-\frac{1}{2} n |\Sigma^{\frac{1}{2}} t/\sqrt{n}|_2^2} \right| \\ &\leq n \left| \lambda_{t/\sqrt{n}} - e^{-\frac{1}{2} |\Sigma^{\frac{1}{2}} t/\sqrt{n}|_2^2} \right| \leq n o \left(|t/\sqrt{n}|_2^2 \right) = o(1). \end{aligned}$$

- (b) If Proposition 5.11 allows to prove that $\lambda_t - 1 \sim \mathbb{E}[e^{it \cdot W_1}] - 1$, where $(W_n)_{n \geq 1}$ is a sequence of i.i.d. random variables such that $(A_n \sum_{k=1}^n W_k)_n$ converges in distribution to \mathcal{W} , then Theorem 6.2 applies.

This assumption ensures (see e.g. [27, Theorem 2.6.5]) the existence of $\alpha_0 \in [1, 2]$ and of $|\beta| < \tan(\alpha_0 \pi / 2)$ such that

$$\mathbb{E}[e^{it \cdot \mathcal{W}}] = e^{-c_0 |t|^{\alpha_0} (1 - i\beta \text{sgn}(t))}, \quad \lambda_t - 1 \sim c_0 |t|^{\alpha_0} (1 - i\beta \text{sgn}(t)) L_0(|t|^{-1})$$

as $t \rightarrow 0$, with L_1 slowly varying at infinity, $A_n \rightarrow 0$ and $\lim_{n \rightarrow +\infty} n |A_n|^{\alpha_0} L_0(A_n^{-1}) = 1$. Indeed, proceeding as previously, we observe that, for every $t \in \mathbb{R}$,

$$\begin{aligned} \left| \lambda_{A_n t}^n - \mathbb{E}[e^{it \cdot \mathcal{W}}] \right| &= \left| \lambda_{A_n t}^n - \mathbb{E}[e^{in^{-\frac{1}{\alpha_0}} t \cdot \mathcal{W}}] \right| \\ &\leq n \left| \lambda_{A_n t} - \mathbb{E}[e^{in^{-\frac{1}{\alpha_0}} t \cdot \mathcal{W}}] \right| = n |t|^{\alpha_0} |c_0 (1 - i\beta \text{sgn}(t))| \left| A_n^{\alpha_0} L_0(|t|/A_n) - n^{-1} \right| \\ &\leq |t|^{\alpha_0} |c_0 (1 - i\beta \text{sgn}(t))| |n A_n^{\alpha_0} L_0(|t|/A_n) - 1| \\ &\leq |t|^{\alpha_0} |c_0 (1 - i\beta \text{sgn}(t))| \left| \frac{L_0(|t|/A_n)}{L_0(A_n^{-1})} (1 + o(1)) - 1 \right| = o(1), \end{aligned}$$

since L_0 is slowly varying.

⁸setting $\Sigma^{\frac{1}{2}}$ for the nonnegative symmetric matrix the square of which is Σ .

Example 6.4 (Knudsen gas, convergence to gaussian or stable distributions). *Consider the simple Knudsen gas model introduced in Example 5.2. The continuity of $t \mapsto \Phi_t = \mathbb{E}_\mu[\Pi_t(\mathbf{1})]$ as well as the fact that $\sup_{|t| < \delta} \mathbb{E}_\mu[N_t^n(\mathbf{1})]$ decays exponentially fast as $n \rightarrow +\infty$ have been proved in Example 5.4 thanks to Theorem 3.3 and Remark 6.1.*

Furthermore

- *If $r \geq 1$, then $(\sum_{k=0}^{n-1} f(X_k)/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable \mathcal{W} with variance matrix $\Sigma = \frac{1+\alpha}{1-\alpha} \mathbb{E}[f^{\otimes 2}]$. This follows from Theorem 6.2 combined with the first item of Remark 6.3, the desired second order expansion being proved in Example 5.9.*
- *Consider now the situation of Example 5.12, that is $r = 0$, $d = 1$ and there exists $\alpha_0 \in]1, 2]$ and a function L_0 slowly varying at infinity such that the characteristic function φ_f of f with respect to μ satisfies*

$$\varphi_f(t) - 1 \sim |t|^{\alpha_0} (1 - i\beta \operatorname{sgn}(t)) L_0(|t|^{-1}),$$

as $t \rightarrow 0$, for some $|\beta| < \tan(\alpha_0 \pi/2)$. Then $(A_n \sum_{k=0}^{n-1} f(X_k))_{n \geq 1}$ converges in distribution to the stable random variable \mathcal{W} with characteristic function given by $\mathbb{E}[e^{it\mathcal{W}}] = e^{-c|t|^{\alpha_0}(1-i\beta \operatorname{sgn}(t))}$, for $c = (1-\alpha) \sum_{n \geq 0} \alpha^n (n+1)^{\alpha_0}$ and where $A_n \rightarrow 0$ is so that $\lim_{n \rightarrow +\infty} n|A_n|^{\alpha_0} L_0(A_n^{-1}) = 1$.

Indeed, due to Example 5.12, $\lambda_t - 1 \sim c|t|^{\alpha_0} L_0(|t|^{-1})(1 - i\beta \operatorname{sgn}(t))$ and Theorem 6.2 applies with the second item of Remark 6.3.

Theorem 6.5 (Local Limit Theorem). *Assume Assumptions of Theorem 6.2 with A_n invertible, $\sup_{|t| < \delta_0} |\Phi(t)| < \infty$, $\sup_{|t| < \delta_0} |M_{t,n}| = o(\det A_n)$ and $|\lambda_{A_n^* t}^n| \leq g(t)$ if $|A_n^* t| < \delta_0$, with g integrable on \mathbb{R}^d . Assume moreover that \mathcal{W} has density $h_{\mathcal{W}}$ and integrable characteristic function, that f takes its values in \mathbb{Z}^d and that $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]| = o(\det A_n)$. Then*

$$\sup_{k \in \mathbb{Z}} |\mathbb{P}(S_n = k) - \det(A_n) h_{\mathcal{W}}(A_n k)| = o(\det(A_n)).$$

Proof of Theorem 6.5. Observe that

$$\begin{aligned} \mathbb{P}(S_n = k) &= \mathbb{E}[\mathbf{1}_{\{S_n - k = 0\}}] = \mathbb{E}\left[\frac{1}{(2\pi)^d} \int_{[-\pi, \pi]^d} e^{it \cdot (S_n - k)} dt\right] \\ &= \frac{1}{(2\pi)^d} \int_{[-\pi, \pi]^d} e^{-itk} \mathbb{E}[e^{it \cdot S_n}] dt \\ &= \frac{1}{(2\pi)^d} \int_{B(0, \delta_0)} e^{-itk} \lambda_t^n \Phi_t dt + o(\det A_n), \end{aligned}$$

where we used the Fubini theorem for integrable functions and the fact that $\sup_{|t| < \delta_0} |M_{t,n}| = o(\det A_n)$ and that $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]| = o(\det A_n)$. Now, making the change of variable $t = A_n^* s$, we obtain

$$\mathbb{P}(S_n = k) = \frac{\det(A_n)}{(2\pi)^d} \int_{(A_n^*)^{-1} B(0, \delta_0)} e^{-iA_n^* s \cdot k} \lambda_{A_n^* s}^n \Phi_{A_n^* s} ds + o(\det A_n),$$

and

$$\begin{aligned} \sup_{k \in \mathbb{Z}} \left| \int_{(A_n^*)^{-1} B(0, \delta_0)} e^{-iA_n^* s \cdot k} \lambda_{A_n^* s}^n \Phi_{A_n^* s} ds - \int_{\mathbb{R}^d} e^{-is \cdot A_n k} \mathbb{E}[e^{is \cdot \mathcal{W}}] ds \right| \\ \leq \int_{\mathbb{R}^d} |\mathbf{1}_{(A_n^*)^{-1} B(0, \delta_0)} \lambda_{A_n^* s}^n \Phi_{A_n^* s} - \mathbb{E}[e^{is \cdot \mathcal{W}}]| ds = o(1), \end{aligned}$$

due to the dominated convergence theorem since g and $s \mapsto \mathbb{E}[e^{-is\mathcal{W}}]$ are integrable. We end the proof by using $h_{\mathcal{W}}(u) = \frac{1}{(2\pi)^d} \int_{\mathbb{R}^d} e^{-is \cdot u} \mathbb{E}[e^{is \cdot \mathcal{W}}] ds$. \square

Remark 6.6. Consider the context of Remark 6.3 with $(\det(A_n))_n$ subexponential in n . We have already seen in Remark 6.1 that $\sup_{|t| < \delta_0} |\Phi(t)| < \infty$ and $\sup_{|t| < \delta_0} |M_{t,n}| = o(\det A_n)$ for some $\delta_0 > 0$.

The integrability of $t \mapsto \sup_n |\lambda_{A_n^* t}^n \mathbf{1}_{|A_n^* t| < \delta_0}|$ follows in practice from the control of $|\lambda_t - 1|$ (e.g. in case (a) of Remark 6.3, if \mathbf{a} is invertible, $|\lambda_t| \leq e^{-\frac{\mathbf{a} \cdot t \otimes 2}{4}}$ as soon as $|t|$ is small enough). Finally, the condition $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]| = o(\det A_n)$ is usually a consequence of the fact that, for $t \notin 2\pi\mathbb{Z}$, $\rho_{\text{ess}}(P_t) < 1$, $\rho(P_t) \leq 1$ and that P_t admits no eigenvalue of modulus 1, which implies that $\|P_t^{n_t}\| < 1$ (for some n_t) which, combined with a continuity argument of $t \mapsto \|P_t\|$ on the compact $[-\pi, \pi]^d$, leads to the existence of a positive integer n' such that $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} \|P_t^{n'}\| < 1$ and implies the exponential decay of $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]|$ as $n \rightarrow +\infty$.

To complete this remark, let us indicate that, under the assumptions of Theorem 3.3 (Keller and Liverani theorem), it has been proved in [26, Propositions 5.3 and 5.4] that the nonlattice property together with an additional reasonable condition imply the exponential decay of $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]|$ as n goes to infinity.

Example 6.7 (Knudsen gas, LLT). Theorem 6.5 applies in the situation considered in Example 6.4, provided f takes its values in \mathbb{Z}^d but is not supported by a sublattice of \mathbb{Z}^d .

Proof of Example 6.7. The fact that the assumptions of Theorem 6.2 hold has already been proved in Example 6.4 with the use of Remarks 6.1 and 6.3. First, we observe that, either $A_n = \text{Id}/\sqrt{n}$ (if $r \geq 1$) or $A_n = n^{\frac{1}{\alpha_0}} \tilde{L}(n)$ with \tilde{L} slowly varying at infinity (in the situation of Example 5.12). The fact that $\lim_{n \rightarrow +\infty} \lambda_{A_n^* t}^n = \mathbb{E}[e^{it \cdot \mathcal{W}}]$ with either \mathcal{W} as in Example 6.4 has been proved in this example using Theorem 6.2 and Remark 6.3. Let us prove the domination of $|\lambda_{A_n^* t}^n|$ using the expansion of λ_t around $t = 0$.

If $r \geq 1$. The nonlattice assumption ensures the invertibility of the matrix Σ . Example 5.9 ensures that $\lambda_t - 1 \sim -\frac{\Sigma}{2} \cdot t^{\otimes 2} + o(|t|^2)$. Thus, there exists $\delta > 0$, such that, for all $|t| < \delta$, $|\lambda_t| \leq g(t) := e^{-\frac{1}{4} \Sigma \cdot t^{\otimes 2}}$, and so $|\lambda_{t/\sqrt{n}}^n| \leq (g(t/\sqrt{n}))^n = g(t)$ and g is integrable.

Assume now that $\lambda_t - 1 \sim c|t|^{\alpha_0} L_0(|t|^{-1})(1 - i\beta \text{sgn}(t))$ with $\alpha_0 \in [1, 2]$ and with L_0 slowly varying at infinity. Then

$$\forall t \in \mathbb{R}, \quad \lambda_{A_n t}^n = e^{-nc|A_n t|^{\alpha_0}(1-i\beta \text{sgn}(t))L_0(|A_n t|^{-1})} = e^{-c|t|^{\alpha_0}(1-i\beta \text{sgn}(t))\frac{L_0(|A_n t|^{-1})}{L_0(A_n^{-1})}}.$$

Due to Karamata's representation theorem, there exists $u_0 > 0$ and two functions c, ε_0 such that $\lim_{s \rightarrow +\infty} c(s) > 0$, $\lim_{s \rightarrow +\infty} \varepsilon_0(s) = 0$ and such that $L_0(u) = c(u)e^{\int_{u_0}^u \frac{\varepsilon_0(s)}{s} ds}$. Let δ_0 be such that λ_u is well defined for all u satisfying $|u| < \delta_0$. Thus, if $|A_n t| < \delta_0 < 1$ and $|A_n| < \delta_0$,

$$\frac{L_0(|A_n t|^{-1})}{L_0(A_n^{-1})} = \frac{c(|A_n t|^{-1})}{c(|A_n|^{-1})} e^{\int_{A_n^{-1}}^{|A_n t|^{-1}} \frac{\varepsilon_0(s)}{s} ds} \geq \frac{\inf_{|s| > \delta^{-1}} c(s)}{\sup_{|s| > \delta^{-1}} c(s)} \min(|t|, |t|^{-1})^{\inf_{|s| > \delta^{-1}} |\varepsilon_0(s)|}.$$

We can choose $\delta_0 < 1$, such that this quantity is larger than $\frac{1}{2} \min(|t|^{\frac{\alpha_0}{2}}, |t|^{-\frac{\alpha_0}{2}})$ and so, for every t and n such that $|A_n t| < \delta_0$ and $|A_n| < \delta_0$,

$$|\lambda_{A_n t}^n| \leq g(t) := e^{-\frac{c}{2} \min(|t|^{\frac{\alpha_0}{2}}, |t|^{-\frac{\alpha_0}{2}})}.$$

It remains to prove that $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]| = o(\det A_n)$. To this end, we follow the strategy explained in Remark 6.6. Let $t \in \mathbb{R} \setminus 2\pi\mathbb{Z}$ and $p \in]1, +\infty[$. Since $\|P_t\|_{\mathbb{L}^p(\mu)} \leq \alpha\|h\| + (1-\alpha)\|h\|_{\mathbb{L}^1(\mu)}$, it follows by standard arguments (see e.g. [28, 18]) that the essential spectral radius of P_t is strictly smaller than 1. Consider now $h \in \mathbb{L}^p(\mu)$ and $\lambda = e^{id} \in \mathbb{C}$ with $d \in \mathbb{R}$ and $\lambda h = P_t(h) = \alpha e^{it \cdot f} h + (1-\alpha)\mathbb{E}_\mu[e^{it \cdot f} h]$. Then $\mathbb{E}_\mu[|h|] \leq \mathbb{E}_\mu[|\alpha e^{it \cdot f} h|] + (1-\alpha)\mathbb{E}_\mu[|e^{it \cdot f} h|] \leq \mathbb{E}_\mu[|h|]$ and we conclude that $\lambda h = e^{it \cdot f} h = \mathbb{E}_\mu[e^{it \cdot f} h]$ μ -a.e., thus $h = e^{-id} e^{it \cdot f} h$ is constant. So either $h = 0$ or $e^{id} = e^{it \cdot f}$. But $e^{id} = e^{it \cdot f}$ would mean that $t \cdot f \in d + 2\pi\mathbb{Z}$, which would contradict the fact that f is not contained in a sublattice of \mathbb{Z}^d . Thus $h = 0$.

Since P_t has an essential spectral radius strictly smaller than 1 and does not admit any eigenvalue of modulus 1, we conclude that its spectral radius is strictly smaller than 1. So there exists $n_t \geq 1$ such that $\|\mathbb{E}_{\mathcal{P}_\mu}[e^{it S_{n_t}}]\| \leq \|P_t^{n_t} \mathbf{1}\|_{\mathbb{L}^p(\mu)} < 1$. But $u \mapsto \mathbb{E}_{\mathcal{P}_\mu}[e^{iu S_{n_t}}] = \mathbb{E}_\mu[P_u^{n_t}(\mathbf{1})]$ is continuous at t (since $u \mapsto P_u \in \mathcal{L}(\mathbb{L}^{a'}(\mu), \mathbb{L}^{b'}(\mu))$ is continuous, for all $a' > b' \geq 1$, using also the fact that $\sup_u \|P_u\|_{\mathcal{L}(\mathbb{L}^a(\mu))} \leq 1$). We conclude by compactity. \square

6.2. Edgeworth expansions ($d = 1$). We assume $d = 1$ throughout this section. We recall now some general Edgeworth expansions results coming from [10]. We first introduce some assumptions.

Assumption $(\alpha')[\hat{r}]$ (Smoothness): Assume λ , Φ and M enjoys the following Taylor expansions

$$\lambda_t = 1 - \frac{\sigma^2 t^2}{2} + \sum_{k=3}^{\hat{r}+2} \alpha_k t^k + o(t^{\hat{r}+2}), \quad \Phi_t = \sum_{k=0}^{\hat{r}} B_k t^k + \mathcal{O}(|t|^{\hat{r}+1}),$$

$$\exists \delta > 0, \quad \forall t \in [-\delta, \delta], \quad \left| \lambda_t^{-n} M_{t,n} - \sum_{k=0}^{\hat{r}} C_{k,n} t^k \right| \leq K_n \lambda_t^{-\frac{n}{2}} |t|^{\hat{r}+1},$$

with $\sigma^2 > 0$, $\sup_{k,n} (|C_{k,n}|) = \mathcal{O}(n^{-p})$ and $K_n = \mathcal{O}(n^p)$ for all $p > 0$.

Remark 6.8. In the Markovian context of Section 5, Assumption $(\alpha')[\hat{r}]$ will follow from Theorem 5.5 and Corollary 5.6 with $r := \hat{r} + 1$ (up to assume the positivity of σ^2 , the expression of which is given by \mathbf{a} of Proposition 5.8). Indeed the Taylor expansion of $M_{t,n}$ combined with the one of λ_t coming from Theorem 5.5 leads to a Taylor expansion of order $r + 1$ of $\lambda_t^{-n} M_{t,n}$ with coefficients $C_{k,n} = \mathcal{O}(n^k a^n)$ and with error term in $\mathcal{O}(\lambda_t^{-n} n^{r+1} |t|^{\hat{r}+1} a^n)$.

Assumption (β') (Non-arithmeticity): For any compact K of $\mathbb{X}^* \setminus \{0\}$,

$$\forall p > 0, \quad \sup_{s \in K} |\mathbb{E}[e^{is S_n}]| = \mathcal{O}(n^{-p}),$$

As already mentioned in Remark 6.6, in the Markovian context of Section 5, Assumption (β') will follow from the fact that $\rho_{ess}(P_t) < 1$ (which can be established by using (4) if $\mathcal{B}_1 \hookrightarrow \mathcal{B}_2$ is compact, due to a Theorem by Hennion [18]), combined with the fact that $\rho(P_t) \leq 1$ and that P_t admits no eigenvalue of modulus 1 (for t in the compact K).

Assumption $(\gamma')[\alpha'_1, \alpha_1]$: Either $\mathbb{X} = \mathbb{Z}$, or there exist $\hat{\delta}$ such that, for all p ,

$$\forall |s| > K, \quad |\mathbb{E}[e^{is S_n}]| = \mathcal{O}\left(n^{-p} + |s|^{1+\alpha'_1} e^{-n\alpha_1 \hat{\delta} |s|^{-\alpha'_1}}\right).$$

Assumption $(\delta')[r']$: $\mathbb{X} = \mathbb{R}$ and for any $B > 0$, there exists $K > 0$ such that

$$\int_{K < |s| < Bn^{\frac{r'-1}{2}}} \frac{|\mathbb{E}_\mu(e^{is S_n})|}{|s|} ds = o(n^{-r'/2}).$$

Remark 6.9. Note that, both Conditions $(\gamma')(\alpha'_1, 1)$ and $(\delta')[r']$ hold true provided there exist positive $s_0, n_0, \widehat{C} > 0$, such that $r' < 2\alpha'_1{}^{-1} + 1$ and such that

$$\forall |s| > s_0, \forall n \geq n_0, \quad \left| \mathbb{E}[e^{isS_n}] \right| < e^{-n\widehat{C}|s|^{-\alpha'_1}}. \quad (42)$$

In the context of Section 5 with $\mathbb{E}[e^{itS_n}] = \mathbb{E}_\nu[gP_t^n(h)]$, (42) holds true if there exist Banach spaces $\mathcal{B}_1 \hookrightarrow \mathcal{B}_2$ containing h and the duals of which contains $\mathbb{E}_\nu[g \cdot]$ such that

$$\forall n \geq n_0, \forall |s| > s_1, \quad \|P_s^n\|_{\mathcal{L}(\mathcal{B}_1, \mathcal{B}_2)} \leq Ce^{-Cn|s|^{-\alpha'_1}}.$$

(see [10, Lemma 4.7]). Note that this holds true if $\|P_s^{n_0}\|_{\mathcal{B}_1} < 1 - \frac{C}{|s|^{\alpha'_1}}$. This condition generalizes the α -Diophantine property of $\text{supp } X$: $\mathbb{E}[e^{isY_1}] < 1 - \frac{\widehat{C}}{|s|^{\alpha'_1}}$ of the i.i.d. case (see [8]).

Example 6.10 (Knudsen gas). Consider again Example 5.2. Recall that $P(h - \mathbb{E}_\mu[h]) = \alpha(h - \mathbb{E}_\mu[h])$ with $\alpha \in]0, 1[$ and that $Y_n = f(X_n)$, with $f : \Omega \rightarrow \mathbb{R}$ centered. Assume f admits moments of order $(\widehat{r} + 2)$ with respect to $\nu = \mu$.

The fact that Assumption $(\alpha')[\widehat{r}]$ holds true follows from Theorem 5.5 and Example 5.7.

Assumption (β') will hold true if, for all $t \in \mathbb{X}^* \setminus \{0\}$, $|\mathbb{E}_\mu[e^{itf}]| < 1$. Indeed, as seen in Example 6.7, it is enough to prove that P_t (for $t \in \mathbb{X}^*$) admits no eigenvalue λ of modulus 1. If it was the case, there would exist $h \in \mathcal{B}_1 \setminus \{0\}$ such that $\lambda h = e^{itf}h = \mathbb{E}_\mu[e^{itf}h]$ μ -almost surely, contradicting $|\mathbb{E}_\mu[e^{itf}]| < 1$.

Finally, when $\mathbb{X} = \mathbb{R}$, Assumptions $(\gamma')[\alpha'_1, 1]$ and $(\delta')[\lceil 1 + 2\alpha'_1{}^{-1} \rceil - 1]$ hold true as soon as $\mathbb{E}_\mu[e^{itf}] < 1 - \frac{\widehat{C}_0}{|t|^{\alpha'_1}}$. Indeed

$$P_t^2(h)(x) = \alpha^2 e^{i2tf(x)}h(x) + \alpha(1-\alpha)e^{itf}\mathbb{E}_\mu[e^{itf}h] + \alpha(1-\alpha)\mathbb{E}_\mu[e^{i2tf}h] + (1-\alpha)^2\mathbb{E}_\mu[e^{itf}]\mathbb{E}_\mu[e^{itf}h],$$

and so, for all $p \in [1, +\infty]$, it follows that

$$\begin{aligned} \|P_t^2(h)\|_{\mathbb{L}^p(\mu)} &\leq \left(1 - (1-\alpha)^2\right) \|h\|_{\mathbb{L}^p(\mu)} + (1-\alpha)^2 \left(1 - \frac{\widehat{C}_0}{|t|^{\alpha'_1}}\right) |\mathbb{E}_\mu[e^{itf}h]| \\ &\leq \left(1 - (1-\alpha)^2 \frac{\widehat{C}_0}{|t|^{\alpha'_1}}\right) \|h\|_{\mathbb{L}^p(\mu)}. \end{aligned}$$

Set $\widehat{g}(s) := \int_{\mathbb{X}} e^{-isx}g(x) d\lambda(x)$ for $s \in \mathbb{X}^*$, where λ is the Lebesgue measure if $\mathbb{X} = \mathbb{R}$ and where λ is the counting measure if $\mathbb{X} = \mathbb{Z}$. If $\mathbb{X} = \mathbb{R}$, we say $g \in \mathfrak{F}_k^m$ if $g : \mathbb{R} \rightarrow \mathbb{R}$ is continuous, λ -integrable and if $\widehat{g} : \mathbb{X}^* \rightarrow \mathbb{C}$ is k times continuously differentiable with

$$C_k^m(g) := C^m(g) + C_k(g) < \infty,$$

$$\text{with } C^m(g) := \sup_{s \in \mathbb{X}^*} \frac{|\widehat{g}(s)|}{\min(1, |s|^{-m})} \quad \text{and} \quad C_k(g) := \|\widehat{g}^{(k)}\|_\infty.$$

If $\mathbb{X} = \mathbb{Z}$, $\mathfrak{F}_k^m = \mathfrak{F}_k^0$ is the set of functions $g : \mathbb{Z} \rightarrow \mathbb{C}$ satisfying the following summability condition

$$\sum_{n \in \mathbb{Z}} |n|^k |g(n)| < \infty.$$

Note that $C_k(g) \leq \max_{0 \leq j \leq k} \int_{\mathbb{X}} |x|^j |g(x)| d\lambda(x)$. When $\mathbb{X} = \mathbb{R}$, $C^m(g) \leq \max_{0 \leq j \leq m} \|g^{(j)}\|_{L^1(\mathbb{R})}$. Under our assumptions, we set \mathfrak{N} for the distribution function of a centered Gaussian random variable with variance σ^2 and \mathfrak{n} for the corresponding probability density function (that is \mathfrak{n} is the derivative of \mathfrak{N}). Let us recall now the general results of [10].

Theorem 6.11. [10, Theorem 1.1] *Let \hat{r} be a nonnegative integer, $\alpha'_1 \geq 0$, $\alpha_1 > 0$ and $q > \alpha'_1(1 + \frac{\hat{r}}{2\alpha_1})$. Assume $(\alpha')[\hat{r}]$, (β') and $(\gamma')[\alpha'_1, \alpha_1]$ hold. Then there exist polynomials R_j such that, for all $g \in \mathfrak{F}_0^{q+2}$,*

$$\mathbb{E}[g(S_n)] = \sum_{j=0}^{\hat{r}} \frac{1}{n^{(j-1)/2}} \int_{\mathbb{X}} (R_j \cdot \mathbf{n})(x/\sqrt{n}) g(x) d\lambda(x) + C^{q+2}(g) \cdot o(n^{-\hat{r}/2}).$$

Theorem 6.12. [10, Theorem 1.2] *Let \hat{r} be a nonnegative integer, $\alpha'_1 \geq 0$, $\alpha_1 > 0$. Let $q > \alpha'_1(1 + \frac{\hat{r}+1}{2\alpha_1})$. Assume $(\alpha')[\hat{r}]$, (β') and $(\gamma')[\alpha'_1, \alpha_1]$ hold. Then there exist polynomials Q_j such that, for all $g \in \mathfrak{F}_{\hat{r}+1}^{q+2}$,*

$$\sqrt{n}\mathbb{E}[g(S_n)] = \sum_{j=0}^{\lfloor \hat{r}/2 \rfloor} \frac{1}{n^j} \int_{\mathbb{X}} g(x) Q_j(x) d\lambda(x) + C_{\hat{r}+1}^{q+2}(g) \cdot o(n^{-\hat{r}/2}).$$

Theorem 6.13. [10, Theorem 1.7] *Let \hat{r} be a positive integer and $r' \in [1, \hat{r}]$ be a real number. Assume $(\alpha')[\hat{r}]$, (β') and $(\delta')[r']$ hold. Then there exist polynomials P_k such that*

$$\sup_{x \in \mathbb{R}} \left| \mathbb{P}\left(\frac{S_n}{\sqrt{n}} \leq x\right) - \mathfrak{N}(x) - \mathbf{n}(x) \sum_{k=1}^{\lfloor r' \rfloor} \frac{P_k(x)}{n^{k/2}} \right| = o(n^{-r'/2}).$$

Corollary 6.14. [10, Corollary 1.8] *Assume $(\alpha')[1]$ and (β') hold with $\mathbb{X} = \mathbb{R}$. Then*

$$\sup_{x \in \mathbb{R}} \left| \mathbb{P}\left(\frac{S_n}{\sqrt{n}} \leq x\right) - \mathfrak{N}(x) - \frac{P_1(x)}{n^{1/2}} \mathbf{n}(x) \right| = o(n^{-1/2}).$$

Corollary 6.15. [10, Corollary 1.9] *Assume $(\alpha')[2]$, (β') and $(\delta')[r_0]$ hold for some real number $r_0 \in (1, 2)$. Then*

$$\sup_{x \in \mathbb{R}} \left| \mathbb{P}\left(\frac{S_n}{\sqrt{n}} \leq x\right) - \mathfrak{N}(x) - \frac{P_1(x)}{\sqrt{n}} \mathbf{n}(x) \right| = o(n^{-r_0/2}).$$

Example 6.16 (Knudsen gas, Edgeworth expansions in CLT, LLT). *Consider Example 5.2. Recall that Assumptions (α') , (β') , (γ') and (δ') have been checked in Example 6.10. Assume f admits moments of order $(\hat{r}+2)$ with respect to $\nu = \mu$ and that $|\mathbb{E}_\mu[e^{itf}]| < 1$ for all $t \in \mathbb{X}^* \setminus \{0\}$.*

- (Expansions of order $\hat{r} - 1$ and \hat{r} in the LLT) *Assume either $\mathbb{X} = \mathbb{Z}$ or $|\mathbb{E}_\mu[e^{itf}]| < 1 - \frac{\hat{C}_0}{|t|^{\alpha'_1}}$ for some $\alpha'_1 > 0$, then the conclusions of Theorems 6.11 and 6.12 hold true with respectively $q > \alpha'_1(1 + \frac{\hat{r}}{2})$ and $q > \alpha'_1(1 + \frac{\hat{r}+1}{2})$.*
- (Edgeworth expansion of order \hat{r}) *Assume $\mathbb{X} = \mathbb{R}$ and $|\mathbb{E}_\mu[e^{itf}]| < 1 - \frac{\hat{C}_0}{|t|^{\alpha'_1}}$ for some $\alpha'_1 < \left(\frac{r'-1}{2}\right)^{-1}$. Then the conclusion of Theorem 6.13 holds true.*
- (First order Edgeworth expansion) *Assume $\mathbb{X} = \mathbb{R}$ and $\hat{r} = 1$. Then the conclusion of Theorem 6.14 holds true.*
- (Edgeworth expansion of order $r_0 \in]1, 2[$) *Assume $\mathbb{X} = \mathbb{R}$, $\hat{r} = 2$ and that $|\mathbb{E}_\mu[e^{itf}]| < 1 - \frac{\hat{C}_0}{|t|^{\alpha'_1}}$ for some $\alpha'_1 < (\frac{r_0-1}{2})^{-1}$. Then the conclusion of Corollary 6.15 holds true.*

7. LIMIT THEOREMS FOR MARKOV RANDOM WALKS

We focus again in this section on the context of Markov random walks, that is the general Markovian setting of Section 5. Recall that $(X_n)_{n \geq 0}$ is a Markov chain with states space Ω and

with invariant distribution μ and initial distribution ν and that $(Z_k)_{k \geq 1}$ is a sequence of independent identically distributed random variables with common distribution \mathbf{P} and independent of the Markov chain $(X_n)_{n \geq 0}$. Recall that we are interested in the behaviour of $S_n := \sum_{k=1}^n Y_k$, with $Y_k = f(X_{k-1}, X_k, Z_k)$.

In a first subsection, we establish probabilistic limit theorems in the general context as a direct consequence of the results of Section 6. In the three following subsections, we apply our approach for classical families of Markov chains: the ρ -mixing Markov chains, the V -geometrically ergodic Markov chains and Lipschitz iterative model. More precisely, we prove Theorem 1.1 in these three last subsections.

7.1. General results. We set \mathcal{P}_ν for the Markov distribution with transition operator P and initial probability measure ν . We assume that $((x_k)_{k \geq 0}, \omega) \mapsto f(x_0, x_1, \omega)$ is $\mathcal{P}_\mu \otimes \mathbf{P}$ -centered. We set

$$P_t(h)(x) = \mathbb{E} \left[e^{it \cdot Y_1} h(X_1) | X_0 \right].$$

We establish probabilistic limit theorems under the assumptions of Theorem 5.5 that we recall in the following statement for reader's convenience.

Theorem 7.1. *Let $\delta_0 > 0$. Let r be a positive integer and $Y_1 \in \mathbb{L}^{r+1}(\mathcal{P}_\mu \otimes \mathbf{P})$. Let $(\mathcal{B}_j, \|\cdot\|_{(j)})$, $j = 0, \dots, r+1$ be a chain of $(r+2)$ Banach spaces such that $\mathbf{1} \in \mathcal{B}_0$ and that for all $j = 1, \dots, r+1$, $\mathcal{B}_{j-1} \hookrightarrow \mathcal{B}_j$, $\mathcal{B}_j \hookrightarrow \mathbb{L}^{\frac{r+1}{j}}(\mu)$. Assume that P_t (for $|t| < \delta_0$) acts continuously on $\mathcal{B}_1, \dots, \mathcal{B}_{r+1}$ and that P_0 acts continuously on \mathcal{B}_0 and that, for all $m = 0, \dots, r$, $P_t - \sum_{k=0}^m \frac{P_0^{(k)}}{k!} \cdot t^{\otimes k}$ is both $\mathcal{O}(t^m)$ in $\bigcap_{j=0}^{r-m} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+m})$ and $o(t^m)$ in $\mathcal{L}(\mathcal{B}_{r-m}, \mathcal{B}_{r+1})$, with*

$$P_0^{(k)}(h)(x) = i^k \mathcal{Q}_k(h)(x) := \int_E P \left((if(x, \cdot, \omega))^{\otimes k} h(\cdot) \right) (x) d\mathbf{P}(\omega) \in \bigcap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k}^{\otimes k}).$$

Finally we assume that, in $\mathcal{L}(\mathcal{B}_1)$,

$$\exists \vartheta_1 \in]0, 1[, \forall n \in \mathbb{N}^*, \forall |t| < \delta_0, \quad P_t^n = \lambda_t^n \Pi_t + N_t^n, \quad \text{with } \sup_{|t| < \delta_0} \|N_t^n\|_{\mathcal{B}_1} = \mathcal{O}(\vartheta_1^n), \quad (43)$$

with $\Pi_0 = \mathbb{E}_\mu[\cdot] \mathbf{1}$, $\lambda_0 = 1$ and

$$\Pi_t := \frac{1}{2i\pi} \int_{\Gamma_1} (z \text{Id} - P_t)^{-1} dz, \quad N_t^n := \frac{1}{2i\pi} \int_{\Gamma_0} z^n (z \text{Id} - P_t)^{-1} dz, \quad (44)$$

with Γ_1 the oriented circle $\mathcal{C}(1, \delta)$ and Γ_0 the oriented circle $\mathcal{C}(0, a)$, with $\vartheta_1 < a < a + \delta < 1$ and that⁹

$$\sup_{j=1, \dots, r+1} \sup_{|t| < \delta_0} \sup_{z \in \Gamma_0 \cup \Gamma_1} \left(\|(z \text{Id} - P_t)^{-1}\|_{\mathcal{B}_j} + \|(z \text{Id} - P)^{-1}\|_{\mathcal{B}_0} \right) < \infty. \quad (45)$$

Assume either $\mathbb{E}_\nu \in \mathcal{B}_1^*$ or, more generally, that there exists some Banach space $\tilde{\mathcal{B}}_0$ (that can be intermediate between \mathcal{B}_0 and \mathcal{B}_1) such that and $\mathbb{E}_\nu \in \tilde{\mathcal{B}}_0^*$ and such that $t \mapsto \Pi_t \in \mathcal{L}(\mathcal{B}_0, \tilde{\mathcal{B}}_0)$ is continuous at 0 (using e.g. Theorem 3.3), then

- (i) Theorem 6.2 (CLT) applies and we conclude that $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable with variance matrix $\sum_{n \in \mathbb{Z}} \text{Cov}_{\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}} (Y_1, Y_{|n|+1})$.
- (ii) Theorem 6.5 (LLT) applies if Y_1 is \mathbb{Z}^d valued and if the non-arithmetic condition (β') is satisfied.

If moreover $\mathbb{E}_\nu \in \mathcal{B}_r^*$, then Assumption $(\alpha')[\hat{r} := r-1]$ is satisfied. In particular:

⁹recall that, in practice, (43), (44) and the first part of (45) follow from the Keller-Liverani theorem (Theorem 3.3) applied with the Banach spaces $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$ up to consider an additional Banach space \mathcal{B}_{r+2} .

- (iii) ($d = 1$) Theorems 6.11 and 6.12 (Expansions of order $\hat{r} - 1 = r - 2$ and $\hat{r} = r - 1$ in the LLT) apply if Y_1 is \mathbb{Z} -valued and if the non-arithmetic condition (β') is satisfied.
- (iii') ($d = 1$) Theorems 6.11 and 6.12 (Expansions of order $\hat{r} - 1 = r - 2$ and $\hat{r} = r - 1$ in the LLT) apply if conditions (β') (non-arithmeticity) and $(\gamma')[\alpha'_1, \alpha_1]$ (Diophantine-type condition) are satisfied.
- (iv) ($d = 1$) Corollary 6.14 (First order Edgeworth expansion) holds true if $r \geq 2$ and if conditions (β') (non-arithmeticity) is satisfied on $\mathbb{X} = \mathbb{R}$.
- (v) ($d = 1$) Theorem 6.13 (Edgeworth expansion of order $\hat{r} = r - 1$) applies if $r \geq 2$ and if the non-arithmetic condition (β') and Condition $(\delta')[r']$ hold true.

7.2. ρ -mixing Markov chains. We consider here the case of Markov chains that are ρ -mixing, i.e. i.e. when the transfer operator P is geometrically ergodic on $\mathbb{L}^2(\mu)$, that is satisfies

$$\exists C > 0, \exists \vartheta \in]0, 1[, \quad \forall g \in \mathbb{L}^2(\mu), \quad \|P^n(g) - \mathbb{E}_\mu[g] \mathbf{1}\|_{\mathbb{L}^2(\mu)} \leq C \vartheta^n \|g\|_{\mathbb{L}^2(\mu)}.$$

Recall that this also implies the geometric ergodicity on each $\mathbb{L}^p(\mu)$ for $p \in]1, +\infty[$ (see [38]). Let us observe that the study of Markov random walks driven by a ρ -mixing Markov chain $(X_k)_k$ can be simplified in an additive function of a ρ -mixing Markov chain.¹⁰

Proposition 7.2. *If $(X_k)_{k \geq 0}$ is ρ -mixing, then the Markov chain $(\tilde{X}_k := (X_{k-1}, X_k, Z_k))_{k \geq 0}$ with invariant probability measure $\tilde{\mu}$ the distribution of (X_0, X_1, Z_0) with respect to \mathcal{P}_μ , which is given by*

$$\tilde{\mu}(A \times B \times C) = \mathbb{E}_\mu [P(\mathbf{1}_B) \mathbf{1}_A] \mathbf{P}(C)$$

is also ρ -mixing, with same rate.

Proof. Let \tilde{P} be the transfer operator of \tilde{X} . Then, for all $n \geq 1$,

$$\begin{aligned} \tilde{P}^n(G)(x, y, z) &= \mathbb{E} [G(X_{n-1}, X_n, Z_n) | X_0 = y] \\ &= \int_E (P^{n-1}(H(\cdot, z)))(y) d\mathbf{P}(z), \end{aligned}$$

where $H(x, z) = \mathbb{E} [G(X_0, X_1, z) | X_0 = x]$, and so

$$\begin{aligned} \|\tilde{P}^n(G) - \mathbb{E}_{\tilde{\mu}}[G]\|_{L^2(\tilde{\mu})}^2 &= \int_\Omega \left| \int_E \left(P^{n-1}(H(\cdot, z))(y) - \mathbb{E}_\mu[H(\cdot, z)] \right) d\mathbf{P}(z) \right|^2 d\mu(y) \\ &\leq \int_\Omega \int_E \left| P^{n-1}(H(\cdot, z))(y) - \mathbb{E}_\mu[H(\cdot, z)] \right|^2 d\mathbf{P}(z) d\mu(y) \\ &\leq \int_E \left\| P^{n-1}(H(\cdot, z))(y) - \mathbb{E}_\mu[H(\cdot, z)] \right\|_{\mathbb{L}^2(\mu)}^2 d\mathbf{P}(z) \\ &\leq \int_E C^2 \vartheta^{2(n-1)} \|H(\cdot, z)\|_{\mathbb{L}^2(\mu)}^2 d\mathbf{P}(z) \\ &= \int_E C^2 \vartheta^{2(n-1)} \|\mathbb{E} [G(X_0, X_1, z) | X_0]\|_{\mathbb{L}^2(\mu)}^2 d\mathbf{P}(z) \leq C^2 \vartheta^{2(n-1)} \|G\|_{\mathbb{L}^2(\tilde{\mu})}^2. \end{aligned}$$

□

Thus, without any loss of generality, from now on, in this subsection, we replace $f(x, y, z)$ by $f(y)$ (up to replace the Markov chain X by the Markov chain \tilde{X}). Note that this replacement changes the notion of non-lattice, which can be corrected by using [24].

¹⁰Note that this change induces a change in the definition of non-lattice.

Theorem 7.3. Assume P is geometrically ergodic on \mathbb{L}^2 . Assume the initial measure ν is the stationary measure μ and $Y_k = f(X_k)$, with $f : \Omega \rightarrow \mathbb{R}^d$. Let r be a positive integer. Assume f is μ -centered and in $\mathbb{L}^{r+1}(\mu)$. Then the assumptions of Theorem 7.1 hold true with $\mathcal{B}_0 = \mathbb{C}.1$ endowed with the infinite norm and with $\mathcal{B}_j = \mathbb{L}^{\frac{r+1}{j}}(\mu)$ for all $j \in \{1, \dots, r\}$ and $\mathcal{B}_{r+1} = \mathbb{L}^{\frac{2r+1}{2r}}(\mu)$ (note that $1 < \frac{2r+1}{2r} < \frac{r+1}{r}$). In particular

- (i) Theorem 6.2 (CLT) applies and we conclude that $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable with variance matrix $\sum_{n \in \mathbb{Z}} \text{Cov}_{\mathcal{P}_\mu}(Y_1, Y_{|n|+1})$.
- (ii) Theorem 6.5 (LLT) applies if Y_1 is \mathbb{Z}^d -valued and non-lattice in \mathbb{Z}^d .¹¹
- (iii) ($d = 1$) Theorems 6.11 and 6.12 (Expansions of order $\hat{r} - 1 = r - 2$ and $\hat{r} = r - 1$ in the LLT) apply if Y_1 is \mathbb{Z} -valued and non-lattice in \mathbb{Z} .
- (iii') ($d = 1$) Theorems 6.11 and 6.12 (Expansions of order $\hat{r} - 1 = r - 2$ and $\hat{r} = r - 1$ in the LLT) apply if Assumption $(\gamma')[\alpha'_1, \alpha_1]$ holds true and if Y_1 is non-lattice in \mathbb{R} .¹²
- (iv) ($d = 1$) Corollary 6.14 (First order Edgeworth expansion) holds true if $r \geq 2$ and if Y_1 is non-lattice in \mathbb{R} .
- (v) ($d = 1$) Theorem 6.13 (Edgeworth expansion of order $r' \in [1, \hat{r} = r - 1]$) applies if $r \geq 2$, if Y_1 is non-lattice in \mathbb{R} and if Assumption $(\delta')[r']$ holds true.

Proof. Since $P = \text{Id}$ on \mathcal{B}_0 , 1 is the single spectral value of $P|_{\mathcal{B}_0}$ and $\|(z \text{Id} - P)^{-1}\|_{\mathcal{L}(\mathcal{B}_0)} = |z - 1|^{-1}$. We know that for all $p \in]1, +\infty[$, $\|P^n - \mathbb{E}_\mu[\cdot] \mathbf{1}\|_{\mathcal{L}(\mathbb{L}^p)}$ decreases exponentially fast. This implies in particular that P is quasi-compact on \mathcal{B}_1 and on \mathcal{B}_2 with a single dominating eigenvalue 1, which is simple and also that

$$\|P_t^n(g)\|_{\mathbb{L}^p(\mu)} \leq \|P^n(|g|)\|_{\mathbb{L}^p(\mu)} \leq \|P^n - \mathbb{E}_\mu[\cdot] \mathbf{1}\|_{\mathcal{L}(\mathbb{L}^p)} \|g\|_{\mathbb{L}^p(\mu)} + \|g\|_{\mathbb{L}^1(\mu)},$$

implying the uniform Doeblin-Fortet inequality (4) for $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$. Moreover, for all $j = 0, \dots, r$ and $m \in \{1, \dots, r - j\}$,

$$\begin{aligned} \left\| P_t(g) - \sum_{k=0}^m \frac{P((if)^k g)}{k!} t^k \right\|_{\mathcal{B}_{j+m}} &\leq \left\| P \left(\left(e^{itf} - \sum_{k=0}^m \frac{(if)^k}{k!} t^k \right) g \right) \right\|_{\mathcal{B}_{j+m}} \\ &\leq \left\| e^{itf} - \sum_{k=0}^m \frac{(if)^k}{k!} t^k \right\|_{\mathcal{B}_m} \|g\|_{\mathcal{B}_j} \leq o(t^m) \|g\|_{\mathcal{B}_j}, \end{aligned}$$

due to the dominated convergence theorem. Note that this last inequality implies in particular the continuity assumption of $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+1})$ required in Theorem 3.3. Moreover $P_0^{(k)} = P((if)^k \cdot) \in \bigcap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k})$. We also apply Theorem 3.3 with $\mathcal{B}_r \hookrightarrow \mathcal{B}_{r+1}$ and with $\mathcal{B}_{r+1} \hookrightarrow \mathbb{L}^1(\mu)$. The above inequalities combined with the Keller and Liverani perturbation theorem (Theorem 3.3) ensure that the assumptions of Items (C) and (D) of Proposition 4.1. Thus Corollary 5.6 with $\nu = \mu$ and Proposition 5.8 apply. For (i), we apply Theorem 6.2 thanks to the previous facts and we identify the variance matrix of the limit using Proposition 5.8.

Condition (P) of [26] is satisfied since the $\mathbb{L}^p(\mu)$ are contained and dense in $\mathbb{L}^1(\mu)$. Condition (\widehat{K}) of [26] follows from the ρ -mixing. Thus, due to [26, Proposition 5.4], the non-lattice property implies the non-arithmeticity (exponential decay of $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]|$ as $n \rightarrow +\infty$). We conclude by applying Theorem 7.1. \square

¹¹non-latticity for \mathbb{Z}^d -valued observables means that there exist no triple (a, H, θ) with $a \in \mathbb{Z}^d$, $H \neq \mathbb{Z}^d$ a closed subgroup in \mathbb{Z}^d and $\theta : \Omega \rightarrow \mathbb{Z}^d$ such that $Y_1 + \theta(X_1) - \theta(X_0) \in a + H$ $\mathcal{P}_\mu \times \mathbf{P}$ -a.s.

¹²non-latticity in \mathbb{R} means that there exist no triple $a, b \in \mathbb{R}$ and $\theta : \Omega \rightarrow \mathbb{R}$ such that $Y_1 + \theta(X_1) - \theta(X_0) \in a + b\mathbb{Z}$ $\mathcal{P}_\mu \times \mathbf{P}$ -a.s..

We illustrate the previous theorem with the following explicit example with a smooth transition density and an observable that does not admit moment of every order.

Example 7.4. Let $\kappa > 0$. We consider the following von Mises Markov chain $(X_n)_{n \geq 0}$ on $\Omega = [-\frac{1}{2}, \frac{1}{2}]$ endowed with the Lebesgue measure μ (we also take $\nu = \mu$). We assume that conditionally to X_k , X_{k+1} has von Mises distribution with mean X_k and with concentration κ . Let r be a positive integer and $\gamma \in]r+1, r+2[$. We consider $S_n = \sum_{k=1}^n Y_k$, with $Y_k = f(X_k)$ and $f(y) := \text{sgn}(y)|y|^{-\frac{1}{\gamma}}$. Then $Y_k = f(X_k) \in \mathbb{L}^{r+1} \setminus \mathbb{L}^{r+2}$ and

- (i) Theorem 6.2 (CLT) applies and so $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable with variance matrix $\sum_{n \in \mathbb{Z}} \text{Cov}_{\mathcal{P}_\mu}(Y_1, Y_{|n|+1})$.
- (ii) Theorem 6.5 (LLT) applies.
- (iii') Theorems 6.11 and 6.12 (Expansions of order $\hat{r} - 1 = r - 2$ and $\hat{r} = r - 1$ in the LLT) apply (for any α_1 and α'_1).
- (iv) Corollary 6.14 (First order Edgeworth expansion) holds true if $r \geq 2$.
- (v) Theorem 6.13 (Edgeworth expansion of order $\hat{r} = r - 1$) applies if $r \geq 2$ (for $r' = \hat{r} = r - 1$).

Proof. We apply Theorem 7.3. This Markov chain has a smooth transition density. The family of Fourier-perturbed operators $(P_t)_t$ is given by

$$\forall t \in \mathbb{R}, \quad P_t(u)(x) = \int_{\Omega} p(x, y) e^{itf(y)} u(y) dy,$$

with $p(x, y) = \frac{e^{\kappa \cos(2\pi(y-x))}}{I_0(\kappa)}$, where I_0 is the modified Bessel function of order 0.

- We first observe that the Lebesgue measure μ on $\Omega := [-\frac{1}{2}, \frac{1}{2}]$ is invariant. This follows from the fact that $\int_{\Omega} p(x, y) dx = 1$.
- We observe that

$$\forall |t| > 2^{\frac{1}{\gamma}}, \quad \mu(|f|^k > t) = \mu(|f|^k > t^{\frac{1}{k}}) = \mu(|\text{Id}| < t^{-\frac{\gamma}{k}}) = 2t^{-\frac{\gamma}{k}}.$$

Thus

$$\int_{\Omega} |f|^k d\mu = \int_0^{+\infty} \mu(|f|^k > t) dt$$

is finite if and only if $\gamma > k$. This ensures that $f \in \mathbb{L}^{r+1}(\mu) \setminus \mathbb{L}^{r+2}(\mu)$.

- Now let us prove that the Markov chain $(X_n)_n$ is ρ -mixing. To this end, we establish a Doeblin-Fortet inequality and study the peripheral spectrum. For any $n \in \mathbb{N}$ and $x \in \Omega$, the following inequalities hold true:

$$\begin{aligned} P^n(u)(x) &\leq \int_{\Omega^n} \left(\prod_{i=0}^{n-1} p(x_i, x_{i+1}) \right) |u(x_n)| dx_1 \dots dx_n \\ &\leq \frac{e^\kappa}{I_0(\kappa)} \|u\|_{L^1(\mu)}, \end{aligned}$$

with the convention $x_0 = x$, where we used the fact that $p(x_{n-1}, x_n) \leq \frac{e^\kappa}{I_0(\kappa)}$ and that, for all $i = 0, \dots, n-2$, $p(x_i, \cdot)$ are probability density functions. Thus we have proved that, for any $p \in [1, +\infty]$ and any $n \in \mathbb{N}$,

$$\|P^n(u)\|_{\mathbb{L}^p(\mu)} \leq \frac{e^\kappa}{I_0(\kappa)} \|u\|_{L^1(\mu)}. \quad (46)$$

This Doeblin-Fortet inequality combined with e.g. Hennion's theorem [18] ensures the quasi-compactness of P on $\mathbb{L}^p(\mu)$ for all $p \in]1, +\infty[$ and the nullity of its spectral radius.

Furthermore $P\mathbf{1} = \mathbf{1}$ and if $h \in \mathbb{L}^p(\mu)$ and $\lambda \in \mathbb{C}$ of modulus one are such that $Ph = \lambda h$, then

$$|h(x)| \leq \int_{\Omega} p(x, y) |h(y)| dy.$$

But

$$\int_{\Omega} |h(x)| dx = \int_{\Omega^2} p(x, y) |h(y)| dy dx,$$

since μ is invariant. Since $p(x, y) > 0$, this implies that $|h|$ is almost surely constant, and the relation

$$\lambda h(x) = \int_{\Omega} p(x, y) h(y) dy$$

then implies that h is constant and so also that $\lambda = 1$. Hence 1 is the only eigenvalue of modulus one of P and its eigenspace consists in constant functions.

Finally (46) ensures that $\sup_n \|P^n\|_{\mathbb{L}^2(\mu)} < \infty$ and so that the generalized eigenspace of P associated to 1 coincide with its eigenspace, ending the proof of the simplicity of 1 as an eigenvalue of P . Thus we have proved that the Markov chain $(X_n)_n$ is ρ -mixing.

- Now let us prove Assumptions (β') , (γ') and (δ') . To this end we will prove that $\lim_{t \rightarrow +\infty} \|P_t^2\|_{\infty} = 0$. Observe that $P_t^2(u)(x) = \int_{\Omega} q_t(x, z) e^{itf(z)} u(z) dz$, with

$$q_t(x, z) = \int_{\Omega} p(x, y) p(y, z) e^{itf(y)} dy = \int_{|y| > 2^{\frac{1}{\gamma}}} p(x, y^{-\gamma}) p(y^{-\gamma}, z) e^{ity\gamma|y|^{-\gamma-1}} dy.$$

The Riemann-Lebesgue theorem ensures that $\lim_{t \rightarrow +\infty} q_t(x, z) = 0$. By uniform continuity of p and since Ω is compact, $(x, z) \mapsto q_t(x, z)$ is also uniformly continuous (uniformly in (x, z, t)). Since Ω is compact, we conclude that $q_t(x, z)$ converges to 0 uniformly in (x, z) , as $t \rightarrow +\infty$. Thus $\lim_{t \rightarrow +\infty} \|P_t^2\|_{\infty} = 0$, ensuring (42) and so Assumptions $(\gamma')[\alpha'_1, \alpha_1]$ for any α_1, α'_1 and $(\delta')[r']$ for any r' and also Assumption (β') (noticing that the lattice condition $f(y) + \theta(y) - \theta(x) \in a + b\mathbb{Z}$ (with $b \neq 0$) would imply that $P_{\frac{2\pi k}{b}}(e^{\frac{2i\pi k}{b}\theta}) = e^{\frac{2i\pi k}{b}\theta}$).

We conclude by applying 7.3. □

7.3. V -geometrically ergodic Markov chains. Let $V : \Omega \rightarrow [1, +\infty[$ be an unbounded measurable function. The random walk $(X_n)_n$, or equivalently its transition operator P , is said to be V -geometrically ergodic if there exist $C > 0$ and $\vartheta \in]0, 1[$ such that

$$\forall n \geq 1, \quad \left\| \frac{P^n(\cdot) - \mathbb{E}_{\pi}[\cdot]}{V} \right\|_{\infty} \leq C \vartheta^n \|\cdot / V\|_{\infty}.$$

Again the study of Markov random walks driven by a V -geometrically ergodic random walk $(X_k)_{k \geq 0}$ can be reduced to an additive function of a V -geometrically ergodic Markov chain.

Proposition 7.5. *If the Markov chain $(X_k)_{k \geq 1}$ is V -geometrically ergodic, then the Markov chain $(\tilde{X}_k = (X_{k-1}, X_k, Z_k))_{k \geq 0}$ with invariant measure $\tilde{\mu}$ defined in Proposition 7.2 is \tilde{V} -geometrically ergodic with same rate, with $\tilde{V}(x, y, z) = V(x) + V(y)$.*

Proof. Let $G : \Omega^2 \times E \rightarrow \mathbb{C}$ be a bounded measurable function. We have seen in the proof of Proposition 7.2 that the transfer operator \tilde{P} of \tilde{X} satisfies

$$\forall n \geq 2, \quad \tilde{P}^n((G\tilde{V})(x, y, z)) = \int_E (P^{n-1}(H(\cdot, z)))(y) d\mathbf{P}(z),$$

where $H(x, z) = \mathbb{E}[G(X_0, X_1, z)(V(x) + V(X_1)) | X_0 = x] =: K(x, z)V(x)$. Moreover

$$\begin{aligned} \|K\|_\infty &\leq \frac{\mathbb{E}[\|G\|_\infty(V(X_0) + V(X_1)) | X_0 = x]}{V(x)} \\ &\leq \|G\|_\infty (1 + \|P(V)/V\|_\infty) . \end{aligned}$$

Thus

$$\begin{aligned} \left\| \frac{\tilde{P}^n(G\tilde{V}) - \mathbb{E}_{\tilde{\mu}}[G\tilde{V}]}{\tilde{V}} \right\|_\infty &= \left\| \frac{1}{\tilde{V}} \int_E \left(P^{n-1}(K(\cdot, z)V) - \mathbb{E}_\mu[K(\cdot, z)V] \right) d\mathbf{P}(z) \right\|_\infty \\ &\leq C\vartheta^{n-1} \int_E \|K(\cdot, z)\|_\infty d\mathbf{P}(z) \leq C\vartheta^{n-1} \|G\|_\infty (1 + \|P(V)/V\|_\infty) . \end{aligned}$$

This ends the proof of the proposition since P acts continuously on the space $V.L^\infty$ endowed with the norm $\|\cdot/V\|_\infty$. \square

Thus, without any loss of restriction, from now on, in this subsection, we replace $f(x, y, z)$ by $f(y)$ (again this complicates the notion of non-lattice, see [25] for a simple one).

Theorem 7.6. *Let r be a positive integer, $r' \in]r, r+1]$ be a real number and $f : \Omega \rightarrow \mathbb{R}^d$ be a μ -centered function belonging to $\mathbb{L}^{r+1}(\mu)$. Let $V : \Omega \rightarrow [1, +\infty[$ be an unbounded measurable function such that $\mathbb{E}_\mu[V] < \infty$. Assume that P is V -geometrically ergodic and that*

$$\max_{u \in \{0, r'-r\}} \sup_{j=0, \dots, r} \sup_{m=1, \dots, r-j} \left\| V^{-\frac{j+m+u}{r+1}} P \left(|f|^{m+u} V^{\frac{j}{r+1}} \right) \right\|_\infty < \infty , \quad (47)$$

then the conclusions of Items (C) and (D) of Propositions 4.1, of Corollary 5.6 (for $\nu = \mu$) and of Proposition 5.8 hold true with $\mathcal{B}_0 = \mathbb{C} \cdot \mathbf{1}$ endowed with the infinite norm and with, for $j = 1, \dots, r$, $\mathcal{B}_j := V^{\frac{j}{r+1}}.L^\infty$ and $\mathcal{B}_{r+1} := V^{\frac{r'}{r+1}}.L^\infty$ endowed with the respective norms $\|\cdot\|_{(j)} := \|\cdot/V^{\frac{j}{r+1}}\|_\infty$ and $\|\cdot\|_{(r+1)} := \|\cdot/V^{\frac{r'}{r+1}}\|_\infty$. If $\mathbb{E}_\nu[V^\varepsilon] < \infty$ for some $\varepsilon > 0$, then

- (i) Theorem 6.2 (CLT) applies and we conclude that $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable with variance matrix $\sum_{n \in \mathbb{Z}} \text{Cov}_{\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}} (Y_1, Y_{|n|+1})$.
- (ii) Theorem 6.5 (LLT) applies if Y_1 is \mathbb{Z}^d valued and non-lattice in \mathbb{Z}^d .

If moreover $\mathbb{E}_\nu[V^{\frac{r}{r+1}}] < \infty$, then

- (iii) ($d = 1$) Theorems 6.11 and 6.12 apply with $\hat{r} = r - 1$ if f is \mathbb{Z} -valued and non-lattice in \mathbb{Z} .
- (iii') ($d = 1$) Theorems 6.11 and 6.12 (Expansions of order r' in the LLT) apply with $\hat{r} = r - 1$ if Assumption $(\gamma')[\alpha'_1, \alpha_1]$ holds true for $\mathbb{E}_{\mathcal{P}_\nu}[e^{isS_n}]$ and if Y_1 is non-lattice in \mathbb{R} .
- (iv) ($d = 1$) Corollary 6.14 (First order Edgeworth expansion) holds true if $r = 2$ and f is non-lattice in \mathbb{R} .
- (v) ($d = 1$) Theorem 6.13 (Edgeworth expansion of order $r' \in [1, \hat{r} = r - 1]$) applies if Y_1 is non-lattice in \mathbb{R} and if Assumption $(\delta')[r']$ holds true for $\mathbb{E}_{\mathcal{P}_\nu}[e^{isS_n}]$.

Proof. Recall that P is also V^γ -geometrically ergodic for any $\gamma \in]0, 1]$ and that Theorem 3.3 applies with the Banach spaces $V^\gamma L^\infty \hookrightarrow \mathbb{L}^1(\mu)$ (see e.g. [26, Lemma 10.1]). If $V \in \mathbb{L}^1(\mathbf{m})$, then $\mathcal{B}_j \subset \mathbb{L}^{\frac{r+1}{j}}(\mathbf{m})$ since

$$\|h\|_{\mathbb{L}^{\frac{r+1}{j}}(\mathbf{m})} = \left\| V^{\frac{j}{r+1}} \frac{h}{V^{\frac{j}{r+1}}} \right\|_{\mathbb{L}^{\frac{r+1}{j}}(\mathbf{m})} \leq \left\| V^{\frac{j}{r+1}} \right\|_{\mathbb{L}^{\frac{r+1}{j}}(\mathbf{m})} \left\| \frac{h}{V^{\frac{j}{r+1}}} \right\|_\infty = \|V\|_{\mathbb{L}^1(\mathbf{m})}^{\frac{j}{r+1}} \|h\|_{\mathcal{B}_j} .$$

In particular, if $V \in \mathbb{L}^1(\nu)$, then ν defines a linear continuous form on \mathcal{B}_{r+1} . The fact that $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$ follows from $V^{-\frac{j+1}{r}} \leq V^{-\frac{j}{r}}$. There exist $C > 0$ and $\vartheta \in]0, 1[$ such that, for any $j \in \{1, \dots, r+1\}$,

$$\|P^n(g) - \mathbb{E}_\pi[g]\|_{(j)} \leq C\vartheta^n \|g\|_{(j)},$$

$$\|P_t^n(g)\|_{(j)} \leq \|P^n(|g|)\|_{(j)} \leq C\vartheta^n \|g\|_{(j)} + \|g\|_1,$$

ensuring the Doeblin-Fortet estimate (4) for $\mathcal{B}_j \hookrightarrow \mathcal{B}_{j+1}$ for $j = 1, \dots, r$. For $k = 0, \dots, r$, $P_0^{(k)} := P((if)^k \cdot) \in \cap_{j=0}^{r-k} \mathcal{L}(\mathcal{B}_j, \mathcal{B}_{j+k})$ since

$$\begin{aligned} \|P_0^{(k)}(h)\|_{(j+k)} &\leq \left\| V^{-\frac{j+k}{r+1}} P(|f|^k |h|) \right\|_\infty \\ &\leq \left\| V^{-\frac{j+k}{r+1}} P(|f|^k V^{\frac{j}{r+1}}) \right\|_\infty \|h\|_{(j)} \end{aligned}$$

and due to (47). It remains to check the regularity assumptions. For all $j = 0, \dots, r$ and $m \in \{1, \dots, r-j\}$,

$$\begin{aligned} \left\| P_t(g) - \sum_{k=0}^m \frac{P((if)^k g)}{k!} t^k \right\|_{(j+m)} &\leq \left\| P \left(\left(e^{itf} - \sum_{k=0}^m \frac{(if)^k}{k!} t^k \right) g \right) \right\|_{(j+m)} \\ &\leq \left\| V^{-\frac{j+m}{r+1}} P \left(|tf|^m \|g\|_{(j)} V^{\frac{j}{r+1}} \right) \right\|_\infty \\ &\leq \mathcal{O}(t^m) \|g\|_{(j)}, \end{aligned}$$

and

$$\begin{aligned} \left\| P_t(g) - \sum_{k=0}^m \frac{P((if)^k g)}{k!} t^k \right\|_{(r+1)} &\leq \left\| P \left(|tf|^{m+r'-r} V^{\frac{j}{r+1}} \right) V^{-\frac{j+m+r'-r}{r+1}} \right\|_\infty \|g\|_{(j)} \\ &\leq \mathcal{O}(t^{m+r'-r}) \|g\|_{(j)} = o(t^m) \|g\|_{(j)}. \end{aligned}$$

Condition (P) of [26] is satisfied since the Banach spaces are stables under complex modulus and Condition (\widehat{K}) of [26] follows from the V -geometric ergodicity. Thus, due to [26, Proposition 5.4], the non-lattice property implies the exponential decay of $\sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} |\mathbb{E}[e^{it \cdot S_n}]|$ as $n \rightarrow +\infty$. We end the proof by applying Theorem 7.1 with $\tilde{\mathcal{B}}_0 := V^\varepsilon.L^\infty$ for (applying Theorem 3.3 for $V^{\frac{\varepsilon}{2}}.L^\infty \hookrightarrow \tilde{\mathcal{B}}_0 = V^\varepsilon.L^\infty$). \square

Remark 7.7. Assume $\nu = \mu$. Observe that (47) is satisfied as soon as $\|f^{r+1}/V\|_\infty < \infty$ (which also ensures that $f \in \mathbb{L}^{r+1}(\mu)$) with $r \geq 1$, since for $j+m \leq r$ and $u \in \{0, r'-r\}$:

$$\left\| V^{-\frac{j+m+u}{r+1}} P \left(|f|^{m+u} V^{\frac{j}{r+1}} \right) \right\|_\infty \leq \left\| V^{-\frac{j+m+u}{r+1}} P \left(V^{\frac{j+m+u}{r+1}} \right) \right\|_\infty \|f^{r+1}/V\|_\infty^{\frac{m+u}{r+1}} < \infty.$$

Moreover, for the Markov chain (X_{k-1}, X_k, Z_k) considered in Proposition 7.5 with the reference function \tilde{V} , (47) holds true if $\left\| \int_E |f(\cdot, \cdot, \omega)|^{r+1} d\mathbf{P}(\omega) / \tilde{V} \right\|_\infty < \infty$. Indeed, setting $g(x, y) =$

$$\begin{aligned}
& \frac{\int_E |f(x, y, \omega)|^{r+1} d\mathbf{P}(\omega)}{V(x) + V(y)}, \\
& \left| \tilde{P} \left(|f|^{m+u} \tilde{V}^{\frac{j}{r+1}} \right) (x) \right| = \left| P \left(\int_E |f(x, \cdot, \omega)|^{m+u} d\mathbf{P}(\omega) (V(x) + V(\cdot))^{\frac{j}{r+1}} \right) (x) \right| \\
& \leq \left| P \left(\left(\int_E |f(x, \cdot, \omega)|^{r+1} d\mathbf{P}(\omega) \right)^{\frac{m+u}{r+1}} (V(x) + V(\cdot))^{\frac{j}{r+1}} \right) (x) \right| \\
& \leq \left| P \left((g(x, \cdot))^{\frac{m+u}{r+1}} (V(x) + V(\cdot))^{\frac{j+m+u}{r+1}} \right) (x) \right| \\
& \leq \|g\|_{\infty}^{\frac{m+u}{r+1}} \left(1 + \left\| V^{-\frac{j+m+u}{r+1}} P(V^{\frac{j+m+u}{r+1}}) \right\|_{\infty} \right) V(x)^{\frac{j+m+u}{r+1}},
\end{aligned}$$

since P is $V^{\frac{j+m+u}{r+1}}$ -geometrically ergodic.

7.4. Lipschitz iterative models. We consider a non-compact metric space (E, d) in which every closed ball is compact, and endow it with its Borel σ -algebra \mathcal{E} . Let (G, \mathcal{G}) be a measurable space, let $(\theta_n)_{n \geq 1}$ be a sequence of independent identically distributed G -valued random variables. Let $F : E \times G \rightarrow E$ be a measurable function such that, for all $g \in G$, $F_g := F(\cdot, g)$ is Lipschitz continuous with Lipschitz constant \mathcal{C}_g . We consider the random walk defined by

$$\forall n \geq 1, \quad X_n = F(X_{n-1}, \theta_n),$$

the sequence $(\theta_n)_{n \geq 1}$ being independent of the initial value X_0 of the random walk. This random walk is called an iterative Lipschitz model [5] [7] and has transition operator

$$P(g)(x) = \mathbb{E}[g(F(x, \theta_1))].$$

Note that this context includes the autoregressive chains on \mathbb{R}^d of the form

$$\forall n \geq 1, \quad X_n = A_n X_{n-1} + \theta_n,$$

where $(A_n, \theta_n)_{n \geq 1}$ is an i.i.d. sequence of r.v. taking values in $\mathcal{M}_d(\mathbb{R}) \times \mathbb{R}^d$, independent of X_0 ($\mathcal{M}_d(\mathbb{R})$ denotes the set of real $d \times d$ -matrices.) with e.g. $|A_1| < 1$ ($|\cdot|$ being a matrix norm). For one-dimensional autoregressive chains, convergence to stable laws has been investigated in [17] for $f(x) = x$. Let x_0 be a fixed point in E . For $x \in E$, we set $p(x) = 1 + d(x, x_0)$. We set, for all $g \in G$, $\mathcal{M}_g := \mathcal{C}_g + p(F(x_0, g))$. We are interested in the asymptotic behaviour of $S_n = \sum_{k=1}^n f(X_k)$, with $f : \Omega \rightarrow \mathbb{R}^d$ satisfying the following condition for some $C_1, s \geq 0$:

$$\forall (x, y) \in E \times E, \quad |f(x) - f(y)| \leq C_1 d(x, y) (p(x) + p(y))^s. \quad (48)$$

Recall that the case f is Lipschitz continuous (i.e. $s = 0$ in (48)) has been studied e.g. in [7] and in [2]. Fix $\alpha \in]0, 1]$. We set $\tilde{\mathcal{B}}_{\alpha, b, a}^{(0)}$ for the set of functions $g : E \rightarrow \mathbb{C}$ such that $|f|_a^{(0)} + m_{\alpha, b}^{(0)}(f) < \infty$, with

$$|f|_a^{(0)} := \|f/p^a\|_{\infty} \quad \text{and} \quad m_{\alpha, b}^{(0)}(f) := \sup_{x, y \in E, x \neq y} \frac{|f(x) - f(y)|}{d(x, y)^{\alpha} \max(p(x), p(y))^b}.$$

Remark 7.8. Observe that (48) means that the coordinates of f belong to $\tilde{\mathcal{B}}_{1, s, s+1}^{(0)} \subset \tilde{\mathcal{B}}_{\alpha, s+1-\alpha, s+1}^{(0)}$.

In the sequel, we set $\mathcal{M} := \mathcal{M}_{\theta_1}$ and $\mathcal{C} := \mathcal{C}_{\theta_1}$. Recall that it has been proved in [20, Th. I] that if

$$\mathbb{E} [\mathcal{M}^{\alpha(s+1)}] < \infty \quad \text{and} \quad \mathbb{E} [\mathcal{C}^{\alpha} \max\{\mathcal{C}, 1\}^{\alpha s}] < 1,$$

then the Markov chain $(X_n)_n$ admits a unique stationary distribution μ and $p \in \mathbb{L}^{\alpha(s+1)}(\mu)$ (see also e.g. [5] [7]), which implies in particular that $f \in \mathbb{L}^{\alpha}(\mu)$.

Theorem 7.9. *Let r be a positive integer and a real number $r' \in]r, r+1]$. Let ν be a probability measure on Ω . Assume that $p \in \mathbb{L}^{(s+1)(r+1)}(\mu)$ and $\alpha \in]0, \frac{s+1}{s+2}]$. Assume that f is μ -centered, satisfies (48) and that¹³*

$$\mathbb{E} \left[\mathcal{M}^{(s+1)(r+1)} + \mathcal{C}^\alpha \mathcal{M}^{(s+1)r'+s\alpha} \right] < +\infty \quad (49)$$

and

$$\mathbb{E} \left[\mathcal{C}^\alpha \max\{\mathcal{C}, 1\}^{(s+1)r'+s\alpha} \right] < 1. \quad (50)$$

If $p \in \mathbb{L}^\varepsilon(\nu)$ for some $\varepsilon > 0$, then

- (i) Theorem 6.2 (CLT) applies and we conclude that $(S_n/\sqrt{n})_{n \geq 1}$ converges in distribution to a centered Gaussian random variable with variance matrix $\sum_{n \in \mathbb{Z}} \text{Cov}_{\mu \otimes \mathbf{P}^{\otimes \mathbb{N}}} (Y_1, Y_{|n|+1})$.
- (ii) Theorem 6.5 (LLT) applies if Y_1 is \mathbb{Z}^d valued and non-lattice in \mathbb{Z}^d .

If moreover $p \in \mathbb{L}^{(s+1)r}(\nu)$, then

- (iii) ($d = 1$) Theorems 6.11 and 6.12 apply with $\hat{r} = r - 1$ if f is \mathbb{Z} -valued and non-lattice in \mathbb{Z} .
- (iii') ($d = 1$) Theorems 6.11 and 6.12 (Expansions of order r in the LLT) apply with $\hat{r} = r - 1$ if Assumption (γ') holds true for $\mathbb{E}_\nu[e^{isS_n}]$ and if Y_1 is non-lattice in \mathbb{R} .
- (iv) ($d = 1$) Corollary 6.14 (First order Edgeworth expansion) holds true if $r \geq 2$ and if f is non-lattice in \mathbb{R} .
- (v) ($d = 1$) Theorem 6.13 (Edgeworth expansion of order $r' \in [1, \hat{r} = r - 1]$) applies if $r \geq 2$, if Y_1 is non-lattice in \mathbb{R} and if Assumption $(\delta')[r']$ holds true for $\mathbb{E}_\nu[e^{isS_n}]$.

Let us prove this result. We consider the following notion of weighted Hölder-type spaces due to D. Guibourg [14] and used in [26] generalizing those introduced [32] (used also in [34, 37]). For positive real numbers β, γ such that $0 < \beta \leq \gamma$ and for $(x, y) \in E^2$, we set

$$\Delta_{\alpha, \beta, \gamma}(x, y) := p(x)^{\alpha\gamma} p(y)^{\alpha\beta} + p(x)^{\alpha\beta} p(y)^{\alpha\gamma}.$$

Then $\mathcal{B}_{\alpha, \beta, \gamma}$ denotes the space of \mathbb{C} -valued functions g on E satisfying the following condition

$$m_{\alpha, \beta, \gamma}(g) := \sup \left\{ \frac{|g(x) - g(y)|}{d(x, y)^\alpha \Delta_{\alpha, \beta, \gamma}(x, y)}, \quad x, y \in E, \quad x \neq y \right\} < +\infty.$$

Set $|g|_{\alpha, \gamma} := \sup_{x \in E} \frac{|g(x)|}{p(x)^{\alpha(\gamma+1)}}$ and $\|g\|_{\alpha, \beta, \gamma} := m_{\alpha, \beta, \gamma}(g) + |g|_{\alpha, \gamma}$. Then $(\mathcal{B}_{\alpha, \beta, \gamma}, \|\cdot\|_{\alpha, \beta, \gamma})$ is a Banach space. Moreover

$$\gamma \leq \gamma' \quad \Rightarrow \quad \mathcal{B}_{\alpha, \beta, \gamma} \hookrightarrow \mathcal{B}_{\alpha, \beta, \gamma'}, \quad (51)$$

and $\mathcal{B}_{\alpha, \beta, \gamma} \hookrightarrow \mathbb{L}^{\frac{K}{\gamma+1}}(\mu)$ if $p \in \mathbb{L}^{\alpha K}(\mu)$. Theorem 7.9 relies on the following propositions combined with Theorem 7.1

Proposition 7.10. *Assume the Assumptions of Theorem 7.9. Then the assumptions of Corollary 5.6 are satisfied with $\mathcal{B}_0 = \mathbb{C} \cdot \mathbf{1}$, $\mathcal{B}_j = \mathcal{B}_{\alpha, s+1, \frac{j(s+1)}{\alpha} - 1}$ for $j = 1, \dots, r$, $\mathcal{B}_{r+1} = \mathcal{B}_{\alpha, s+1, \frac{r'(s+1)}{\alpha} - 1}$, $\mathcal{B}_{r+2} = \mathbb{L}^1(\mu)$.*

Proof. Since $\alpha \in]0, \frac{s+1}{s+2}]$, $s+1 \leq \frac{s+1}{\alpha} - 1$. Due to Lemmas 7.11 and 7.12, we know that, for all for $j = 1, \dots, r+1$, P is geometrically ergodic on \mathcal{B}_j and that $(P_t)_t$ satisfies the Doeblin Fortet assumption of Theorem 3.3 for the Banach spaces $\mathcal{B}_j \hookrightarrow \mathbb{L}^1(\mu)$. Moreover $\mathbf{1} \in \mathcal{B}_0$, and P coincide with Id on \mathcal{B}_0 and so $(z \text{Id} - P)^{-1} = (z - 1)^{-1} \text{Id}$ on \mathcal{B}_0 . Due to Lemma 7.13, $P_t - \sum_{j=0}^k \frac{P_0^{(j)} \cdot t^{\otimes j}}{j!}$

¹³Observe that these assumptions ensure that $p \in \mathbb{L}^{(s+1)(r+1)}(\mu)$, thus that $f \in \mathbb{L}^{r+1}(\mu)$.

is in $\mathcal{O}(|t|^k)$ in $\mathcal{L}(\mathcal{B}_i, \mathcal{B}_{i+k})$ for all $i = 0, \dots, r-k$ and in $\mathcal{O}(|t|^{k+r'-r}) = o(t^k)$ in $\mathcal{L}(\mathcal{B}_{r-k}, \mathcal{B}_{r+1})$, with $P_0^{(j)} := P((if)^{\otimes j})$.

Moreover $\mathcal{B}_j \hookrightarrow p^{(s+1)j} L^\infty \hookrightarrow \mathbb{L}^{\frac{r+1}{j}}(\mu)$ since $p^{(s+1)(r+1)} \in \mathbb{L}^1(\mu)$; and $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_{r+1}, \mathcal{B}_{r+2})$ is continuous (we can use e.g. Lemma 7.13). \square

Lemma 7.11. ([26, Lemma B.1]) *Let $a \geq \alpha$ and $b \in [0, a + \alpha\gamma]$ and $0 \leq \beta \leq \gamma$. Assume*

$$\mathbb{E} \left[\mathcal{M}^{a+\alpha(\gamma+1)} + \mathcal{C}^\alpha \mathcal{M}^{a+\alpha(\gamma+\beta)} + \mathcal{C}^\alpha \mathcal{M}^{b+\alpha(\gamma+1)} \right] < \infty.$$

Then there exists $C' > 0$ such that, for all $g \in \tilde{\mathcal{B}}_{\alpha, a, b}^{(0)}$,

$$\|P(g \cdot)\|_{\mathcal{L}(\mathcal{B}_{\alpha, \beta, \gamma}, \mathcal{B}_{\alpha, \max(\beta, \frac{b-a+\alpha}{\alpha}), \gamma + \frac{a}{\alpha}})} \leq C' \|g\|_{\tilde{\mathcal{B}}_{\alpha, b, a}^{(0)}},$$

and

$$\|P(g \cdot)\|_{\mathcal{L}(\mathbb{C}, \tilde{\mathcal{B}}_{\alpha, b, a}^{(0)})} \leq C' \|g\|_{\tilde{\mathcal{B}}_{\alpha, b, a}^{(0)}}.$$

More precisely

$$|P(gh)|_{\alpha, \gamma + \frac{a}{\alpha}} \leq \mathbb{E}[\mathcal{M}^{a+\alpha(\gamma+1)}] |g|_a^{(0)} |h|_{\alpha, \gamma} \quad (52)$$

and

$$m_{\alpha, \max(\beta, \frac{b-a+\alpha}{\alpha}), \gamma + \frac{a}{\alpha}}(gh) \leq \mathbb{E}[\mathcal{C}^\alpha \mathcal{C}_1^{a+\alpha(\gamma+\beta)}] |g|_a^{(0)} m_{\alpha, \beta, \gamma}(h) + \mathbb{E}[\mathcal{C}^\alpha \mathcal{M}^{b+\alpha(\gamma+1)}] m_{\alpha, b}^{(0)}(g) |h|_{\alpha, \gamma}, \quad (53)$$

with $\mathcal{C}_1 := \max(\mathcal{C}, 1) + d(F(x_0, \theta_1), x_0) \leq \mathcal{M}$.

Proof. The fact that $\mathcal{C}_1 \leq \mathcal{M}$ follows from [20, p.1945]. We start from [26, Lemma B.1] which states that

$$|P(gh)| \leq \mathbb{E}[\mathcal{M}^{a+\alpha(\gamma+1)}] |g|_a^{(0)} |h|_{\alpha, \gamma} p^{a+\alpha(\gamma+1)}$$

(ensuring (52)) and that

$$\begin{aligned} \frac{|P(gh)(x) - P(gh)(y)|}{d(x, y)^\alpha} &\leq \mathbb{E}[\mathcal{C}^\alpha \mathcal{C}_1^{a+\alpha(\gamma+\beta)}] |g|_a^{(0)} m_{\alpha, \beta, \gamma}(h) p(x)^a \Delta_{\alpha, \beta, \gamma}(x, y) \\ &\quad + \mathbb{E}[\mathcal{C}^\alpha \mathcal{M}^{b+\alpha(\gamma+1)}] m_{\alpha, b}^{(0)}(g) |h|_{\alpha, \gamma} p(x)^b p(y)^{\alpha(\gamma+1)}, \end{aligned}$$

if $p(y) \leq p(x)$. Still assuming that $p(y) \leq p(x)$, we conclude by noticing that

$$\begin{aligned} p(x)^a \Delta_{\alpha, \beta, \gamma}(x, y) &= p(x)^a (p(x)^{\alpha\beta} p(y)^{\alpha\gamma} + p(y)^{\alpha\beta} p(x)^{\alpha\gamma}) \\ &\leq p(x)^{a+\alpha\gamma} p(y)^{\alpha\beta} + p(y)^{\alpha\beta} p(x)^{a+\alpha\gamma} \end{aligned}$$

since $\beta \leq \gamma$ and

$$p(x)^b p(y)^{\alpha(\gamma+1)} \leq p(x)^{a+\alpha\gamma} p(y)^{\alpha+b-a}$$

since $\alpha(\gamma+1) \geq \alpha + b - a$. \square

Lemma 7.12. ([20, Theorem 11.5], [26, Propositions 11.2 and 11.4]) *Assume $s+1 \leq \beta \leq \gamma$ and*

$$\mathbb{E} \left[\mathcal{M}^{\alpha(\gamma+1)} + \mathcal{C}^\alpha \mathcal{M}^{\alpha(\gamma+\beta)} \right] < +\infty \quad \text{and} \quad \mathbb{E} \left[\mathcal{C}^\alpha \max\{\mathcal{C}, 1\}^{\alpha(\gamma+\beta)} \right] < 1. \quad (54)$$

Then P is geometrically ergodic on $\mathcal{B}_{\alpha, \beta, \gamma}$: $\|P^n - \mathbb{E}_\mu\|_{\mathcal{L}(\mathcal{B}_{\alpha, \beta, \gamma})} = \mathcal{O}(\vartheta^n)$ and $(P_t)_t$ satisfies the assumptions of Theorem 3.3 with $\mathcal{B}_1 = \mathcal{B}_{\alpha, \beta, \gamma}$ and $\mathcal{B}_2 = \mathbb{L}^1(\mu)$.

We prove the following result (generalizing [26, Proposition 11.6], see also ([26, Propositions 11.5 and 11.7] for \mathcal{C}^m -smoothness).

Lemma 7.13. *If $s+1 \leq \beta \leq \gamma < \gamma'$ satisfies $\alpha(\gamma' - \gamma) \leq s+1$ and*

$$\mathbb{E}[\mathcal{M}^{\alpha(\gamma'+1)} + \mathcal{C}^\alpha \mathcal{M}^{\alpha(\gamma'+\beta)}] < \infty,$$

then $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_{\alpha,\beta,\gamma}, \mathcal{B}_{\alpha,\beta,\gamma'})$ is continuous at 0.

If m is a positive integer, if $m' \in [m, m+1]$ and if β, γ are such that $0 \leq \beta \leq \gamma$ and

$$\mathbb{E} \left[\mathcal{M}^{m'(s+1)+\alpha(\gamma+1)} + \mathcal{C}^\alpha \mathcal{M}^{m'(s+1)+\alpha(\gamma+\beta)} + \mathcal{C}^\alpha \mathcal{M}^{m'(s+1)+\alpha\gamma} \right] < \infty.$$

Then

$$P_t - \sum_{k=0}^m \frac{P_0^{(k)} \cdot t^{\otimes k}}{k!} = \mathcal{O}(t^{m'})$$

in $\mathcal{L} \left(\mathcal{B}_{\alpha,\beta,\gamma}, \mathcal{B}_{\alpha,\beta,\gamma+\frac{m'(s+1)}{\alpha}} \right)$ and in $\mathcal{L} \left(\mathbb{C} \cdot \mathbf{1}, \tilde{\mathcal{B}}_{\alpha,m'(s+1)-\alpha,m'(s+1)}^{(0)} = \mathcal{B}_{\alpha,0,\frac{m'(s+1)}{\alpha}-1} \right)$, with $P_0^{(k)} := P((if)^{\otimes k})$.

Proof. Due to Lemma 7.11, the continuity at 0 of $t \mapsto P_t \in \mathcal{L}(\mathcal{B}_{\alpha,\beta,\gamma}, \mathcal{B}_{\alpha,\beta,\gamma'})$ will follow from the control of $\|e^{it \cdot f} - 1\|_{\tilde{\mathcal{B}}_{\alpha,s\alpha,a=\alpha(\gamma'-\gamma)}^{(0)}}$ (since $\frac{s\alpha-a+\alpha}{\alpha} \leq s+1 \leq \beta$ and $0 \leq s\alpha \leq a + \alpha\gamma = \alpha\gamma'$).

We conclude by noticing, on the first hand (using $\frac{a}{s+1} \leq 1$), that

$$\begin{aligned} |e^{it \cdot f} - 1| &\leq \min(2, |t| \cdot |f|) \leq \min \left(2, |t| \cdot (|f(x_0)| + C_1 p^{s+1}) \right) \\ &\leq 2^{1-\frac{a}{s+1}} |t|^{\frac{a}{s+1}} (|f(x_0)| + C_1)^{\frac{a}{s+1}} p^{\frac{a}{s+1}(s+1)}, \end{aligned}$$

which implies that $|e^{it \cdot f} - 1|_a^{(0)} = \mathcal{O}(|t|^{\min(\frac{a}{s+1}, \alpha)})$ and, on the second hand, that

$$|e^{it \cdot f(x)} - e^{it \cdot f(y)}| \leq \min(2, |t| \cdot |f(x) - f(y)|) \leq 2^{1-\alpha} |t|^\alpha C_1^\alpha d(x, y)^\alpha (p(x) + p(y))^{s\alpha},$$

ensuring that $m_b^{(0)}(e^{it \cdot f} - 1) = \mathcal{O}(|t|^{\min(\frac{a}{s+1}, \alpha)})$. This ends the proof of the first point. Let us prove the second point. To this end, we first observe that

$$P_t - \sum_{k=0}^m \frac{P_0^{(k)} \cdot t^{\otimes k}}{k!} = P \left(\left(e^{it \cdot f} - \sum_{k=0}^m \frac{(it \cdot f)^k}{k!} \right) \cdot \right).$$

Due to Lemma 7.11, it is enough to prove that $\left\| e^{it \cdot f} - \sum_{k=0}^m \frac{(it \cdot f)^k}{k!} \right\|_{\tilde{\mathcal{B}}_{\alpha,m'(s+1)-\alpha,m'(s+1)}^{(0)}} = \mathcal{O}(|t|^{m'})$.

We set $h(u) := e^{iu} - \sum_{k=0}^m \frac{(iu)^k}{k!}$ and notice that

$$|h(t \cdot f)| \leq 2 \frac{|t \cdot f|^{m'}}{m!} \leq \frac{2|t|^{m'} \left(|f|_{s+1}^{(0)} \right)^{m'}}{m!} p^{(s+1)m'}.$$

Observe that there exists $C'_0 > 0$ such that $|h'(iu)| = \left| i \left(e^{iu} - \sum_{k=0}^{m-1} \frac{(iu)^k}{k!} \right) \right| \leq C'_0 |u|^{m'-1}$ for all $u \in \mathbb{R}$. Thus

$$\begin{aligned} |h(t \cdot f(x)) - h(t \cdot f(y))| &\leq C'_0 \max(|t \cdot f(x)|, |t \cdot f(y)|)^{m'-1} |t \cdot (f(x) - f(y))| \\ &\leq C'_1 |t|^{m'-1} \max(p(x)^{s+1}, p(y)^{s+1})^{m'-1} |t| d(x, y)^\alpha \max(p(x), p(y))^{s+1-\alpha} \\ &\leq |t|^{m'} d(x, y)^\alpha \max(p(x), p(y))^{m'(s+1)-\alpha}. \end{aligned}$$

□

Proof of Theorem 7.9. Proposition 7.10 ensures that assumptions of Theorem 7.1 (except maybe those on ν) are satisfied. Observe that the condition $p \in \mathbb{L}^{(s+1)r}(\nu)$ ensures that $\mathbb{E}_\nu \in \mathcal{B}_r^*$ since $\mathcal{B}_r \hookrightarrow p^{r(s+1)}.L^\infty$.

For Items (i)-(ii) (TCL and LLT), assuming $\varepsilon < 1/3$ and $\varepsilon < \alpha(2s+2)$, we take $\tilde{\mathcal{B}}_0 := \mathcal{B}_{\frac{\varepsilon}{2s+2}, s+1, 2s+1}$ and apply Theorem 3.3 for $\mathcal{B}_{\frac{\varepsilon}{2s+2}, s+1, s+1} \hookrightarrow \tilde{\mathcal{B}}_0 = \mathcal{B}_{\frac{\varepsilon}{2s+2}, s+1, 2s+1}$. This will ensure the continuity of $t \mapsto \Pi_t \in \mathcal{L}(\mathcal{B}_0, \tilde{\mathcal{B}}_0)$. The fact that $\tilde{\mathcal{B}}_0 := \mathcal{B}_{\frac{\varepsilon}{2s+2}, s+1, 2s+1} \hookrightarrow \mathbb{L}^1(\nu)$ will follow from $\tilde{\mathcal{B}}_0 \hookrightarrow p^\varepsilon.L^\infty$ and $p^\varepsilon \in \mathbb{L}^1(\nu)$. Let us check that the assumptions of Theorem 3.3 are satisfied for $\mathcal{B}_{\frac{\varepsilon}{2s+2}, s+1, s+1} \hookrightarrow \tilde{\mathcal{B}}_0$. The continuity assumption of Theorem 3.3 is ensured by Lemma 7.13 (the integrability conditions follows from $\varepsilon < 1/3$). The other assumptions follow from Lemma 7.12. Indeed Conditions (49) and (50) imply (54) for $\beta = \gamma = s+1$ since $\alpha \in]0, \frac{s+1}{s+2}]$ implies $\alpha(2s+2) \leq (s+1)r' + s\alpha$ and $\alpha(s+2) \leq (s+1)(r+1)$; and moreover (54) implies that (54) holds also true with α being replaced by $\varepsilon/(2s+2) < \alpha$ (due to the Hölder inequality). Finally it has been proved in [26, Proposition 11.8] that the non-lattice property implies the exponential decay of $\max_{j=1, \dots, r+1} \sup_{t \in [-\pi, \pi]^d \setminus B(0, \delta_0)} \|P_t^n\|_{\mathcal{B}_j}$ as $n \rightarrow +\infty$. \square

APPENDIX A. MARKOV ADDITIVE PROCESSES

In [24], the authors considered Markov processes $(X_n, \tilde{S}_n)_n$ on $\Omega \times \mathbb{R}^d$ such that

$$\mathbb{E} \left[h(X_n, \tilde{S}_n) \middle| X_{n-1} = x, \tilde{S}_{n-1} = y \right] = \mathbb{E} \left[h(X_n, y + \tilde{S}_n) \middle| X_{n-1} = x, \tilde{S}_{n-1} = 0 \right]. \quad (55)$$

They also considered a continuous version of this form. We concentrate here on the discrete time process. The first coordinate $(X_n)_{n \geq 0}$ of this process is also a Markov process (driving the process $(\tilde{S}_n)_{n \geq 1}$).

We explain here how, starting with the process (X, \tilde{S}) and using an independent sequence $(Z_k)_{k \geq 1}$ of i.i.d. random variables, we can construct a process $(S_n)_{n \geq 0}$ with

$$S_n := \sum_{k=1}^n Y_k, \quad Y_k := f(X_{k-1}, X_k, Z_k)$$

such that $(X_n, \tilde{S}_0 + S_n)_{n \geq 0}$ is Markov with the same transition operator as $(X_n, \tilde{S}_n)_{n \geq 0}$. Let $(Z_n = (Z_n^{(1)}, \dots, Z_n^{(d)}))_{n \geq 1}$ be a sequence of i.i.d. random variables uniformly distributed on $]0, 1[^d$ and independent of $(\tilde{S}_0, (X_n)_{n \geq 0})$. Let $F_k(x, x', s_1, \dots, s_{k-1}, \cdot)$ be the distribution function of $\tilde{S}_1^{(k)} | (X_0 = x, X_1 = x', \tilde{S}_0 = 0, \tilde{S}_1^{(1)} = s_1, \dots, \tilde{S}_1^{(k-1)} = s_{k-1})$ (Jirina's desintegration theorem ensures it is well defined). Consider $f_k(x, x', s_1, \dots, s_{k-1}, \cdot)$ its inverse, i.e.

$$f_k(x, x', s_1, \dots, s_{k-1}, u) = \inf \{ z \in \mathbb{R} | F_k(x, x', s_1, \dots, s_{k-1}, z) \geq u \}.$$

Then we define the coordinates $Y_n^{(k)}$ of $Y_n = (Y_n^{(1)}, \dots, Y_n^{(d)})$ inductively on k by setting, for all $k = 1, \dots, d$,

$$Y_n^{(k)} := f_k(X_{n-1}, X_n, Y_n^{(1)}, \dots, Y_n^{(k-1)}, Z_n^{(k)}).$$

A crucial well known fact is that $f_k(\dots, z) \leq u \Leftrightarrow z \leq F_k(\dots, u)$. In particular

$$Y_n^{(k)} \leq u_k \Leftrightarrow Z_n^{(k)} \leq F_k(X_{n-1}, X_n, Y_n^{(1)}, \dots, Y_n^{(k-1)}, u_k).$$

This ensures that the distribution of $Y_n^{(k)}$ given $(X_{n-1}, X_n, S_{n-1}, (Y_n^{(j)})_{j=1, \dots, k-1})$ is the same as the one of $\tilde{S}_n^{(k)}$ given $(X_{n-1}, X_n, \tilde{S}_{n-1} = 0, (\tilde{S}_n^{(j)} = Y_n^{(j)})_{j=1, \dots, k-1})$. Thus, we prove by induction that, for all $k = 1, \dots, d$, the distribution of $(Y_n^{(1)}, \dots, Y_n^{(k)})$ given X_{n-1}, X_n, S_{n-1} coincides with

the one of $(\tilde{S}_n^{(1)}, \dots, \tilde{S}_n^{(k)})$ given $(X_{n-1}, X_n, \tilde{S}_{n-1} = 0)$. Indeed this holds true for $k = 1$ and, moreover assuming the result holds true at some rank $k \in \{1, \dots, d-1\}$, it follows that

$$\begin{aligned}
& \mathbb{E} \left[\prod_{j=1}^{k+1} \mathbf{1}_{]-\infty, u_j]}(Y_n^{(j)}) \middle| X_{n-1}, X_n, S_{n-1} \right] \\
&= \mathbb{E} \left[\prod_{j=1}^k \mathbf{1}_{]-\infty, u_j]}(Y_n^{(j)}) \mathbf{1}_{\{Z_n^{(k+1)} \leq F_{k+1}(X_{n-1}, X_n, Y_n^{(1)}, \dots, Y_n^{(k)}, u_{k+1})\}} \middle| X_{n-1}, X_n, S_{n-1} \right] \\
&= \mathbb{E} \left[\prod_{j=1}^k \mathbf{1}_{]-\infty, u_j]}(Y_n^{(j)}) F_{k+1}(X_{n-1}, X_n, Y_n^{(1)}, \dots, Y_n^{(k)}, u_{k+1}) \middle| X_{n-1}, X_n, S_{n-1} \right] \\
&= \mathbb{E} \left[\prod_{j=1}^k \mathbf{1}_{]-\infty, u_j]}(\tilde{S}_n^{(j)}) F_{k+1}(X_{n-1}, X_n, \tilde{S}_n^{(1)}, \dots, \tilde{S}_n^{(k)}, u_{k+1}) \middle| X_{n-1}, X_n, \tilde{S}_{n-1} = 0 \right] \\
&= \mathbb{E} \left[\prod_{j=1}^{k+1} \mathbf{1}_{]-\infty, u_j]}(\tilde{S}_n^{(j)}) \middle| X_{n-1}, X_n, S_{n-1} \right],
\end{aligned}$$

and thus the result holds also at rank $k+1$. Therefore the distribution of Y_n given X_{n-1}, X_n, S_{n-1} coincides with the one of \tilde{S}_n given $(X_{n-1}, X_n, \tilde{S}_{n-1} = 0)$. Moreover $(X_n, S_n)_n$ is a Markov process.

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