ABSTRACT<br>Title of dissertation: HIGHER ORDER ASYMPTOTICS FOR THE CENTRAL LIMIT THEOREM AND LARGE DEVIATION PRINCIPLES<br>Buddhima Kasun Fernando Akurugodage Doctor of Philosophy, 2018<br>Dissertation directed by: Professor Dmitry Dolgopyat<br>Department of Mathematics

First, we present results that extend the classical theory of Edgeworth expansions to independent identically distributed non-lattice discrete random variables. We consider sums of independent identically distributed random variables whose distributions have $d+1$ atoms and show that such distributions never admit an Edgeworth expansion of order $d$ but for almost all parameters the Edgeworth expansion of order $d-1$ is valid and the error of the order $d-1$ Edgeworth expansion is typically $\mathcal{O}\left(n^{-d / 2}\right)$ but the $\mathcal{O}\left(n^{-d / 2}\right)$ terms have wild oscillations.

Next, going a step further, we introduce a general theory of Edgeworth expansions for weakly dependent random variables. This gives us higher order asymptotics for the Central Limit Theorem for strongly ergodic Markov chains and for piecewise expanding maps. In addition, alternative versions of asymptotic expansions are introduced in order to estimate errors when the classical expansions fail to hold. As applications, we obtain Local Limit Theorems and a Moderate Deviation Principle.

Finally, we introduce asymptotic expansions for large deviations. For sufficiently regular weakly dependent random variables, we obtain higher order asymptotics (similar to Edgeworth Expansions) for Large Deviation Principles. In particular, we obtain asymptotic expansions for Cramér's classical Large Deviation Principle for independent identically distributed random variables, and for the Large Deviation Principle for strongly ergodic Markov chains.

# HIGHER ORDER ASYMPTOTICS FOR THE CENTRAL LIMIT THEOREM AND LARGE DEVIATION PRINCIPLES 

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Dissertation submitted to the Faculty of the Graduate School of the University of Maryland, College Park in partial fulfillment of the requirements for the degree of<br>Doctor of Philosophy<br>2018

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## Dedication

To the memory of my uncle, Ivan Fernando, who urged me to seek truth.

## Acknowledgments

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## List of Abbreviations

CLT Central Limit Theorem
i.i.d. independent and identically distributed

LCLT Local Central Limit Theorem
LDP Large Deviation Principle
LDCT Lebesgue Dominated Convergence Theorem
LHS Left hand side
LLT Local Limit Theorem
PDE Partial differential equation
RHS Right hand side
WLOG Without loss of generality

## Chapter 1: Introduction

The Central Limit Theorem (CLT) is one of the most fundamental concepts in probability which was introduced by the work of Laplace and Bernoulli. It describes the long term behaviour of random trials repeated under uniform conditions.

Let $S_{N}=\sum_{n=1}^{N} X_{n}$ be a sum of random variables. We say that $S_{N}$ satisfies the CLT if there are real constants $A$ and $\sigma>0$ such that

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \mathbb{P}\left(\frac{S_{N}-N A}{\sqrt{N}} \leq z\right)=\mathfrak{N}(z) \tag{1.1}
\end{equation*}
$$

where $\mathfrak{N}(z)=\int_{-\infty}^{z} \mathfrak{n}(y) d y$ and $\mathfrak{n}(y)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} e^{-\frac{y^{2}}{2 \sigma^{2}}}$.
The usefulness of the CLT and related limit theorems depends on rapid convergence of distributions of normalized partial sums to the limiting distribution. This is because limit theorems are primarily used for approximating distributions of sums of large but finite number of random variables. Therefore, an important problem is to estimate the rate of convergence of (1.1).

In this regard, an asymptotic expansion as a series of increasing powers of order $n^{-1 / 2}$ (now commonly referred to as the Edgeworth expansion) was formally derived by Chebyshev in [8]. Kolmogorov and Gnedenko emphasize the importance of these expansion in their monograph [23] by stating that the Edgeworth Expansion is "the most powerful and general method of finding such corrections."

Definition 1. $S_{N}$ admits Edgeworth expansion of order $r$ if there are polynomials $P_{1}(z), \ldots, P_{r}(z)$ such that

$$
\begin{equation*}
\mathbb{P}\left(\frac{S_{N}-N A}{\sqrt{N}} \leq z\right)=\underbrace{\mathfrak{N}(x)+\sum_{p=1}^{r} \frac{P_{p}(z)}{N^{p / 2}} \mathfrak{n}(z)}_{\mathcal{E}_{r, N}(z)}+o\left(N^{-r / 2}\right) \tag{1.2}
\end{equation*}
$$

uniformly for $z \in \mathbb{R}$.
Remark 1.1. It is an easy observation that Edgeworth expansion of $S_{N}$, if it exists, is unique. Suppose $\left\{P_{p}(z)\right\}_{p}$ and $\left\{\tilde{P}_{p}(z)\right\}_{p}, 1 \leq p \leq r$ are polynomials corresponding to two Edgeworth expansions. Then,

$$
\sum_{p=1}^{r} \frac{P_{p}(z)}{N^{p / 2}} \mathfrak{n}(z)=\sum_{p=1}^{r} \frac{\tilde{P}_{p}(z)}{N^{p / 2}} \mathfrak{n}(z)+o\left(N^{-r / 2}\right)
$$

Multiplying by $\sqrt{N}$ taking the limit $N \rightarrow \infty$ we have $P_{1}(z)=\tilde{P}_{1}(z)$. Therefore,

$$
\sum_{p=2}^{r} \frac{P_{p}(z)}{N^{p / 2}} \mathfrak{n}(z)=\sum_{p=2}^{r} \frac{\tilde{P}_{p}(z)}{N^{p / 2}} \mathfrak{n}(z)+o\left(N^{-r / 2}\right)
$$

Then, multiplying by $N$ and taking $N \rightarrow \infty, P_{2}(z)=\tilde{P}_{2}(z)$. Continuing this $r$ times we can conclude $P_{p}(z)=\tilde{P}_{p}(z)$ for $1 \leq p \leq r$.

Here and in what follows, $A$ is the asymptotic mean i.e. $A=\lim _{N \rightarrow \infty} \mathbb{E}\left(\frac{S_{N}}{N}\right)$.
Work of Lyapunov, Edgeworth and Cramér focus on the problem of finding higher order asymptotics in the CLT. Their main focus was on independent and identically distributed (i.i.d.) sequences of random variables. In 1928, Cramér introduced a theory of Edgeworth expansions for a broad class of random variables. For the first rigorous derivation of this expansion see [10]. The monograph [11] by Cramér also gives a detailed account of his theory of Edgeworth expansions.

Theorem 1.1 (Cramér). Let $X$ be a centred random variable with $\mathbb{E}\left(X^{2}\right)=\sigma^{2}>0$ and $r+2$ absolute moments. Let $X_{1}, \ldots, X_{N}, \ldots$ be sequence of i.i.d. copies of $X$.

Assume further that

$$
\begin{equation*}
\limsup _{|t| \rightarrow \infty}\left|\mathbb{E}\left(e^{i t X}\right)\right|<1 \tag{1.3}
\end{equation*}
$$

Then, $S_{N}$ satisfies (1.2).
Many refinements of this result appear in later literature. A good introduction to this theory and later developments can be found in $[3,11,20,23]$.

In the i.i.d. case, $P_{p}$ 's are polynomials such that the characteristic function $\phi(t)=\mathbb{E}\left(e^{i t X}\right)$ and the Fourier transform $\hat{\mathcal{E}}_{r, N}$ of $\mathcal{E}_{r, N}$ satisfy

$$
\phi\left(\frac{t}{\sigma \sqrt{N}}\right)^{N}-\hat{\mathcal{E}}_{r, N}(t)=o\left(N^{-r / 2}\right) .
$$

For example, $\mathcal{E}_{1, n}(z)=\mathfrak{N}(z)+\mathfrak{n}(z) \frac{\mathbb{E}\left(X^{3}\right)}{6 \sigma^{3} \sqrt{n}}\left(1-z^{2}\right)$ and

$$
\begin{aligned}
\mathcal{E}_{2, n}(z)=\mathfrak{N}(z)+\mathfrak{n}(z)\left[\frac{\mathbb{E}\left(X^{3}\right)}{6 \sqrt{n} \sigma^{3}}\left(1-z^{2}\right)+\right. & \frac{\mathbb{E}\left(X^{4}\right)-3 \sigma^{4}}{24 n \sigma^{4}}\left(3 z-z^{3}\right) \\
& \left.-\frac{\mathbb{E}\left(X^{3}\right)^{2}}{72 n \sigma^{6}}\left(15 z-10 z^{3}+z^{5}\right)\right] .
\end{aligned}
$$

Since all distributions with an absolutely continuous component satisfy (1.3), this theorem covers a large class of random variables. However, (1.3) indicates that the common distribution of $X_{n}$ 's is far from being discrete. In fact, (1.3) fails when random variables are purely discrete. Surprisingly, not much had been explored in the case of discrete random variables, except in the lattice case. The purpose of my first project [16], joint with Dmitry Dolgopyat, was to address this issue. A detailed discussion about this can found in Chapter 2.

When $X_{n}$ 's are i.i.d., it is known that the order 1 Edgeworth expansion exists if and only if the distribution is non-lattice (see [19]). Therefore, the following asymptotic expansion for the Local Central Limit Theorem (LCLT) for lattice random
variables is also useful.

Definition 2. Suppose that $X_{n}$ 's are integer valued. We say that $S_{N}$ admits a lattice Edgeworth expansion of order $r$, if there are polynomials $P_{0, d}, \ldots, P_{r, d}$ and a number $A$ such that

$$
\sqrt{N} \mathbb{P}\left(S_{N}=k\right)=\mathfrak{n}\left(\frac{k-N A}{\sqrt{N}}\right) \sum_{p=0}^{r} \frac{P_{p, d}((k-N A) / \sqrt{N})}{N^{p / 2}}+o\left(N^{-r / 2}\right)
$$

uniformly for $k \in \mathbb{Z}$.
Remark 1.2. Here, the subscript $d$ in $P_{p, d}$ refers to the fact that the expansion is for discrete lattice-valued random variables. A priori, there is no reason for the polynomials $P_{p}$ in Definition 1 to be related to $P_{p, d}$. In Section 3.3, we show how these two polynomials are related.

As in remark 1.1, we can prove the uniqueness of this expansion. Because $P_{p, d}$ 's have finite degree, say at most $q$, choose $N$ large enough so that $S_{N}$ has more than $q$ values. Then the argument in remark 1.1 applies.

During the 20th century, the work of Lyapunov, Edgeworth, Cramér, Kolmogorov, Esséen, Petrov, Bhattacharya and many others led to the development of the theory of these two asymptotic expansions. See $[26,31]$ and references therein, for more details.

It is immediate that $S_{N}$ admits an order $r$ Edgeworth expansion if

$$
\begin{equation*}
\lim _{N \rightarrow \infty} N^{r / 2}\left[\mathbb{P}\left(\frac{S_{N}-N A}{\sqrt{N}} \leq z\right)-\mathcal{E}_{r, N}(z)\right]=0 \tag{1.4}
\end{equation*}
$$

uniformly in $z$. [3, 4] discuss weak Edgeworth expansions where the LHS of (1.4) is convolved with smooth compactly supported functions. These expansions yield the asymptotics of $\mathbb{E}\left(f\left(S_{N}\right)\right)$.

To introduce these expansions, suppose $(\mathcal{F},\|\cdot\|)$ is a function space.
Definition 3. $S_{N}$ admits weak global Edgeworth expansion of order $r$ for $f \in \mathcal{F}$ if there are polynomials $P_{0, g}(z), \ldots P_{r, g}(z)$ and $A$ (which are independent of $f$ ) such that

$$
\mathbb{E}\left(f\left(S_{N}-N A\right)\right)=\sum_{p=0}^{r} \frac{1}{N^{\frac{p}{2}}} \int P_{p, g}(z) \mathfrak{n}(z) f(z \sqrt{N}) d z+\|f\| \cdot o\left(N^{-(r+1) / 2}\right)
$$

Definition 4. $S_{N}$ admits weak local Edgeworth expansion of order $r$ for $f \in \mathcal{F}$ if there are polynomials $P_{0, l}(z), \ldots P_{r, l}(z)$ and $A($ which are independent of $f)$ such that

$$
\sqrt{N} \mathbb{E}\left(f\left(S_{N}-N A\right)\right)=\frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p}} \int P_{p, l}(z) f(z) d z+\|f\| \cdot o\left(N^{-r / 2}\right)
$$

We also introduce the following asymptotic expansion which yields an averaged form of the error of approximation.

Definition 5. $S_{N}$ admits averaged Edgeworth expansion of order $r$ if there are polynomials $P_{1, a}(z), \ldots P_{r, a}(z)$ and numbers $k, m$ such that for $f \in \mathcal{F}$ we have

$$
\begin{aligned}
& \int\left[\mathbb{P}\left(\frac{S_{N}-N A}{\sqrt{N}} \leq z+\frac{y}{\sqrt{N}}\right)-\mathfrak{N}\left(z+\frac{y}{\sqrt{N}}\right)\right] f(y) d y \\
& \quad=\sum_{p=1}^{r} \frac{1}{N^{p / 2}} \int P_{p, a}\left(z+\frac{y}{\sqrt{N}}\right) \mathfrak{n}\left(z+\frac{y}{\sqrt{N}}\right) f(y) d y+\|f\| \cdot o\left(N^{-r / 2}\right)
\end{aligned}
$$

Remark 1.3. Here, the subscripts $g, l$, a refer to global, local and averaged respectively and used to distinguish the polynomials appearing each definition. In Section 3.3, we show how these two polynomials are related.

All of these weak forms of expansions are unique provided that $\mathcal{F}$ is dense in $C_{c}^{\infty}$. If there are two different weak global expansions with polynomials $\left\{P_{p, g}\right\}$ and
$\left\{\tilde{P}_{p, g}\right\}$, the argument in remark 1.1 yields,

$$
\int P_{p, g}(z) \mathfrak{n}(z) f(z \sqrt{N}) d z=\int \tilde{P}_{p, g}(z) \mathfrak{n}(z) f(z \sqrt{N}) d z
$$

for all $f \in C_{c}^{\infty}$ which gives us the equality, $P_{p, g}(z)=\tilde{P}_{p, g}(z)$. The same idea works for the other two expansions.

We have seen that these asymptotic expansions are unique. They also form a hierarchy. We discuss this in Appendix A.2. Due to this hierarchy, in the absence of one, others can be useful in extracting information about the rate of convergence in (1.1).

Previous results on existence of Edgeworth expansions, for example in [11, 20, 23], assume independence of random variables $X_{n}$. For many applications the independence assumption of random variables is too restrictive. Because of this reason, there have been attempts to develop a theory of Edgeworth expansions for weakly dependent random variables where weak dependence often refers to asymptotic decorrelation. See $[9,22,29,40,41]$ for such examples. Their focus is on the classical expansions introduced in Definition 1 and Definition 2.

Except in [9], the sequences of random variables considered are uniformly ergodic Markov processes with strong recurrent properties or processes approximated by such Markov processes. In [9], the authors consider aperiodic subshifts of finite type endowed with a stationary equilibrium state and give explicit construction of the order 1 Edgeworth expansion. They also prove the existence of higher order classical Edgeworth expansions under a rapid decay assumption on the tail of the characteristic function.

The goal of [21], a joint work with Carlangelo Liverani, is to generalize these results and to provide sufficient conditions that guarantee the existence of Edgeworth expansions for weakly dependent random variables including observations arising from sufficiently chaotic dynamical systems, and strongly ergodic Markov chains. In fact, we introduce a widely applicable theory for both classical and weak forms of Edgeworth expansions and significantly improve existing results. This work is discussed in detail in Chapter 3.

The CLT and related asymptotic expansions provide accurate descriptions only of typical events. For example, if $X_{n}$ 's are centered i.i.d. random variables then for all $a>0, \lim _{N \rightarrow \infty} \mathbb{P}\left(S_{N} \geq a N\right)=0$, due to the Law of Large Numbers i.e. $\frac{S_{N}}{N} \rightarrow 0$ in probability. Large Deviation Principles (LDPs) give better descriptions of these non-typical events by specifying the exponential rate at which their probabilities decay.

Before we present results related to LDPs, we recall the following definitions, and facts whose proofs can be found in $[17,30]$. Given a function $f: \mathbb{R} \rightarrow(-\infty, \infty]$ with $f \neq \infty$, define its effective domain to be $D_{f}=\{x \in \mathbb{R} \mid f(x)<\infty\}$ and its Legendre transform by $f^{*}(x)=\sup _{t \in \mathbb{R}}[t x-f(t)]$. Then, $f^{*}$ is convex and lower semi-continuous. Therefore, $D_{f^{*}}$ is an interval and $f^{*}$ is continuous on $\bar{D}_{f^{*}}$.

In addition, suppose $f$ is convex, lower semi-continuous with $\stackrel{\circ}{D}_{f}=(a, b)$ and $f \in C^{2}(a, b)$ with $f^{\prime \prime}>0$ on $(a, b)$ (possibly $a=-\infty$ or $b=+\infty$ ). Then, ${\stackrel{\circ}{f^{*}}}=$ $(A, B)$ where $A=\lim _{t \rightarrow a+} f^{\prime}(t)$ and $B=\lim _{t \rightarrow b-} f^{\prime}(t), f^{*}$ is continuously differentiable on $(A, B)$ and $\left(f^{*}\right)^{\prime}=\left(f^{\prime}\right)^{-1}$. For any $f$ satisfying the above properties, for any $x \in \stackrel{\circ}{D}_{f^{*}}$ the supremum in the definition of $f^{*}(x)$ is achieved at the unique point
$t \in \check{D}_{f}$ which solves $f^{\prime}(t)=x$ and hence, $f^{*}(x)=\sup _{t \in D_{f}}[t x-f(t)]$. Also, $f$ is called steep if $\lim _{t \rightarrow a}\left|f^{\prime}(t)\right|=\lim _{t \rightarrow b}\left|f^{\prime}(t)\right|=\infty$.

The following classical result, due to Cramér, is one of the fundamental results in the theory of Large Deviations.

Theorem 1.2 (Cramér). Let $X$ be a real valued random variable with mean $A$ and variance $\sigma^{2}>0$. Suppose that the logarithmic moment generating function of $X$, $\log \mathbb{E}\left(e^{t X}\right)$, is finite in a neighbourhood of 0 . Let $X_{n}$ be a sequence of i.i.d. copies of X. Then,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{P}\left(S_{N} \geq N z\right)=-I(z), \text { if } z>A
$$

and

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{P}\left(S_{N} \leq N z\right)=-I(z), \text { if } z<A
$$

where $I$ is given by $I(z)=\sup _{\lambda \in \mathbb{R}}\left[\lambda z-\log \mathbb{E}\left(e^{\lambda X}\right)\right]$ (the Legendre transform of the logarithmic moment generating function of $X$ ).

From the hypothesis it is immediate that $I$ is convex and lower semi-continuous. Also, $I^{\prime \prime}>0$ on $\grave{D}_{I}=(\inf (\operatorname{supp} X), \sup (\operatorname{supp} X))$, therefore $I$ is strictly convex on $\stackrel{\circ}{D}_{I}, I(z)=0 \Longleftrightarrow z=\mu$ and there is a unique $\lambda^{*}$ such that $I(z)=$ $\lambda^{*} z-\log \mathbb{E}\left(e^{\lambda^{*} X}\right)$.

Cramér's LDP has an extension to the non-i.i.d. case. We refer the reader to [6][Chapter V.6] for a proof of the following result.

Theorem 1.3 (Local Gärtner-Ellis). Let $X_{n}$ be a sequence of random variables not necessarily i.i.d. Suppose there exists $\delta>0$ such that for $\lambda \in(-\delta, \delta)$,

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{E}\left(e^{\lambda S_{N}}\right)=\Omega(\lambda) \tag{1.5}
\end{equation*}
$$

where $\Omega$ is strictly convex continuously differentiable function with $\Omega^{\prime}(0)=0$. Then, for all $z \in\left(0, \frac{\Omega(\delta)}{\delta}\right)$,

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{P}\left(S_{N} \geq N z\right)=-I(z) \tag{1.6}
\end{equation*}
$$

where $I(z)=\sup _{\lambda \in(-\delta, \delta)}[z \lambda-\Omega(\lambda)]$.

## Remark 1.4.

1. If the limit (1.5) exists for all $\lambda \in \mathbb{R}$. Then, $B=\lim _{\delta \rightarrow \infty} \frac{\Omega(\delta)}{\delta}$ exists and (1.6) holds for all $z \in(0, B)$.
2. The function I appearing in the theorem is called the rate function because it gives us the exponential rate of decay of tail probabilities.

In an on-going joint work with Pratima Hebbar, we develop a theory of higher order asymptotics for LDPs, using the weak forms of Edgeworth expansions and extensions of results in [27, Chapter VIII]. As in the CLT case, higher order asymptotics are given as expansions.

Definition 6. Suppose $S_{N}$ satisfies an LDP with rate function I. Then, $S_{N}$ admits strong asymptotic expansion of order $r$ for large deviations in the range $(0, L)$ if there are functions $C_{p}:(0, L) \rightarrow \mathbb{R}$ for $0 \leq p<\frac{r}{2}$ and $A>0$ such that for each $a \in(0, L)$,

$$
\mathbb{P}\left(S_{N}-A N \geq a N\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{C_{p}(a)}{N^{p+1 / 2}}+C_{r, a} \cdot o\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

These expansions are in the spirit of the higher order expansions found [1] for i.i.d. sequences of random variables. In [7], authors refer to these expansions as strong large deviation results. $[7,32]$ establish the order 1 expansions under certain assumptions on the behaviour of the moment generating functions. These strengthen
the results of [1] but only in the order 1 case. Here, we give an alternative way to establish the so-called strong large deviation results of all orders. We also manage to recover the results in [1] in the non-lattice setting. For applications of these results to statisitcs, see examples listed in $[1,7,32]$ and references therein.

We also introduce the following weak form of the expansion for LDPs. As in the CLT case, we define these expansions over a function space $(\mathcal{F},\|\cdot\|)$.

Definition 7. Suppose $S_{N}$ satisfies an LDP with rate function $I$. Then, $S_{N}$ admits weak asymptotic expansion of order $r$ for large deviations in the range $(0, L)$ for $f \in \mathcal{F}$, if there are functions $D_{p}:(0, L) \rightarrow \mathbb{R}$ for $0 \leq p<\frac{r}{2}$ and $A>0$ such that for each $a \in(0, L)$,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{D_{p}(a)}{N^{p+1 / 2}}+C_{r, a} \cdot o\left(\frac{1}{N^{\frac{r+1}{2}}}\right) .
$$

In fact, our results show that for a sequence $X_{n}$ of i.i.d. $l$-Diophantine random variables with all exponential moments, for every $r, S_{N}$ admits weak asymptotic expansions of order $r$ for large deviations on $(0, \infty)$ for sufficiently regular $f$. This is a refinement of the LDP by Cramér for a broad class of random variables.

We also obtain similar results for certain classes of non-i.i.d. random variables. As an application, we obtain asymptotic expansions for the LDP in the case of Markov chains with smooth densities. In particular, let $x_{n}$ be a time homogeneous Markov chain on a compact connected manifold $\mathcal{M}$ with a smooth transition density and $h: \mathcal{M} \times \mathcal{M} \rightarrow \mathbb{R}$ be smooth with non-degenerate critical points. Then $X_{n}=$ $h\left(x_{n}, x_{n-1}\right)$ admits asymptotic expansions for large deviations of all orders. These results are presented in Chapter 4.

## Chapter 2: Central Limit Theorem: Discrete Random Variables.

### 2.1 Overview and main results.

Let $X$ be a random variable with zero mean and positive variance $\sigma^{2}$. Let $S_{n}=\sum_{n=1}^{n} X_{j}$ where $X_{j}$ are independent identically distributed and have the same distribution as $X$. Then, it is well-known that $S_{n}$ satisfies the CLT with $A=0$ and $\sigma$ as in (1.1).

In this chapter, we consider a case which is opposite to $X$ having a density, namely we suppose that $X$ has a discrete distribution with $d+1$ atoms where $d \geq 2$. $d=2$ is the simplest non-trivial case since distributions with two atoms are lattice and as a result they do not admit even the first order Edgeworth expansion.

Thus, we suppose that $X$ takes values $a_{1}, \ldots, a_{d+1}$ with probabilities $p_{1}, \ldots, p_{d+1}$ respectively. Since $X$ should have zero mean we suppose that our $2(d+1)$-tuple $(\mathbf{a}, \mathbf{p})$ belongs to the set

$$
\Omega=\left\{p_{i}>0, \quad p_{1}+\cdots+p_{d+1}=1, \quad p_{1} a_{1}+\cdots+p_{d+1} a_{d+1}=0\right\} .
$$

It is easy to see that $S_{n}$ never admits the order $d$ Edgeworth expansion. Indeed

$$
\begin{equation*}
\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(S_{n} \leq z\right)=\sum_{\substack{m_{i} \geq 0, \sum m_{i}=n \\ \sum_{i} m_{i} a_{i} \leq z}} \frac{n!}{m_{1}!\ldots m_{d+1}!} p_{1}^{m_{1}} \ldots p_{d+1}^{m_{d+1}} \tag{2.1}
\end{equation*}
$$

Applying the Local Central Limit Theorem to the time homogeneous $\mathbb{Z}^{d}$-random walk which jumps to $\mathbf{e}_{i}$ from the origin $\mathbf{0}$ with probability $p_{i}$ for $i=1, \ldots, d$ and stays at $\mathbf{0}$ with probability $p_{d+1}$ we conclude that if

$$
\sum m_{i} a_{i}=n \sum a_{i} p_{i}+\mathcal{O}(\sqrt{n})
$$

then

$$
n^{d / 2} \frac{n!}{m_{1}!\ldots m_{d+1}!} p_{1}^{m_{1}} \ldots p_{d+1}^{m_{d+1}}
$$

is uniformly bounded from below. Accordingly $\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(S_{n} \leq z\right)$ has jumps of order $n^{-d / 2}$. On the other hand $\mathcal{E}_{d}(z)$ is a smooth function of $z$. So, it cannot approximate both $\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(S_{n} \leq z-0\right)$ and $\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(S_{n} \leq z+0\right)$ at the points of jumps.

Here we show that for typical $(\mathbf{a}, \mathbf{p})$ the order $d$ Edgeworth expansion just barely fails. We present two results in this direction. For the first result let

$$
b_{j}=a_{j}-a_{1}, \text { for } j=2 \ldots d+1
$$

Set

$$
d(s)=\max _{j \in\{2, \ldots d+1\}} \operatorname{dist}\left(b_{j} s, 2 \pi \mathbb{Z}\right)
$$

We say that $\mathbf{a}$ is $\beta$-Diophantine if there is a constant $K$ such that for $|s|>1$,

$$
d(s) \geq \frac{K}{|s|^{\beta}} .
$$

It is easy to see that almost all $\mathbf{a}$ is $\beta$-Diophantine provided that $\beta>(d-1)^{-1}$ (see $[36,47])$.

Theorem 2.1.1. If $\mathbf{a}$ is $\beta$-Diophantine and

$$
\begin{equation*}
2\left(R-\frac{1}{2}\right) \beta<1 \tag{2.2}
\end{equation*}
$$

then

$$
\lim _{n \rightarrow \infty} n^{R}\left[\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)-\mathcal{E}_{d-1}(z)\right]=0
$$

Thus for almost every a the order $(d-1)$ Edgeworth expansion approximates the distribution of $\frac{S_{n}}{\sigma \sqrt{n}}$ with error $\mathcal{O}\left(n^{\varepsilon-d / 2}\right)$ for any $\varepsilon$.

Note that Theorem 2.1.1 applies for all $\beta \mathrm{s}$, in particular for $\beta \mathrm{s}$ which are much larger than $(d-1)^{-1}$. However if $\beta$ is large, then the statement of the theorem can be simplified. Namely, let $r$ be the integer such that $r<2 R \leq r+1$. (Note that (2.2) can be rewritten as $2 R<\frac{1}{\beta}+1$ so provided that $2 R$ is suffciently close to $\frac{1}{\beta}+1$ we can take $r=\left\lfloor\beta^{-1}\right\rfloor+1$. Then,

$$
\begin{aligned}
\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right) & =\mathcal{E}_{d-1}(z)+o\left(\frac{1}{n^{R}}\right) \\
& =\mathcal{E}_{r}(z)+o\left(\frac{1}{n^{R}}\right)+\mathcal{O}\left(\mathcal{E}_{d-1}(z)-\mathcal{E}_{r}(z)\right) .
\end{aligned}
$$

Since $\frac{r+1}{2}>R$ the first error term dominates the second and we obtain the following result.

## Corollary 2.1.1.

$$
\lim _{n \rightarrow \infty} n^{R}\left[\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)-\mathcal{E}_{r}(z)\right]=0
$$

provided that $\mathbf{a}$ is $\beta$-Diophantine, $r=1+\left\lfloor\beta^{-1}\right\rfloor$, and $r<2 R<\frac{1}{\beta}+1$.
Theorem 2.1.1 shows that for almost every a and for $r \in\{1, \ldots, d-1\}$, the order $r$ Edgeworth expansion is valid. Results that follow show that,

$$
\begin{equation*}
\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)-\mathcal{E}_{d}(z) \tag{2.3}
\end{equation*}
$$

is typically of order $\mathcal{O}\left(n^{-d / 2}\right)$ but the $\mathcal{O}\left(n^{-d / 2}\right)$ term has wild oscillations. To formulate this result precisely we suppose that our $2(d+1)$-tuple is chosen at random
according to an absolutely continuous distribution $\mathbf{P}$ on $\Omega$. Thus (2.3) becomes a random variable.

Theorem 2.1.2. There exists a smooth function $\Lambda(\mathbf{a}, \mathbf{p})$ such that for each $z$ the random variable

$$
e^{z^{2} / 2} \frac{n^{d / 2}}{\Lambda(\mathbf{a}, \mathbf{p})}\left[\mathcal{E}_{d}(z)-\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)\right]
$$

converges in law to a non-trivial random variable $\mathcal{X}$.

More precisely we have,

$$
\begin{equation*}
\Lambda(\mathbf{a}, \mathbf{p})=\frac{\left|a_{d+1}-a_{1}\right|}{2^{d} \pi^{d+\frac{1}{2}} \sqrt{\operatorname{det}\left(D_{\mathbf{a}, \mathbf{p}}\right)} \sigma(\mathbf{a}, \mathbf{p})} \tag{2.4}
\end{equation*}
$$

where $D_{\mathbf{a}, \mathbf{p}}$ is a $(d-1) \times(d-1)$ matrix defined by equations (2.37)-(2.38) of Section 2.5, $\sigma(\mathbf{a}, \mathbf{p})$ denotes the standard deviation of the distribution of the random variable taking value $a_{j}$ with probability $p_{j}$ and $\mathcal{X}$ is defined as follows.

Let $\mathbb{M}$ be the space of pairs $(\mathcal{L}, \chi)$ where $\mathcal{L}$ is a unimodular lattice in $\mathbb{R}^{d}$ and $\chi$ is a homeomorphism $\chi: \mathcal{L} \rightarrow \mathbb{T}$. In the formulas below, we identify $\mathbb{T}$ with segment $[0,1)$ equipped with addition modulo one. Given a vector $\mathbf{w} \in \mathbb{R}^{d}$ we denote by $y(\mathbf{w})$ its first coordinate and by $\mathbf{x}(\mathbf{w})$ its last $d-1$ coordinates.

Lemma 2.1.2. For almost every pair $(\mathcal{L}, \chi) \in \mathbb{M}$ with respect to the Haar measure the following limit exists

$$
\begin{equation*}
\mathcal{X}(\mathcal{L}, \chi)=\lim _{R \rightarrow \infty} \sum_{\mathbf{w} \in \mathcal{L} \backslash\{0\},\|\mathbf{w}\| \leq R} \frac{\sin (2 \pi \chi(\mathbf{w}))}{y(\mathbf{w})} e^{-\|\mathbf{x}(\mathbf{w})\|^{2}} \tag{2.5}
\end{equation*}
$$

In order to simplify the notation we will abbreviate expressions such as (2.5)
by

$$
\begin{equation*}
\mathcal{X}(\mathcal{L}, \chi)=\sum_{\mathbf{w} \in \mathcal{L} \backslash\{0\}} \frac{\sin (2 \pi \chi(\mathbf{w}))}{y(\mathbf{w})} e^{-\|\mathbf{x}(\mathbf{w})\|^{2}} \tag{2.6}
\end{equation*}
$$

The Haar measure on $\mathbb{M}$ can be defined in two equivalent ways. First, note that $\chi$ is of the form $\chi(\mathbf{w})=e^{i \tilde{\chi}(\mathbf{w})}$ for some linear functional $\tilde{\chi} \in\left(\mathbb{R}^{d}\right)^{*} . S L_{d}(\mathbb{R})$ acts on $\mathbb{R}^{d} \oplus\left(\mathbb{R}^{d}\right)^{*}$ by the formula,

$$
A(\mathbf{w}, \tilde{\chi})=\left(A \mathbf{w}, \tilde{\chi} A^{-1}\right)
$$

Observe that if $A(\mathbf{w}, \tilde{\chi})=(\hat{\mathbf{w}}, \hat{\chi})$ then,

$$
\begin{equation*}
\tilde{\chi}(\mathbf{w})=\hat{\mathbf{w}}(\hat{\chi}) \tag{2.7}
\end{equation*}
$$

The above action of $S L_{d}(\mathbb{R})$ induces the following action of $S L_{d}(\mathbb{R}) \ltimes\left(\mathbb{R}^{d}\right)^{*}$ on $\mathbb{M}$ given by,

$$
(A, \tilde{\chi})(\mathcal{L}, \chi)=\left(A \mathcal{L}, e^{2 \pi i t \tilde{\chi}} \cdot\left(\chi \circ A^{-1}\right)\right)
$$

This action is transitive because $S L_{d}(\mathbb{R})$ acts transitively on unimodular lattices and $\left(\mathbb{R}^{d}\right)^{*}$ acts transitively on characters. This allows us to identify $\mathbb{M}$ with $\left(S L_{d}(\mathbb{R}) \ltimes\right.$ $\left.\mathbb{R}^{d}\right) /\left(S L_{d}(\mathbb{Z}) \ltimes \mathbb{Z}^{d}\right)$ and so $\mathbb{M}$ inherits the Haar measure from $S L_{d}(\mathbb{R}) \ltimes \mathbb{R}^{d}$.

The second way to define the Haar measure is to note that the space $\mathcal{M}$ of unimodular lattices is naturally identified with $S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z})$ and so it inherits the Haar measure from $S L_{d}(\mathbb{R})$. Next for a fixed $\mathcal{L}$ the set of homeomorphisms $\chi: \mathcal{L} \rightarrow \mathbb{T}$ is a $d$ dimensional torus so it comes with its own Haar measure.

Now, if we want to compute the average of a function $\Phi(\mathcal{L}, \chi)$ with respect to the Haar measure then we can first compute its average $\bar{\Phi}(\mathcal{L})$ in each fiber and then integrate the result with respect to the Haar measure on the space of lattices. In the proof of Lemma 2.1.2 given in Section A. 1 the averaging inside a fiber will be denoted by $\mathbf{E}_{\chi}$ and the averaging with respect to the Haar measure on the space of lattices will be denoted by $\mathbf{E}_{\mathcal{L}}$.

If we assume that the pair $(\mathcal{L}, \chi)$ is distributed according to the Haar measure on $\mathbb{M}$ then $\mathcal{X}$, defined in Lemma 2.1.2, becomes a random variable. This is the variable mentioned in Theorem 2.1.2. Note that the distribution of $\mathcal{X}$ depends neither on $\mathbf{P}$ nor on $z$.

Using the second representation of the Haar measure we can also describe $\mathcal{X}$ as follows. Let $\mathbf{w}_{1}, \ldots, \mathbf{w}_{d}$ be the shortest spanning set of $\mathcal{L}$. That is $\mathbf{w}_{1}$ is the shortest non zero vector in $\mathcal{L}$ and, for $j>1, \mathbf{w}_{j}$ is the shortest vector which is linearly independent of $\mathbf{w}_{1}, \ldots, \mathbf{w}_{j-1}$. Given $\mathbf{m}=\left(m_{1}, \ldots, m_{d}\right) \in \mathbb{Z}^{d}$ let $(y, \mathbf{x})(\mathbf{m})$, $y \in \mathbb{R}$ and $\mathbf{x} \in \mathbb{R}^{d-1}$, denote the point

$$
\begin{equation*}
m_{1} \mathbf{w}_{1}+\cdots+m_{d} \mathbf{w}_{d} \in \mathcal{L} \tag{2.8}
\end{equation*}
$$

Let $\theta_{j}=\chi\left(\mathbf{w}_{j}\right)$. Then $\theta_{j}$ are uniformly distributed on $\mathbb{T}$ and independent of each other. Set $\theta(\mathbf{m})=m_{1} \theta_{1}+\cdots+m_{d} \theta_{d}$. Now $\mathcal{X}$ (see definition in Lemma 2.1.2) can be rewritten as

$$
\begin{equation*}
\mathcal{X}=\sum_{\mathbf{m} \in \mathbb{Z}^{d} \backslash\{\mathbf{0}\}} \frac{\sin (2 \pi \theta(\mathbf{m}))}{y(\mathbf{m})} e^{-\|\mathbf{x}(\mathbf{m})\|^{2}} \tag{2.9}
\end{equation*}
$$

where $\mathcal{L}$ is uniformly distributed on the space of lattices, $(y, \mathbf{x})(\mathbf{m})$ is defined by (2.8), and $\left(\theta_{1}, \ldots \theta_{d}\right)$ is uniformly distributed on $\mathbb{T}^{d}$ and independent of $\mathcal{L}$.

Theorems 2.1.1 and 2.1.2 have analogues when we consider probabilities that $S_{n}$ belongs to finite intervals. In particular, our results have applications to the Local Limit Theorem.

Theorem 2.1.3. Let $z_{1}(n)$ and $z_{2}(n)$ be two uniformly bounded sequences such that $\left|z_{1}(n)-z_{2}(n)\right| n^{d / 2} \rightarrow \infty$. Then, the random vector,

$$
\begin{equation*}
\frac{n^{d / 2}}{\Lambda(\mathbf{a}, \mathbf{p})}\left(e^{z_{1}^{2} / 2}\left[\mathcal{E}_{d}\left(z_{1}\right)-\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z_{1}\right)\right], e^{z_{2}^{2} / 2}\left[\mathcal{E}_{d}\left(z_{2}\right)-\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z_{2}\right)\right]\right) \tag{2.10}
\end{equation*}
$$

converges in law to a random vector $\left(\mathcal{X}\left(\mathcal{L}, \chi_{1}\right), \mathcal{X}\left(\mathcal{L}, \chi_{2}\right)\right)$ where $\mathcal{X}(\mathcal{L}, \chi)$ is defined by (2.6) and the triple $\left(\mathcal{L}, \chi_{1}, \chi_{2}\right)$ is uniformly distributed on $\left(S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z})\right) \times$ $\mathbb{T}^{d} \times \mathbb{T}^{d}$.

Here and below the uniform distribution of $\left(\mathcal{L}, \chi_{1}, \chi_{2}\right)$ means that $\mathcal{L}$ is uniformly distributed on the space of lattices and for a given lattice, $\chi_{1}$ and $\chi_{2}$ are chosen independently and uniformly from the space of characters.

Theorem 2.1.4. Let $z_{1}(n)<z_{2}(n)$ be two uniformly bounded sequences such that $l_{n}=z_{2}(n)-z_{1}(n) \rightarrow 0$.
(a) If $l_{n} \geq C n^{\varepsilon-d / 2}$ for some $\varepsilon>0$ then

$$
\frac{\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(z_{1}<\frac{S_{n}}{\sigma \sqrt{n}}<z_{2}\right)}{l_{n} \mathfrak{n}\left(z_{1}\right)} \rightarrow 1 \text { almost surely. }
$$

(b) If $l_{n} n^{d / 2} \rightarrow \infty$ then

$$
\frac{\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(z_{1}<\frac{S_{n}}{\sigma \sqrt{n}}<z_{2}\right)}{l_{n} \mathfrak{n}\left(z_{1}\right)} \Rightarrow 1
$$

(here and below " $\Rightarrow$ " denotes the convergence in law).
(c) If $l_{n}=\frac{c\left|a_{d+1}-a_{1}\right|}{\sigma(\mathbf{a}, \mathbf{p}) n^{d / 2}}$ then

$$
2^{d-\frac{3}{2}} \pi^{d} \sqrt{\operatorname{det}\left(D_{\mathbf{a}, \mathbf{p}}\right)}\left[\frac{\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(z_{1}<\frac{S_{n}}{\sigma \sqrt{n}}<z_{2}\right)}{l_{n} \mathfrak{n}\left(z_{1}\right)}-1\right] \Rightarrow \mathcal{Y}
$$

where

$$
\mathcal{Y}(\mathcal{L}, \chi, c)=\sum_{\mathbf{w} \in \mathcal{L} \backslash\{\mathbf{0}\}} \frac{\sin (2 \pi[\chi(\mathbf{w})-c y(\mathbf{w})])-\sin (2 \pi \chi(\mathbf{w}))}{y(\mathbf{w})} e^{-\|\mathbf{x}(\mathbf{w})\|^{2}}
$$

and $\mathcal{L}, \chi$ are as in Theorem 2.1.2 and $D_{\mathbf{a}, \mathbf{p}}$ given by equations (2.37)-(2.38).
The intuition behind this result is the following. Call $y_{n} \delta$-plausible if $\mathbb{P}\left(S_{n}=\right.$ $\left.y_{n}\right) \geq \delta n^{-d / 2}$. The discussion following (2.1) shows that for each $\delta$ there are about
$C(\delta) n^{d / 2} \delta$-plausible values. Therefore, if $l_{n} \ll n^{-d / 2}$ then the interval $\left[z_{1}(n), z_{2}(n)\right]$ would typically contain no plausible values. Hence, we should not expect the LLT to hold on that scale. Theorem 2.1.4 shows that as soon as interval $\left[z_{1}(n), z_{2}(n)\right]$ contains many plausible values then the LLT typically holds for this interval.

Recall that,

$$
\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(S_{n} \in\left[z_{1}, z_{2}\right]\right)=\sum_{\substack{m_{i} \geq 0, \sum_{m_{i}=n}=n \\ z_{1} \leq \Sigma m_{i} a_{i} \leq z_{2}}} \frac{n!}{m_{1}!\ldots m_{d+1}!} p_{1}^{m_{1}} \ldots p_{d+1}^{m_{d+1}}
$$

Thus, in Theorem 2.1.4 we just count the number of visits of a random linear form $\sum m_{i} a_{i}$ to a finite interval with weights given by multinomial coefficients. It is also interesting to consider counting with equal weight. In this case the analogue of Theorem 2.1.4(c) is obtained in [38] while for longer intervals only partial results are available, for example see $[15,34]$.

The chapter is organized as follows. Theorem 2.1.1 is proven in Section 2.2. The proof is a minor modification of the arguments of [20, Chapter XVI]. The bulk of the chapter is devoted to the proof of Theorem 2.1.2. In Section 2.3 we provide an equivalent formula for $\mathcal{X}$. This formula looks more complicated than (2.6) but it is easier to identify with the limit of the error term. Section 2.4 contains preliminary reductions. We show that the density $\rho$ on $\Omega$ could be assumed smooth and the integration in the Fourier inversion formula could be restricted to a finite domain. In Section 2.5, we show that main contribution to the error term comes from resonances where characteristic function of $S_{n}$ is close to 1 in absolute value. The proof relies on several technical estimates which are established in Section 2.6. In Section 2.7, we use dynamics on homogenuous spaces in order to show that the contribution of
resonances converges to (2.6) completing the proof of Theorem 2.1.2. The proofs of Theorems 2.1.3 and 2.1.4 are similar to the proof of Theorem 2.1.2. The necessary modifications are explained in Section 2.8. We postpone the proof of Lemma 2.1.2 till Appendix A.1.

### 2.2 Edgeworth Expansion under Diophantine conditions.

Theorem 2.1.1 is a consequence of Theorem 2.2.1 below and the fact that in our case there is a positive constant $c$ such that

$$
\begin{equation*}
|\phi(s)| \leq 1-c d(s)^{2} \tag{2.11}
\end{equation*}
$$

(2.11) follows from inequality (2.35) proven in Section 2.5 .

Theorem 2.2.1. If the distribution of $X$ has $d+2$ moments and its characteristic function satisfies

$$
\begin{equation*}
|\phi(s)| \leq 1-\frac{K}{|s|^{\gamma}} \tag{2.12}
\end{equation*}
$$

and $R<\frac{d}{2}$ is such that

$$
\begin{equation*}
\left(R-\frac{1}{2}\right) \gamma<1 \tag{2.13}
\end{equation*}
$$

then

$$
\lim _{n \rightarrow \infty} n^{R}\left[\mathbb{P}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)-\mathcal{E}_{d-1}(z)\right]=0
$$

Theorem 2.2.1 follows easily from the estimates in [20, ChapterXVI] but we provide the proof here for completeness.

Proof. Denoting

$$
\bar{\Delta}_{n}(\mathbf{a}, \mathbf{p})=\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)-\mathcal{E}_{d-1}(z)
$$

we get by [20, Chapter XVI] that for each $T$

$$
\begin{equation*}
\left|\bar{\Delta}_{n}(\mathbf{a}, \mathbf{p})\right| \leq \frac{1}{\pi} \int_{-\frac{T}{\sigma \sqrt{n}}}^{\frac{T}{\sigma \sqrt{n}}}\left|\frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d-1}(s \sigma \sqrt{n})}{s}\right| d s+\frac{C}{T} . \tag{2.14}
\end{equation*}
$$

Choose $T=B n^{R}$ with $B=\frac{C}{\varepsilon}$. Then, $\frac{C}{T}=\frac{\varepsilon}{n^{R}}$. Take a small $\delta$ and split the integral in the RHS of (2.14) into two parts.

$$
\begin{equation*}
\frac{1}{\pi} \int_{-\delta}^{\delta}\left|\frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d-1}(s \sigma \sqrt{n})}{s}\right| d s+\frac{1}{\pi} \int_{\delta<|s|<B n^{R-1 / 2 / \sigma}}\left|\frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d-1}(s \sigma \sqrt{n})}{s}\right| d s \tag{2.15}
\end{equation*}
$$

Again, by [20, Chapter XVI], we have that the first integral of (2.15) is $\mathcal{O}\left(n^{-d / 2}\right)$. Also, $\int_{|s|>\delta}\left|\frac{\hat{\mathcal{E}}_{d-1}(s \sigma \sqrt{n})}{s}\right| d s$ has exponential decay as $n \rightarrow \infty$. Put $J=\{s: \delta<$ $\left.|s|<B n^{R-1 / 2} / \sigma\right\}$. Thus, we only need to approximate,

$$
\begin{equation*}
\int_{J}\left|\frac{\phi^{n}(s)}{s}\right| d s \leq \frac{1}{\delta} \int_{J}\left|\phi^{n}(s)\right| d s \leq \frac{C}{\delta} \int_{J} \exp \left(-\bar{c} n^{1-\left(R-\frac{1}{2}\right) \gamma}\right) d s \tag{2.16}
\end{equation*}
$$

where the last inequality is due to (2.12). By (2.13) the integral decay faster than any power of $n$. Because $R<\frac{d}{2}$ the contribution of $|s| \leq \delta$ is also under control.

### 2.3 Change of variables.

Here we deduce Theorem 2.1.2 from:
Theorem 2.1.2*. For each $z$ the random variable

$$
n^{d / 2}\left[\mathcal{E}_{d}(z)-\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)\right]
$$

converges in law to $\hat{\mathcal{X}}$ where

$$
\begin{equation*}
\hat{\mathcal{X}}(\mathfrak{a}, \mathfrak{p}, \mathcal{L}, \chi)=e^{-z^{2} / 2} \frac{\left|\mathfrak{a}_{d+1}-\mathfrak{a}_{1}\right|}{2 \sigma(\mathfrak{a}, \mathfrak{p}) \sqrt{\pi^{3}}} \sum_{\mathbf{w} \in \mathcal{L} \backslash\{0\}} \frac{\sin 2 \pi \chi(\mathbf{w})}{y(\mathbf{w})} e^{-4 \pi^{2} \mathbf{x}(\mathbf{w})^{T} D_{\mathfrak{a}, \mathfrak{p}(\mathbf{x}(\mathbf{w})}} \tag{2.17}
\end{equation*}
$$

$\mathfrak{a}=\left(\mathfrak{a}_{1}, \ldots, \mathfrak{a}_{d+1}\right), \mathfrak{p}=\left(\mathfrak{p}_{1}, \ldots, \mathfrak{p}_{d+1}\right)$ and $(\mathfrak{a}, \mathfrak{p}) \in \Omega$ are distributed according to $\mathbf{P}$ and $D_{\mathfrak{a}, \mathfrak{p}}$ and $\sigma(\mathfrak{a}, \mathfrak{p})$ are defined immediately after (2.4).

In order to deduce Theorem 2.1.2 from Theorem 2.1.2* we need to show that $e^{z^{2} / 2} \frac{\hat{\mathcal{X}}}{\Lambda(\mathfrak{a}, \mathfrak{p})}$ has the same distribution as $\mathcal{X}$. To this end we rewrite the sum in (2.17) as

$$
\begin{equation*}
\frac{1}{(2 \pi)^{d-1} \operatorname{det}\left(\sqrt{D_{\mathfrak{a}, \mathfrak{p}}}\right)} \sum_{\mathbf{w} \in \mathcal{L} \backslash\{0\}} \frac{\sin (2 \pi \chi(\mathbf{w}))}{y(\mathbf{w}) /\left((2 \pi)^{d-1} \operatorname{det}\left(\sqrt{D_{\mathfrak{a}, \mathfrak{p}}}\right)\right.} e^{-4 \pi^{2}\left\|\left(\sqrt{D_{\mathfrak{a}, \mathbf{p}}} \mathbf{x}(\mathbf{w})\right)\right\|^{2}} . \tag{2.18}
\end{equation*}
$$

Let $A$ be the linear map such that

$$
A(y, \mathbf{x})=\left(\frac{y}{(2 \pi)^{d-1} \sqrt{\operatorname{det}\left(D_{\mathfrak{a}, \mathfrak{p}}\right)}}, 2 \pi \sqrt{D_{\mathfrak{a}, \mathfrak{p}}} \mathbf{x}\right)
$$

Put $(\overline{\mathcal{L}}, \bar{\chi})=A(\mathcal{L}, \chi)$. Then, using (2.7), (2.18) can be rewritten as:

$$
\frac{1}{(2 \pi)^{d-1} \operatorname{det}\left(\sqrt{D_{\mathfrak{a}, \mathfrak{p}}}\right)} \sum_{\overline{\mathbf{w}} \in \mathcal{L} \backslash\{0\}} \frac{\sin (2 \pi \bar{\chi}(\overline{\mathbf{w}}))}{y(\overline{\mathbf{w}})} e^{-\| \mathbf{x}(\overline{\mathbf{w}})) \|^{2}}
$$

Since $\operatorname{det}(A)=1$, the pair $(\overline{\mathcal{L}}, \bar{\chi})$ is distributed according to the Haar measure on $\mathbb{M}$ proving our formula for $\mathcal{X}$.

Sections 2.4-2.7 are devoted to the proof of Theorem 2.1.2*. Note that similarly to (2.9) we have

$$
\hat{\mathcal{X}}=e^{-z^{2} / 2} \frac{\left|\mathfrak{a}_{d+1}-\mathfrak{a}_{1}\right|}{2 \sigma(\mathfrak{a}, \mathfrak{p}) \sqrt{\pi^{3}}} \sum_{\mathbf{m} \in \mathbb{Z}^{d} \backslash\{\mathbf{0}\}} \frac{\sin 2 \pi \theta(\mathbf{m})}{y(\mathbf{m})} e^{-4 \pi^{2} \mathbf{x}(\mathbf{m})^{T} D_{\mathfrak{a}, \mathfrak{p} \mathbf{x}(\mathbf{m})}}
$$

The statements of Theorems 2.1.2 and 2.1.2* look similar, however, there is an important distinction. Namely the proof of Theorem 2.1.2* is constructive. In the course of the proof given $n$, a and $z$ we construct a lattice $\mathcal{L}(\mathbf{a}, n)$ and a character $\chi(\mathbf{a}, \mathbf{p}, n, z)$ such that the expression $n^{-d / 2} \hat{\mathcal{X}}(\mathbf{a}, \mathbf{p}, \mathcal{L}(\mathbf{a}, n), \chi(\mathbf{a}, \mathbf{p}, n, z))$ well-approximates the error in the Edgeworth expansion. We believe that such a
construction could be made for more general distributions where the Edgeworth expansion fails, and this will be a subject of a future investigation. So the difference between Theorems 2.1.2 and 2.1.2* is that in the first case we have only an approximation in law while in the second case we are able to obtain an approximation in probability.

### 2.4 Cut off.

### 2.4.1 Density.

Here we show that it is enough to prove Theorem 2.1.2* under the assumption that $\mathbf{P}$ has smooth density supported on a subset

$$
\Omega_{\kappa}=\left\{(\mathbf{a}, \mathbf{p}) \in \Omega: \forall i p_{i} \geq \kappa \quad \text { and } \quad \forall i \neq j\left|a_{i}-a_{j}\right| \geq \kappa\right\}
$$

for some $\kappa>0$. Indeed suppose that the theorem is true for such densities. Let $p(\mathbf{a}, \mathbf{p})$ the original density of $\mathbf{P}$. Let $\phi$ be a bounded continuous test function. Given $\varepsilon$ we can find a smooth density $\tilde{p}(\mathbf{a}, \mathbf{p})$ supported on some $\Omega_{\kappa}$ such that $\|p-\tilde{p}\|_{L^{1}} \leq \varepsilon$. In Section 2.7 we prove that

$$
\begin{equation*}
\int \phi\left(n^{d / 2} \Delta_{n}\right) \tilde{p} d \mathbf{a} d \mathbf{p} \rightarrow \iint \phi(\hat{\mathcal{X}}(\mathbf{a}, \mathbf{p}, \mathcal{L}, \boldsymbol{\theta})) \tilde{p} d \mathbf{a} d \mathbf{p} d \mu(\mathcal{L}, \boldsymbol{\theta}) \tag{2.19}
\end{equation*}
$$

where $\Delta_{n}=\mathcal{E}_{d}(z)-\mathbb{P}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)$ and $\mu$ is the Haar measure on $\left(S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z})\right) \times$ $\mathbb{T}^{d}$. Let $p_{m}(\mathbf{a}, \mathbf{p})$ be the smooth density supported on $\Omega_{\kappa}$ corresponding to $\varepsilon=m^{-1}$. Passing to subsequence, $p_{m} \rightarrow p$ almost surely. Because $\left|p_{m} \phi\right| \leq\|\phi\|\left|p_{m}\right| \in L^{1}$ and $|p \phi| \leq\|\phi\||p| \in L^{1}$ and $\|\phi\|\left|p_{m}\right| \rightarrow\|\phi\||p|$ almost surely, Lebesgue Dominated Convergence Theorem gives

$$
\begin{align*}
\iint \phi(\hat{\mathcal{X}}(\mathbf{a}, \mathbf{p}, \mathcal{L}, \boldsymbol{\theta})) p_{m} d \mathbf{a} d \mathbf{p} d \mu(\mathcal{L}, \boldsymbol{\theta}) & \\
& \rightarrow \iint \phi(\hat{\mathcal{X}}(\mathbf{a}, \mathbf{p}, \mathcal{L}, \boldsymbol{\theta})) p d \mathbf{a} d \mathbf{p} d \mu(\mathcal{L}, \boldsymbol{\theta}) \tag{2.20}
\end{align*}
$$

Combining (2.19) and (2.20) we have that,

$$
\begin{align*}
& \int \phi\left(n^{d / 2} \Delta_{n}\right) p d \mathbf{a} d \mathbf{p}=\int \phi\left(n^{d / 2} \Delta_{n}\right) p_{m} d \mathbf{a} d \mathbf{p}+\mathcal{O}\left(m^{-1}\|\phi\|\right)  \tag{2.21}\\
& \xrightarrow{n \rightarrow \infty} \iint \phi(\hat{\mathcal{X}}(\mathbf{a}, \mathbf{p}, \mathcal{L}, \boldsymbol{\theta})) p_{m} d \mathbf{a} d \mathbf{p} d \mu(\mathcal{L}, \boldsymbol{\theta})+\mathcal{O}\left(m^{-1}\|\phi\|\right) \\
& \xrightarrow{m \rightarrow \infty} \iint \phi(\hat{\mathcal{X}}(\mathbf{a}, \mathbf{p}, \mathcal{L}, \boldsymbol{\theta})) p d \mathbf{a} d \mathbf{p} d \mu(\mathcal{L}, \boldsymbol{\theta})
\end{align*}
$$

### 2.4.2 Fourier transform.

As in the previous section let

$$
\Delta_{n}=\mathcal{E}_{d}(z)-F_{n}(z) \quad \text { where } \quad F_{n}(z)=\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z\right)
$$

Denote by $v_{T}(x)=\frac{1}{\pi} \cdot \frac{1-\cos T x}{T x^{2}}$ and let $\mathcal{V}(s, T)=\left(1-\frac{|s|}{T}\right) \mathbb{1}_{|s| \leq T}$ be its Fourier transform. Using the approach of [20, Section XVI.3] we let $T_{2}=n^{2 d+6}$ and decompose

$$
\begin{equation*}
\Delta_{n}=\left[\mathcal{E}_{d}-F_{n}\right] \star v_{T_{2}}(z)-\left[F_{n}-F_{n} \star v_{T_{2}}\right](z)+\left[\mathcal{E}_{d}-\mathcal{E}_{d} \star v_{T_{2}}\right](z) . \tag{2.22}
\end{equation*}
$$

To estimate the last term we split

$$
\begin{align*}
{\left[\mathcal{E}_{d}-\mathcal{E}_{d} \star v_{T_{2}}\right](z) } & =\int_{|x| \leq 1}\left[\mathcal{E}_{d}(z)-\mathcal{E}_{d}(z-x)\right] v_{T_{2}}(x) d x  \tag{2.23}\\
& +\int_{|x| \geq 1}\left[\mathcal{E}_{d}(z)-\mathcal{E}_{d}(z-x)\right] v_{T_{2}}(x) d x
\end{align*}
$$

Since $v_{T}$ is even the first integral in (2.23) equals to

$$
\int_{|x| \leq 1} \mathcal{E}_{d}^{\prime}(z) x v_{T_{2}}(x) d x-\int_{|x| \leq 1} \frac{\mathcal{E}_{d}^{\prime \prime}(y(z, x))}{2} x^{2} v_{T_{2}}(x) d x
$$

$$
=\int_{|x| \leq 1} \frac{\mathcal{E}_{d}^{\prime \prime}(y(z, x))}{2} \frac{1-\cos T_{2} x}{\pi T_{2}} d x=\mathcal{O}\left(\frac{1}{T_{2}}\right) .
$$

Since both $\mathcal{E}_{\boldsymbol{d}}$ and cosine are bounded the second integral in (2.23) is bounded by

$$
C \int_{|x| \geq 1} \frac{d x}{T_{2} x^{2}}=\frac{C}{T_{2}} .
$$

Thus the last term in (2.22) is $\mathcal{O}\left(T_{2}^{-1}\right)$. To estimate the second term in (2.22) we split the integral in $F_{n} \star v_{T_{2}}$ into regions $\left\{|x| \geq 1 / \sqrt{T_{2}}\right\}$ and $\left\{|x| \leq 1 / \sqrt{T_{2}}\right\}$. The contribution of $\left\{|x| \geq 1 / \sqrt{T_{2}}\right\}$ is bounded by

$$
C \int_{1 / \sqrt{T_{2}}}^{\infty} \frac{d x}{T_{2} x^{2}}=\frac{C}{\sqrt{T_{2}}} .
$$

On the other hand

$$
\int_{|x| \leq 1 / \sqrt{T_{2}}}\left[F_{n}(z)-F_{n}(z-x)\right] V_{T_{2}}(x) d x=0
$$

unless there is a point of increase of $F_{n}$ inside $\left[z-1 / \sqrt{T_{2}}, z+1 / \sqrt{T_{2}}\right]$. The probability that such a point exists is bounded by

$$
\begin{equation*}
\sum_{m_{1}+\cdots+m_{d+1}=n} \mathbf{P}\left(m_{1} a_{1}+\cdots+m_{d+1} a_{d+1} \in\left[z-1 / \sqrt{T_{2}}, z+1 / \sqrt{T_{2}}\right]\right) . \tag{2.24}
\end{equation*}
$$

Note that for each fixed $\left(m_{1}, \ldots, m_{d+1}\right)$ the random variable

$$
m_{1} a_{1}+\cdots+m_{d+1} a_{d+1}
$$

has a bounded density with respect to the uniform distribution on the segment of length $\mathcal{O}\left(\sqrt{m_{1}^{2}+\cdots+m_{d+1}^{2}}\right)$ and so

$$
\mathbb{P}(\mathbf{m} \cdot \mathbf{a} \in J)=\mathcal{O}\left(\frac{|J|}{\|\mathbf{m}\|}\right)
$$

for any interval $J$. Hence each term in (2.24) is $\mathcal{O}\left(\frac{1}{n \sqrt{T_{2}}}\right)$ and so the sum is $\mathcal{O}\left(\frac{n^{d}}{n \sqrt{T_{2}}}\right)$. Thus with probability $1-\mathcal{O}\left(\frac{1}{n^{4}}\right)$ we have that $\Delta_{n}=\Delta_{n, 2}+\mathcal{O}\left(T_{2}^{-1 / 2}\right)$ where

$$
\begin{aligned}
\Delta_{n, 2} & =\frac{1}{2 \pi} \int_{-T_{2}}^{T_{2}} \frac{\left[\phi^{n}\left(\frac{t}{\sqrt{n}}\right)-\hat{\mathcal{E}}_{d}(t)\right]}{i t} \mathcal{V}\left(t, T_{2}\right) e^{-i t z} d t \\
& =\frac{1}{2 \pi} \int_{-\frac{T_{2}}{\sigma \sqrt{n}}}^{\frac{T_{2}}{\sigma \sqrt{n}}} e^{-i s z \sigma \sqrt{n}} \frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d}(s \sigma \sqrt{n})}{i s} \mathcal{V}\left(s, n, T_{2}\right) d s
\end{aligned}
$$

$\mathcal{V}(s, n, T)=1-\left|\frac{s \sigma \sqrt{n}}{T}\right|$ and $\phi(s)$ is the characteristic function of $X$ given by

$$
\phi(s)=p_{1} e^{i s a_{1}}+\cdots+p_{d+1} e^{i s a_{d+1}}
$$

Let $T_{1}=K_{1} n^{d / 2}$ and define

$$
\Delta_{n, 1}=\frac{1}{2 \pi} \int_{-\frac{T_{1}}{\sigma \sqrt{n}}}^{\frac{T_{1}}{\sigma \sqrt{n}}} e^{-i s z \sigma \sqrt{n}} \frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d}(s \sigma \sqrt{n})}{i s} \mathcal{V}\left(s, n, T_{2}\right) d s
$$

Let $\Gamma_{n}=\Delta_{n, 2}-\Delta_{n, 1}$. Put

$$
\tilde{\Gamma}_{n}=\frac{1}{2 \pi} \int_{|s| \in\left[T_{1} /(\sigma \sqrt{n}), T_{2} /(\sigma \sqrt{n})\right]} e^{-i s z \sigma \sqrt{n}} \frac{\phi^{n}(s)}{i s} \mathcal{V}\left(s, n, T_{2}\right) d s
$$

Then, we have $\Gamma_{n}=\tilde{\Gamma}_{n}+\mathcal{O}\left(e^{-\varepsilon T_{1}^{2}}\right)$ due to the exponential decay of $\hat{\mathcal{E}}_{d}$.
The main result of Subsection 2.4.2 is the following.
Proposition 2.4.1.

$$
\begin{equation*}
\left\|\tilde{\Gamma}_{n}\right\|_{L^{2}} \leq \frac{C}{\sqrt{T_{1} n^{d}}} \tag{2.25}
\end{equation*}
$$

Proof.
$\mathbf{E}\left(\tilde{\Gamma}_{n}^{2}\right)=\iint \mathbf{E}\left(e^{-i\left(s_{1}+s_{2}\right) z \sigma \sqrt{n}} \phi^{n}\left(s_{1}\right) \phi^{n}\left(s_{2}\right) \mathcal{V}\left(s_{1}, n, T_{2}\right) \mathcal{V}\left(s_{2}, n, T_{2}\right)\right) \frac{d s_{1}}{s_{1}} \frac{d s_{2}}{s_{2}}$.
We split this integral into two parts.
(1) In the region where $\left|s_{1}+s_{2}\right| \leq 1$ we use Corollary 2.5.2 proven in Section 2.5 to estimate the integral by

$$
\begin{equation*}
\mathcal{O}\left(\int_{|s| \in\left[T_{1} /(\sigma \sqrt{n}), T_{2} /(\sigma \sqrt{n})\right]} \frac{1}{\sqrt{n} s_{1}^{2}} \mathbf{E}\left(\left|\phi^{n}\left(s_{1}\right)\right|\right) d s_{1}\right) \tag{2.26}
\end{equation*}
$$

The next result will be proven in Section 2.6.

## Lemma 2.4.2.

$$
\mathbf{E}\left(\left|\phi^{n}\left(s_{1}\right)\right|\right) \leq \frac{C}{n^{d / 2}} .
$$

Plugging the estimate of Lemma 2.4.2 into (2.26) and integrating we see that the contribution of the first region to $\mathbf{E}\left(\tilde{\Gamma}_{n}^{2}\right)$ is $\mathcal{O}\left(\frac{1}{T_{1} n^{d / 2}}\right)$.
(2) Consider now the region where $\left|s_{1}+s_{2}\right| \geq 1$. Denote

$$
b_{d+1}=a_{d+1}-a_{1}, \ldots, b_{2}=a_{2}-a_{1} .
$$

Then

$$
\phi(s)=e^{i s a_{1}} \psi(s) \quad \text { where } \quad \psi(s)=p_{1}+p_{2} e^{i s b_{2}}+\cdots+p_{d+1} e^{i s b_{d+1}}
$$

Denote $\boldsymbol{\nu}=\left(p_{1}, \ldots, p_{d}, b_{2}, \ldots, b_{d}\right)$. Then there exists a compactly supported density $\rho=\rho\left(a_{1}, \boldsymbol{\nu}\right)$ such that the contribution of the second region is

$$
\iint_{\left|s_{1}+s_{2}\right| \geq 1}\left(\int e^{-i\left(s_{1}+s_{2}\right) z \sigma \sqrt{n}} e^{i n\left(s_{1}+s_{2}\right) a_{1}} \psi^{n}\left(s_{1}\right) \psi^{n}\left(s_{2}\right) \mathcal{V}\left(s_{1}\right) \mathcal{V}\left(s_{2}\right) \rho d a_{1} d \boldsymbol{\nu}\right) \frac{d s_{1}}{s_{1}} \frac{d s_{2}}{s_{2}} .
$$

We are able to use a $2 d$-dimensional coordinate system because on $\Omega$

$$
\begin{equation*}
p_{1}+\cdots+p_{d+1}=1, \quad \text { and } \quad p_{1} a_{1}+\cdots+p_{d+1} a_{d+1}=0 \tag{2.27}
\end{equation*}
$$

To estimate this integral we integrate by parts with respect to $a_{1}$. We use that

$$
e^{i s n a_{1}} d a_{1}=\left[\frac{1}{i s n} \frac{d}{d a_{1}}\right]^{k} d e^{i s n a_{1}}
$$

for some large $k$ (for example we can take $k=2 d+1$ ). The integration by parts amounts to applying $\left(\frac{d}{d a_{1}}\right)^{k}$ to $\left(e^{i s z \sigma \sqrt{n}} \rho\left[\psi\left(s_{1}\right) \psi\left(s_{2}\right)\right]^{n}\right)$ which leads to the terms

$$
\left\{\left(\frac{d}{d a_{1}}\right)^{k_{1}}\left[e^{i\left(s_{1}+s_{2}\right) z \sigma \sqrt{n}}\right]\right\}\left\{\left(\frac{d}{d a_{1}}\right)^{k_{2}}[\rho]\right\}\left\{\left(\frac{d}{d a_{1}}\right)^{k_{3}}\left[\left[\psi\left(s_{1}\right) \psi\left(s_{2}\right)\right]^{n}\right]\right\}
$$

where $k_{1}+k_{2}+k_{3}=k$. (Note that both $\sigma$ and $\psi$ depend on $a_{1}$ implicitly due to the second equation in (2.27)). Thus, the contribution of the above term to the integral is bounded by

$$
C \iint_{\left|s_{1}\right|,\left|s_{2}\right| \in\left[T_{1} / \sigma \sqrt{n}, T_{2} / \sigma \sqrt{n}\right]}^{\left|s_{1}+s_{2}\right| \geq 1} \left\lvert\, \frac{\left(s_{1}+s_{2}\right)^{k_{1}} n^{\left(k_{1} / 2\right)+k_{3}}}{\left(s_{1}+s_{2}\right)^{k} n^{k}} \mathbf{E}\left(\left|\phi^{n}\left(s_{1}\right)\right|\right) \frac{d s_{1}}{s_{1}} \frac{d s_{2}}{s_{2}} .\right.
$$

Using Lemma 2.4.2 again we can estimate the above integral by

$$
\begin{cases}\frac{C}{n^{k / 2}} & \text { if } \quad k_{1} \geq k-2 \\ \frac{C}{T_{1} n^{k+d / 2-k_{1} / 2-k_{3}}} & \text { otherwise. }\end{cases}
$$

Thus the main contribution comes from $k_{1}=k_{2}=0, k_{3}=k$ proving Proposition 2.4.1.

Proposition 2.4.1 shows that the contribution from $\tilde{\Gamma}_{n}$ to the $L^{2}$-limit of $n^{d / 2} \Delta_{n}$ can be made arbitrarily small by choosing $K_{1}$ large. Also, on $|s| \leq T_{1} / \sigma \sqrt{n}$ we have

$$
\mathcal{V}\left(s, n, T_{2}\right)=\left(1-\frac{s \sigma \sqrt{n}}{T_{2}}\right) \mathbb{1}_{|s|<T_{2} / \sigma \sqrt{n}}=1-\frac{s \sigma}{n^{2 d+\frac{11}{2}}}
$$

Hence $\Delta_{n, 1}=\hat{\boldsymbol{\Delta}}_{n}+o\left(n^{-2 d}\right)$ where

$$
\hat{\boldsymbol{\Delta}}_{n}:=\frac{1}{2 \pi} \int_{|s| \leq T_{1} / \sigma \sqrt{n}} \frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d}(s \sigma \sqrt{n})}{i s} e^{-i s z \sigma \sqrt{n}} d s
$$

approximates well $\Delta_{n, 1}$ and hence, $\Delta_{n}$ too. Also, the error from this approximation of $n^{d / 2} \Delta_{n}$ converges to 0 in $L^{2}$. Hence, we only need to analyze $n^{d / 2} \hat{\boldsymbol{\Delta}}_{n}$ for large $n$.

### 2.5 Simplifying the error.

Denote

$$
s_{k}=\frac{2 \pi k}{\left|b_{d+1}\right|}
$$

and let $I_{k}$ be the segment of length $\frac{2 \pi}{\left|b_{d+1}\right|}$ centered at $s_{k}$. Put $K_{2} \gg K_{1}$. Due to the results of the previous section it is sufficient to study

$$
\hat{\boldsymbol{\Delta}}_{n}=\sum_{|k| \leq K_{2} \sqrt{n}} \tilde{\mathcal{I}}_{k}
$$

where

$$
\tilde{\mathcal{I}}_{k}=\frac{1}{2 \pi i} \int_{I_{k}} e^{-i s z \sigma \sqrt{n}} \frac{\phi^{n}(s)-\hat{\mathcal{E}}_{d}(s \sigma \sqrt{n})}{s} d s
$$

$\tilde{\mathcal{I}}_{0}=\mathcal{O}\left(n^{-(d+1) / 2}\right)$ due to [20, Section XVI.2]. Next, $\hat{\mathcal{E}}_{d}(s \sigma \sqrt{n})$ decays exponentially with respect to $n$ outside of $I_{0}$. So, its contribution to $\tilde{\mathcal{I}}_{k}$ is negligible for $k \neq 0$.

Accordingly,

$$
\hat{\boldsymbol{\Delta}}_{n}=\sum_{0<|k| \leq K \sqrt{n}} \mathcal{I}_{k}+\mathcal{O}\left(\frac{1}{n^{(d+1) / 2}}\right)
$$

where

$$
\mathcal{I}_{k}=\frac{1}{2 \pi i} \int_{I_{k}} e^{-i s z \sigma \sqrt{n}} \frac{\phi^{n}(s)}{s} \mathbb{1}_{|s| \leq T_{1} / \sigma \sqrt{n}} d s .
$$

Introduce the following notation

$$
\bar{s}_{k}=\arg \max _{s \in I_{k}}|\phi(s)|, \quad \phi\left(\bar{s}_{k}\right)=r_{k} e^{i \phi_{k}} .
$$

The following lemma is similar to the results of [12, Section 5.2].
Lemma 2.5.1. Suppose that

$$
\begin{equation*}
r_{k}^{n} \geq n^{-100} \tag{2.28}
\end{equation*}
$$

and

$$
\begin{equation*}
\pm \frac{T_{1}}{\sigma \sqrt{n}} \notin I_{k} . \tag{2.29}
\end{equation*}
$$

Then

$$
\mathcal{I}_{k}=\frac{1}{i \sqrt{\pi n} \sigma} \frac{r_{k}^{n}}{\bar{s}_{k}} e^{-z^{2} / 2} e^{i n \phi_{k}-i \bar{s}_{k} z \sigma \sqrt{n}}\left(1+o_{n \rightarrow \infty}(1)\right)
$$

Proof. Let $e^{i \bar{s}_{k} a_{j}}=e^{i\left(\phi_{k}+\beta_{j}(k)\right)}$. Then

$$
\begin{equation*}
r_{k}=\sum_{j=1}^{d+1} p_{j} \cos \left(\beta_{j}(k)\right) \tag{2.30}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{j=1}^{d+1} p_{j} \sin \left(\beta_{j}(k)\right)=0 \tag{2.31}
\end{equation*}
$$

Since (2.28) implies that $r_{k} \geq 1-\frac{C \ln n}{n},(2.30)$ shows that $\left|\beta_{j}(k)\right| \leq C \sqrt{\frac{\ln n}{n}}$ and so (2.31) gives

$$
\begin{equation*}
\sum_{j=1}^{d+1} p_{j} \beta_{j}(k)=\mathcal{O}\left(\frac{\ln ^{3 / 2} n}{n^{3 / 2}}\right) \tag{2.32}
\end{equation*}
$$

Now we use Taylor expansion

$$
\begin{align*}
e^{i\left(\bar{s}_{k}+\delta\right) a_{j}}=e^{i \phi_{k}}\left(1+i \beta_{j}(k)-\frac{\beta_{j}(k)^{2}}{2}\right)\left(1+i a_{j} \delta-\right. & \left.\frac{a_{j}^{2} \delta^{2}}{2}\right) \\
& +\mathcal{O}\left(\frac{\ln ^{3 / 2} n}{n^{3 / 2}}+\delta^{3}\right) \tag{2.33}
\end{align*}
$$

Thus,

$$
\begin{align*}
\phi\left(\bar{s}_{k}+\delta\right) & = & e^{i \phi_{k}} \sum_{j=1}^{d+1} p_{j}\left(\cos \left(\beta_{j}(k)\right)-\frac{a_{j}^{2} \delta^{2}}{2}\right)+\mathcal{O}\left(\frac{\ln ^{3 / 2} n}{n^{3 / 2}}+\delta^{3}\right) \\
& = & r_{k} e^{i \phi_{k}}\left(1-\frac{\sigma^{2} \delta^{2}}{2}\right)+\mathcal{O}\left(\frac{\ln ^{3 / 2} n}{n^{3 / 2}}+\delta^{3}\right) \tag{2.34}
\end{align*}
$$

where we have used (2.32) as well as

$$
p_{1} a_{1}+\cdots+p_{d+1} a_{d+1}=0, \quad p_{1} a_{1}^{2}+\cdots+p_{d+1} a_{d+1}^{2}=\sigma^{2}
$$

Hence for large $n$, the main contribution to $\mathcal{I}_{k}$ equals to

$$
\begin{aligned}
& \frac{r_{k}^{n}}{2 \pi i \bar{s}_{k}} e^{i\left(n \phi_{k}-\sqrt{n} \sigma z \bar{s}_{k}\right)} \int\left(1-\frac{\sigma^{2} \delta^{2}}{2}\right)^{n} e^{-i \sigma z \delta \sqrt{n}} d \delta \\
& \approx \frac{r_{k}^{n}}{2 \pi i \bar{s}_{k}} e^{i\left(n \phi_{k}-\sqrt{n} \sigma z \bar{s}_{k}\right)} \int e^{-\sigma^{2} \delta^{2} n / 2-i \sigma \delta \sqrt{n} z} d \delta .
\end{aligned}
$$

Making the change of variables $\sigma \delta \sqrt{n / 2}=t$ we evaluate the last integral as $\frac{2 \sqrt{\pi} e^{-z^{2} / 2}}{\sigma \sqrt{n}}$.

Corollary 2.5.2. If I is a finite interval of order 1 . Then

$$
\int_{I}\left|\phi^{n}(s)\right| \mathbb{1}_{|s| \leq T_{1} / \sigma \sqrt{n}} d s=\mathcal{O}\left(\frac{1}{\sqrt{n}}\right) .
$$

Proof. We can cover $I$ by a finite number of intervals $I_{k}$. The intervals where $r_{k}^{n} \leq$ $\frac{1}{n^{100}}$ contribute $\mathcal{O}\left(\frac{|I|}{n^{100}}\right)$ while the contribution of the intervals where $r_{k}^{n} \geq \frac{1}{n^{100}}$ is $\mathcal{O}\left(\frac{1}{\sqrt{n}}\right)$ due to Lemma 2.5.1.

Because $r_{k} \approx 1, r_{k}=\left|\psi\left(\bar{s}_{k}\right)\right|=\left|p_{1}+\sum_{j=2}^{d+1} p_{j} e^{i b_{j} \bar{s}_{k}}\right| \approx \sum p_{j}$. Therefore, $a_{j} \bar{s}_{k} \approx$ $a_{1} \bar{s}_{k} \bmod 2 \pi$ for all $j \geq 2$. Thus, $\frac{2 \pi k b_{j}}{b_{d+1}} \approx 0(\bmod 2 \pi)$ for all $2 \leq j \leq d$ and hence, $\phi\left(s_{k}\right) \approx 1$ which means $s_{k}$ and $\bar{s}_{k}$ are close. Define, $\xi_{k}=\bar{s}_{k}-s_{k}, \eta_{j, k}=\frac{2 \pi k b_{j}}{b_{d+1}}+2 \pi l_{j, k}$, $j=1, \ldots, d$ where $l_{j, k}$ is the unique integer such that $\frac{2 \pi k b_{j}}{b_{d+1}}+2 \pi l_{j, k} \approx 0$. Then,

$$
\begin{align*}
r_{k}^{2}=\sum_{j=1}^{d+1} p_{j}^{2}+2 \sum_{l>j, j \neq 1} p_{l} p_{j} \cos \left[\left(b_{l}-b_{j}\right) \xi_{k}+\eta_{l, k}\right. & \left.-\eta_{j, k}\right]+2 p_{d+1} p_{1} \cos b_{d+1} \xi_{k} \\
& +2 \sum_{j=2}^{d} p_{j} p_{1} \cos \left(b_{j} \xi_{k}+\eta_{j, k}\right) \tag{2.35}
\end{align*}
$$

Therefore

$$
r_{k}^{2}=1-\sum_{l>j, j \neq 1} p_{l} p_{j}\left[\left(b_{l}-b_{j}\right) \xi_{k}+\eta_{l, k}-\eta_{j, k}\right]^{2}-p_{d+1} p_{1} b_{d+1}^{2} \xi_{k}^{2}
$$

$$
-\sum_{j=2}^{d} p_{j} p_{1}\left(b_{j} \xi_{k}+\eta_{j, k}\right)^{2}+\mathcal{O}\left(\xi_{k}^{3}+\sum_{l=1}^{d} \eta_{l, k}^{3}\right) .
$$

Taking $\eta_{1, k}=b_{1}=0$ we can write the above as,

$$
\begin{aligned}
r_{k}^{2}=-\xi_{k}^{2} \sum_{l>j} p_{l} p_{j}\left(b_{l}-b_{j}\right)^{2}-2 \xi_{k} \sum_{\substack{l>j \\
(l, j) \neq(d, 1)}} p_{l} p_{j}\left(b_{l}-b_{j}\right)\left(\eta_{l, k}-\eta_{j, k}\right) \\
+1-\sum_{\substack{l>j \\
(l, j) \neq(d, 1)}} p_{l} p_{j}\left(b_{l}-b_{j}\right)\left(\eta_{l, k}-\eta_{j, k}\right)^{2}+\mathcal{O}\left(\xi_{k}^{3}+\sum_{l=1}^{d} \eta_{l, k}^{3}\right) .
\end{aligned}
$$

Since we have $r_{k}^{2}$ approximated by a quadratic polynomial of $\xi_{k}$ (the unknown) we can approximate $\xi_{k}$ by determining the maximizer of $r_{k}^{2}\left(\xi_{k}\right)$, obtaining

$$
\begin{equation*}
\xi_{k}=-\frac{\sum_{\substack{l>j \\(l, j) \neq(d, 1)}} p_{l} p_{j}\left(b_{l}-b_{j}\right)\left(\eta_{l, k}-\eta_{j, k}\right)}{\sum_{l>j} p_{l} p_{j}\left(b_{l}-b_{j}\right)^{2}}+\mathcal{O}\left(\|\left.\boldsymbol{\eta}_{k}\right|^{2}\right) . \tag{2.36}
\end{equation*}
$$

Substituting back we find $r_{k}$ in terms of $\eta_{j, k}$ only. Ignoring higher order terms we compute the maximum to be:

$$
\begin{aligned}
& r_{k}^{2}=1-\sum_{\substack{l>j \\
(l, j) \neq(d, 1)}} p_{l} p_{j}\left(b_{l}-b_{j}\right)\left(\eta_{l, k}-\eta_{j, k}\right)^{2} \\
&+\frac{\left[\sum_{\substack{l>j \\
(l, j) \neq(d, 1)}} p_{l} p_{j}\left(b_{l}-b_{j}\right)\left(\eta_{l, k}-\eta_{j, k}\right)\right]^{2}}{\sum_{l>j} p_{l} p_{j}\left(b_{l}-b_{j}\right)^{2}}+\mathcal{O}\left(\sum_{l=1}^{d} \eta_{l, k}^{3}\right)
\end{aligned}
$$

Put $R=\left[\sum_{l>j} p_{l} p_{j}\left(b_{l}-b_{j}\right)^{2}\right]^{-1}$. Then,

$$
\begin{aligned}
& r_{k}^{2}=1+\sum_{\substack{l>j \\
(l, j) \neq(d, 1)}} p_{l} p_{j}\left(b_{l}-b_{j}\right)\left[p_{l} p_{j}\left(b_{l}-b_{j}\right) R-1\right]\left(\eta_{l, k}-\eta_{j, k}\right)^{2} \\
& +\sum_{\substack{l>j, m>n \\
l \neq m \neq n \\
(l, j),(m, n) \neq(d, 1)}} p_{l} p_{j} p_{m} p_{n}\left(b_{l}-b_{j}\right)\left(b_{m}-b_{n}\right) R\left(\eta_{l, k}-\eta_{j, k}\right)\left(\eta_{m, k}-\eta_{n, k}\right)+\mathcal{O}\left(\sum_{l>j} \eta_{l, j}^{3}\right)
\end{aligned}
$$

$$
\begin{equation*}
:=1-2 \sum_{l, j=2}^{d} D_{l, j}(\mathbf{a}, \mathbf{p}) \eta_{l, k} \eta_{j, k}+\mathcal{O}\left(\sum_{l>j} \eta_{l, j}^{3}\right) \tag{2.37}
\end{equation*}
$$

Thus,

$$
r_{k}=1-\sum_{l, j=2}^{d} D_{l, j}(\mathbf{a}, \mathbf{p}) \eta_{l, k} \eta_{j, k}+\mathcal{O}\left(\sum_{l>j} \eta_{l, j}^{3}\right)=1-\boldsymbol{\eta}_{k}^{T} D_{\mathbf{a}, \mathbf{p}} \boldsymbol{\eta}_{k}+\mathcal{O}\left(\left\|\boldsymbol{\eta}_{k}\right\|^{3}\right)
$$

where $D_{\mathbf{a}, \mathbf{p}}$ is a $(d-1) \times(d-1)$ matrix with

$$
\begin{equation*}
\left[D_{\mathbf{a}, \mathbf{p}}\right]_{i, j}=D_{i, j}(\mathbf{a}, \mathbf{p}) \tag{2.38}
\end{equation*}
$$

and $\boldsymbol{\eta}_{k}^{T}=\left(\eta_{2, k}, \ldots, \eta_{d, k}\right)$. From this we have,

$$
\mathcal{I}_{k}=\frac{e^{-z^{2} / 2}}{i \sqrt{\pi n} \sigma} \frac{\left(1-\boldsymbol{\eta}_{k}^{T} D_{\mathbf{a}, \mathbf{p}} \boldsymbol{\eta}_{k}+\mathcal{O}\left(\left\|\boldsymbol{\eta}_{k}\right\|^{3}\right)\right)^{n}}{s_{k}} e^{i n \phi_{k}-i \bar{s}_{k} z \sigma \sqrt{n}}(1+o(1)) .
$$

Let $\mathcal{B}(\mathbf{a}, \mathbf{p})$ be the contribution of the boundary terms $\pm \frac{T_{1}}{\sigma \sqrt{n}} \in I_{k}$.
Lemma 2.5.3.

$$
\mathbf{E}(|\mathcal{B}|) \leq \frac{C}{n^{(2 d-1) / 2}}
$$

Lemma 2.5.4. Let

$$
\mathcal{I}_{k, l}=\mathcal{I}_{k} \mathbb{1}_{|k|^{\alpha} n^{1 / 4}}\left\|\boldsymbol{\eta}_{k}\right\| \in\left[2^{l}, 2^{l+1}\right] .
$$

with $\alpha=[2(d-1)]^{-1}$. Then there is a constant $\tilde{c}$ such that

$$
\mathbf{E}\left(\sum_{0<|k|<K n^{(d-1) / 2}} \sum_{l>K}\left|\mathcal{I}_{k, l}\right|\right)=\mathcal{O}\left(\frac{1}{n^{d / 2}} 2^{K} \exp \left(-\tilde{c} 2^{2 K}\right)\right) .
$$

Lemmas 2.5.3 and 2.5.4 will be proven in Section 2.6.
Next we prove a lemma that would allow us to further simplify $\hat{\boldsymbol{\Delta}}_{n}$.
Lemma 2.5.5. (a) $\bar{s}_{k}=s_{k}+\boldsymbol{\omega}^{T} \boldsymbol{\eta}_{k}+\mathcal{O}\left(\|\eta\|_{k}^{2}\right)$ where $\boldsymbol{\omega}=\boldsymbol{\omega}(\mathbf{a}, \mathbf{p})$ is a $1 \times(d-1)$ vector.
(b) If $\|\boldsymbol{\eta}\|=\mathcal{O}\left(\frac{\ln n}{\sqrt{n}}\right)$ then $n \phi_{k}=n s_{k} a_{1}+n p_{2} \eta_{2, k}+\cdots+n p_{d} \eta_{d, k}+o(1)$.

Proof. Since $\bar{s}_{k}-s_{k}=\zeta_{k}$ part (a) follows by (2.36).
Next, by (2.34)

$$
\phi_{k}=\arg \phi\left(s_{k}\right)+\mathcal{O}\left(\delta^{3}+\frac{\ln ^{3 / 2} n}{n^{3 / 2}}\right)
$$

Note that,

$$
\phi\left(s_{k}\right)=e^{i s_{k} a_{1}}\left(p_{1}+p_{2} e^{i \eta_{2, k}}+\cdots+p_{d} e^{i \eta_{d, k}}+p_{d+1}\right) .
$$

Thus,

$$
\begin{aligned}
\arg \left(\phi\left(s_{k}\right)\right) & =s_{k} a_{1}+\tan ^{-1}\left(\frac{p_{2} \sin \eta_{2, k}+\cdots+p_{d} \sin \eta_{d, k}}{p_{1}+p_{2} \cos \eta_{2, k}+\cdots+p_{d} \cos \eta_{d, k}+p_{d+1}}\right) \\
& =s_{k} a_{1}+\sum_{l=2}^{d} p_{l} \eta_{l, k}+\mathcal{O}\left(\left\|\boldsymbol{\eta}_{k}\right\|^{3}\right)
\end{aligned}
$$

since the denominator in the first line is $1+\mathcal{O}\left(\|\boldsymbol{\eta}\|^{2}\right)$. Now part (b) follows easily.
Now, we continue the analysis of the leading term in $\hat{\boldsymbol{\Delta}}_{n}$. Pick a small $\delta$ and define

$$
A_{1}=\left\{(\mathbf{a}, \mathbf{p}) \mid \mathcal{I}_{k, l}=0 \forall k, l \text { s.t. }|k|<\delta n^{(d-1) / 2} \text { and } l<K\right\} .
$$

Then

$$
A_{1}^{c}=\left\{(\mathbf{a}, \mathbf{p})|\exists| k\left|<\delta n^{(d-1) / 2},|k|^{\alpha} n^{1 / 4}\left\|\eta_{k}\right\| \leq 2^{K}\right\} .\right.
$$

Thus,

$$
\mathbb{P}\left(A_{1}^{c}\right)=\sum_{|k|<\delta n^{(d-1) / 2}} \frac{C 2^{K}}{|k|^{(d-1) \alpha} n^{(d-1) / 4}}=\mathcal{O}\left(\sqrt{\delta} 2^{K}\right)
$$

if $\alpha=\frac{1}{2(d-1)}$. Hence, for a very large $K$ and $\delta$ such that $\sqrt{\delta} 2^{K}$ is very small, we can approximate $\Delta_{n}$ by the sum of $\mathcal{I}_{k}$ 's with $\delta \leq \frac{|k|}{n^{(d-1) / 2}} \leq K$ and $|k|^{\alpha} n^{1 / 4}\left\|\boldsymbol{\eta}_{k}\right\| \leq$ $2^{K}$.

We define the random vector $X_{k}=\sqrt{n} \boldsymbol{\eta}_{k}$ and $Y_{k}=\frac{k}{n^{(d-1) / 2}}$. Then, combining terms corresponding to $k$ and $-k$, we obtain the following approximation to the distribution of $\Delta_{n}$ for large $n$

$$
\frac{\left|b_{d+1}\right| e^{-z^{2} / 2}}{n^{d / 2} \sigma \sqrt{\pi^{3}}} \sum_{k \in S(n, \delta, K)} \frac{\sin \left(n \phi_{k}-\bar{s}_{k} z \sigma \sqrt{n}\right)}{Y_{k}} e^{-X_{k}^{T} D_{\mathbf{a}, \mathbf{p}} X_{k}}
$$

where $S(n, \delta, K)=\left\{k>\left.0\left|\delta<Y_{k}<K, \quad\right| Y_{k}\right|^{\alpha}\left\|X_{k}\right\|<2^{K}\right\}$.
Define $\mathbf{q}=\left(p_{2}, \ldots, p_{d}\right)$. Then, Lemma 2.5.5 shows that

$$
\begin{aligned}
n \phi_{k}-\bar{s}_{k} z \sigma \sqrt{n} & =s_{k}\left(n a_{1}-z \sigma \sqrt{n}\right)+n \mathbf{q}^{T} \boldsymbol{\eta}_{k}-z \sigma \sqrt{n} \boldsymbol{\omega}^{T} \boldsymbol{\eta}_{k}+o(1) \\
& =\frac{2 \pi n^{d / 2}}{\left|b_{d+1}\right|}\left(\sqrt{n} a_{1}-z \sigma\right) Y_{k}+(\sqrt{n} \mathbf{q}-z \sigma \boldsymbol{\omega})^{T} X_{k}+o(1) .
\end{aligned}
$$

Therefore, for large $n$ and $K$ and $\delta$ such that $\sqrt{\delta} 2^{K}$ is very small, the distribution of $\Delta_{n}$ is well approximated by
$\tilde{\Delta}_{n}(\delta, K)=\frac{\left|b_{d+1}\right| e^{-z^{2} / 2}}{n^{d / 2} \sigma \sqrt{\pi^{3}}} \sum_{k \in S(n, \delta, K)} \frac{\sin \left(\frac{2 \pi n^{d / 2}}{\left|b_{d+1}\right|}\left(\sqrt{n} a_{1}-z \sigma\right) Y_{k}+(\sqrt{n} \mathbf{q}-z \sigma \boldsymbol{\omega})^{T} X_{k}\right)}{Y_{k}} e^{-X_{k}^{T} D_{\mathbf{a}, \mathbf{p}} X_{k}}$.

### 2.6 Expectation of characteristic function.

Proof of Lemma 2.4.2. Recall that $d(s)=\max _{2 \leq j \leq d+1} d\left(b_{j} s, 0\right)$ where the distance is computed on the torus $\mathbb{R} /(2 \pi \mathbb{Z})$. Formula (2.35) shows that there are positive constants $C, c$ such that

$$
\begin{equation*}
\frac{1}{C} \leq \frac{\left|\phi^{n}(s)\right|}{e^{-c n d(s)^{2}}}<C \tag{2.39}
\end{equation*}
$$

To prove the lemma we decompose $\mathbf{E}\left(e^{-c n d(s)^{2}}\right)$ into the pieces where $d(s) \sqrt{n}$ is of order $2^{l}$ for some $l \leq\left(\log _{2} n\right) / 2$. and use the fact that $\partial$ has a bounded density.

$$
\begin{aligned}
\mathbf{E}\left(\phi^{n}(s)\right) & \leq C \mathbf{P}\left(d(s)<\frac{1}{\sqrt{n}}\right)+C \sum_{l=0}^{\left(\log _{2} n\right) / 2} \mathbf{P}\left(d(s) \sqrt{n} \in\left[2^{l}, 2^{l+1}\right]\right) e^{-c 4^{l}} \\
& \leq \frac{C}{n^{d / 2}}+C \sum_{l=0}^{\left(\log _{2} n\right) / 2} \frac{4^{l}}{n^{d / 2}} e^{-c 4^{l}} \leq \frac{C}{n^{d / 2}}
\end{aligned}
$$

completing the proof.
Proof of Lemma 2.5.3. Let $k$ be such that $\frac{T_{1}}{\sigma \sqrt{n}} \in I_{k}$. Then

$$
\mathcal{I}_{k}=\int_{\pi(2 k-1) /\left|b_{d+1}\right|}^{T_{1} / \sigma \sqrt{n}} e^{-i s z \sigma \sqrt{n}} \frac{\phi^{n}(s)}{s} d s
$$

Because $T_{1}=K_{1} n^{d / 2}$ and $s \in\left[\frac{\pi(2 k-1)}{\left|b_{d+1}\right|}, \frac{T_{1}}{\sigma \sqrt{n}}\right]$ we have $s \approx n^{(d-1) / 2}$. Thus

$$
\mathbf{E}\left(\left|\mathcal{I}_{k}\right|\right) \leq \frac{C}{n^{(d-1) / 2}} \mathbf{E}\left(\int_{\pi(2 k-1) /\left|b_{d+1}\right|}^{T_{1} / \sigma \sqrt{n}}\left|\phi^{n}(s)\right| d s\right)
$$

We claim that for all fixed $b_{d}$,

$$
\begin{equation*}
\iint e^{-c n d(s)^{2}} d s d b_{2} \ldots d b_{d-1} \leq \frac{C}{n^{d / 2}} . \tag{2.40}
\end{equation*}
$$

If this is true then using that $\rho$ is a smooth compactly supported density of $b_{d}$ we have that,

$$
\begin{aligned}
\mathbf{E}\left(\int_{\pi(2 k-1) /\left|b_{d+1}\right|}^{T_{1} / \sigma \sqrt{n}}\left|\phi^{n}(s)\right| d s\right) & =\iiint_{\pi(2 k-1) /\left|b_{d+1}\right|}^{T_{1} / \sigma \sqrt{n}}\left|\phi^{n}(s)\right| d s d b_{d} d b_{d-1} \ldots d b_{2} \\
& \leq C \iiint_{\pi(2 k-1) /|x|}^{T_{1} / \sigma \sqrt{n}} e^{-c n d(s)^{2}} \rho(x) d s d x d b_{d-1} \ldots d b_{2} \\
& \leq C \iiint e^{-c n d(s)^{2}} d s d b_{d-1} \ldots d b_{2} \rho(x) d x \\
& \leq \frac{C}{n^{d / 2}} \int \rho(x) d x=\mathcal{O}\left(\frac{1}{n^{d / 2}}\right)
\end{aligned}
$$

Thus

$$
\begin{equation*}
\mathbf{E}\left(\left|\mathcal{I}_{k}\right|\right) \leq \frac{C}{n^{(2 d-1) / 2}} \tag{2.41}
\end{equation*}
$$

Similarly, if $-\frac{T}{\sigma \sqrt{n}} \in \mathcal{I}_{k}$, then(2.41) holds. Hence, $\mathbf{E}(|\mathcal{B}|) \leq \frac{C}{n^{(2 d-1) / 2}}$ as required.
To prove (2.40) we decompose it into pieces where $d(s) \sqrt{n}$ is of order $2^{l}$. Taking $\mu$ to be the product measure $d s d b_{d-1} \ldots d b_{2}$ from (2.39) we have

$$
\begin{aligned}
& \iint e^{-c n d(s)^{2}} d s d b_{d-1} \ldots d b_{2} \leq C \mu\left\{\left(s, b_{2}, \ldots, b_{d-1}\right) \mid d(s)<1 / \sqrt{n}\right\} \\
& +C \sum_{l=0}^{\left(\log _{2} n\right) / 2} \mu\left\{\left(s, b_{2}, \ldots, b_{d-1}\right) \mid d(s) \sqrt{n} \in\left[2^{l}, 2^{l+1}\right]\right\} e^{-c 4^{l}} \\
& \quad \leq \frac{C}{n^{d / 2}}+C \sum_{l=0}^{\left(\log _{2} n\right) / 2} \frac{4^{l}}{n^{d / 2}} e^{-c 4^{l}} \leq \frac{C}{n^{d / 2}}
\end{aligned}
$$

as required.

Proof of Lemma 2.5.4. Because

$$
r_{k}=1-\boldsymbol{\eta}_{k}^{T} D_{\mathbf{a}, \mathbf{p}} \boldsymbol{\eta}_{k}+\mathcal{O}\left(\left\|\boldsymbol{\eta}_{k}\right\|^{3}\right) \text { and }|k|^{\alpha} n^{1 / 4}\left\|\boldsymbol{\eta}_{k}\right\| \in\left[2^{l}, 2^{l+1}\right]
$$

we can write

$$
r_{k}=1-c \frac{4^{l}}{|k|^{2 \alpha} \sqrt{n}}+\mathcal{O}\left(n^{-3 / 4}\right)
$$

Accordingly

$$
r_{k}^{n} \leq C e^{-\frac{c^{2 l} \sqrt{n}}{\left.|k|\right|^{2}}}
$$

Also

$$
\mathbf{P}\left(|k|^{\alpha} n^{1 / 4}\|\boldsymbol{\eta}\| \in\left[2^{l}, 2^{l+1}\right]\right) \leq \frac{C 2^{l}}{\sqrt{|k|} n^{(d-1) / 4}}
$$

Hence,

$$
\mathbf{E}\left(\mathcal{I}_{k, l}\right) \leq \frac{C e^{-\frac{c c^{2 l} \sqrt{n}}{|k|^{2 \alpha}}}}{\sqrt{n}|k|} \frac{2^{l}}{\sqrt{|k|} n^{(d-1) / 4}}=\frac{C 2^{l} e^{-\frac{C c^{2 l} \sqrt{n}}{\left.|k|\right|^{2 \alpha}}}}{|k|^{3 / 2} n^{(d+1) / 4}}
$$

Thus

$$
\sum_{l>K} \mathbf{E}\left(\mathcal{I}_{k, l}\right) \leq \frac{C 2^{K} e^{-\frac{c 2^{2 K} \sqrt{n}}{|k|^{2 \alpha}}}}{|k|^{3 / 2} n^{(d+1) / 4}}
$$

Therefore we need to estimate

$$
\begin{gather*}
\sum_{0<|k|<K n(d-1) / 2} \frac{C 2^{K} e^{-\frac{c c^{2} K \sqrt{n}}{\left.|k|\right|^{2}}}}{|k|^{3 / 2} n^{(d+1) / 4}}= \\
\frac{C}{n^{d / 2}} \sum_{0<|k|<K n^{(d-1) / 2}} \frac{1}{|k|} \sqrt{\frac{2^{2 K} n^{(d-1) / 2}}{|k|}} e^{-\frac{c c^{2 K} \sqrt{n}}{|k|^{2} \alpha}} . \tag{2.42}
\end{gather*}
$$

Split the sum over

$$
\begin{equation*}
|k| \in\left[\frac{K n^{(d-1) / 2}}{2^{s+1}}, \frac{K n^{(d-1) / 2}}{2^{s}}\right] \tag{2.43}
\end{equation*}
$$

for $s \in \mathbb{N}$. Then, for a fixed $s$ we have

$$
|k|^{2 \alpha}=\mathcal{O}\left(\frac{K^{\frac{1}{d-1}} \sqrt{n}}{2^{\frac{s}{d-1}}}\right),
$$

so each term in the sum (2.42) is of order

$$
\frac{2^{K+(3 s / 2)}}{K^{3 / 2} n^{(d-1) / 2}} \exp \left(-\frac{c 2^{2 K+\frac{s}{d-1}}}{K^{\frac{1}{d-1}}}\right) .
$$

But the number of such terms is of order $\frac{n^{(d-1) / 2}}{2^{s}}$. Hence, the sum over $k$ in (2.43) is

$$
\mathcal{O}\left(\frac{2^{K+s / 2}}{K^{3 / 2}} \exp \left(-\frac{c 2^{2 K+\frac{s}{d-1}}}{K^{\frac{1}{d-1}}}\right)\right) .
$$

Summing over $s$ we obtain the result.

### 2.7 Relation to homogeneous flows.

Given $\mathbf{u} \in \mathbb{R}^{d-1}, v \in \mathbb{R}$ consider the following function on space $\mathcal{M}$ of unimodular lattices in $\mathbb{R}^{d}$ :

$$
\begin{equation*}
\mathcal{Z}(L)=\sum_{(y, \mathbf{x}) \in L \backslash\{0\}} \frac{\sin 2 \pi\left(\mathbf{u}^{T} \mathbf{x}+v y\right)}{y} e^{-4 \pi^{2} \mathbf{x}^{T} D_{\mathbf{a}, \mathbf{p}} \mathbb{x}_{\left\{\delta<y<K, y^{\alpha}\|\mathbf{x}\|<2^{K}\right\}} .} \tag{2.44}
\end{equation*}
$$

Define $\boldsymbol{\gamma}=\frac{1}{k} \boldsymbol{\eta}$ and introduce the matrices, $H_{\gamma}=\left(\begin{array}{cc}1 & \boldsymbol{\gamma} \\ \mathbf{0}^{T} & I_{d-1}\end{array}\right), G_{t}=\left(\begin{array}{cc}e^{-(d-1) t} & 0 \\ \mathbf{0}^{T} & e^{t} I_{d-1}\end{array}\right)$.
Then, we have

$$
n^{d / 2} \tilde{\Delta}_{n}=-\frac{\left|b_{d+1}\right| e^{-z^{2} / 2}}{\sigma \sqrt{\pi^{3}}} \mathcal{Z}\left(\mathbb{Z}^{d} H_{\gamma} G_{\frac{\ln n}{2}}\right)
$$

where

$$
\mathbf{u}=\sqrt{n} \mathbf{q}-z \sigma \boldsymbol{\omega} \text { and } v=\frac{n^{d / 2}}{\left|b_{d+1}\right|}\left(\sqrt{n} a_{1}-z \sigma\right)
$$

and $\mathbf{q}$ and $\boldsymbol{\omega}$ are defined at the end of Section 2.5. Let $\mathcal{L}(n, \mathbf{a})$ be the unimodular lattice $\mathbb{Z}^{d} H_{\gamma} G_{\frac{\ln (n)}{2}}$. Let

$$
\mathbf{w}_{j}(n, \mathbf{a})=\left(y_{j}(n, \mathbf{a}), \mathbf{x}_{j}(n, \mathbf{a})\right), j=1, \ldots, d
$$

with $y_{j} \in \mathbb{R}$ and $\mathbf{x}_{j} \in \mathbb{R}^{d-1}$ be the shortest spanning set of $\mathcal{L}$. Put,

$$
\theta_{j}(n,(\mathbf{a}, \mathbf{p}))=\mathbf{u}^{T} \mathbf{x}_{j}(n, \mathbf{a})+v y_{j}(n, \mathbf{a}), j=1, \ldots, d
$$

Proposition 2.7.1. If $(\mathbf{a}, \mathbf{p})$ is distributed according to $\mathbf{P}$ then the distribution of the random vector

$$
((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a}), \boldsymbol{\theta}(n,(\mathbf{a}, \mathbf{p})))
$$

converges to $\mathbf{P} \times \mu$ as $n \rightarrow \infty$, where $\mu$ is the Haar measure on $\left[S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z})\right] \times \mathbb{T}^{d}$.

If we restrict our attention only to $((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a}))$ then the result is standard (see [39, Theorem 5.8], as well as $[18,35,45]$ ). The proof in the general case follows the approach of the proof of Proposition 5.1 in [14].

Proof. We need to show that for each bounded smooth test function $f$,

$$
\begin{equation*}
\int_{\Omega} f((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a}), \boldsymbol{\theta}) d \mathbf{P} \rightarrow \int_{\Omega \times \mathcal{M} \times \mathbb{T}^{d}} f((\mathbf{a}, \mathbf{p}), \mathcal{L}, \boldsymbol{\theta}) d \mathbf{P} d \mathcal{L} d \boldsymbol{\theta} \tag{2.45}
\end{equation*}
$$

as $n \rightarrow \infty$. Write the Fourier series expansion of $f$ :

$$
\begin{equation*}
f((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a}), \boldsymbol{\theta})=\sum_{\mathbf{k}=\left(k_{1}, \ldots, k_{d}\right) \in \mathbb{Z}^{d}} f_{\mathbf{k}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) e^{2 \pi i \mathbf{k}^{T} \boldsymbol{\theta}} \tag{2.46}
\end{equation*}
$$

Then, it is enough to prove (2.45) for individual terms in (2.46).
If $\mathbf{k}=\mathbf{0}$ then by [39, Theorem 5.8] we can conclude that

$$
\int_{\Omega} f_{\mathbf{0}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) d \mathbf{P} \rightarrow \int_{\Omega \times \mathcal{M} \times \mathbb{T}^{d}} f_{\mathbf{0}}((\mathbf{a}, \mathbf{p}), \mathcal{L}) d \mathbf{P} d \mathcal{L} d \boldsymbol{\theta}
$$

Now assume that $\mathbf{k} \neq \mathbf{0}$. Since $\Omega$ is $2 d$ dimensional, we can use $\left(p_{1}, \ldots, p_{d}, a_{1}, b_{2}, \ldots, b_{d}\right)$ as local coordinates. In these coordinates $\mathcal{L}$ is independent of $a_{1}$. Hence, $y_{j}$ 's and $\mathbf{x}_{j}$ 's are independent of $a_{1}$. Put $\boldsymbol{\nu}=\left(p_{1}, \ldots, p_{d}, b_{2}, \ldots, b_{d}\right)$. Then there exists a compactly supported density $\rho$ such that,

$$
\begin{align*}
J_{n, \mathbf{k}}= & \int f_{\mathbf{k}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) e^{2 \pi i \mathbf{k}^{T} \boldsymbol{\theta}} d \mathbf{P}  \tag{2.47}\\
= & \int f_{\mathbf{k}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) \exp 2 \pi i\left(\sqrt{n} \sum k_{j} \mathbf{q}^{T} \mathbf{x}_{j}\right) \\
& \quad \times\left[\int \rho\left(a_{1}, \boldsymbol{\nu}\right) \exp 2 \pi i\left(\frac{n^{d / 2}}{\left|b_{d+1}\right|}\left(\sqrt{n} a_{1}-z \sigma\right) \sum y_{j} k_{j}-z \sigma \sum k_{j} \boldsymbol{\omega}^{T} \mathbf{x}_{j}\right) d a_{1}\right] d \boldsymbol{\nu} .
\end{align*}
$$

Note that,

$$
\int_{\mathbb{T}^{d} \times \Omega \times \mathcal{M}} f_{\mathbf{k}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) e^{2 \pi i \mathbf{k}^{T} \boldsymbol{\theta}} d \theta_{1} \ldots d \theta_{d} d \mathbf{P} d \mathcal{L}=0
$$

because

$$
\int_{\mathbb{T}^{d}} e^{2 \pi i \mathbf{k}^{T} \boldsymbol{\theta}} d \theta_{1} \ldots d \theta_{d}=0
$$

Therefore, it is enough to prove that $J_{n, \mathbf{k}}$ converges to 0 as $n \rightarrow \infty$. To prove this we use integration by parts as follows. Put,

$$
g\left(a_{1}, \boldsymbol{\nu}\right)=\exp i\left(\frac{2 \pi n^{(d+1) / 2} \sum y_{j} k_{j}}{\left|b_{d+1}\right|} a_{1}\right)=\exp i\left(n^{(d+1) / 2} \phi(\boldsymbol{\nu}) a_{1}\right)
$$

where $\phi(\boldsymbol{\nu})=\frac{2 \pi \sum y_{j} k_{j}}{\left|b_{d+1}\right|}$ and,

$$
h\left(a_{1}, \boldsymbol{\nu}\right)=\rho\left(a_{1}, \boldsymbol{\nu}\right) \exp \left[-i\left(\frac{2 \pi n^{d / 2} \sum y_{j} k_{j}}{\left|b_{d+1}\right|}+4 \pi \sum k_{j} \boldsymbol{\omega}^{T} \mathbf{x}_{j}\right) z \sigma\left(a_{1}, \boldsymbol{\nu}\right)\right]
$$

Then, the inner integral in $(2.47)$ is $\int g\left(a_{1}, \boldsymbol{\nu}\right) h\left(a_{1}, \boldsymbol{\nu}\right) d a_{1}$. Let $\varepsilon>0$. On the set $Q_{\mathbf{k}}=\{\phi(\boldsymbol{\nu})>\varepsilon\}$ we can write

$$
g\left(a_{1}, \boldsymbol{\nu}\right) d a_{1}=\frac{1}{i \phi(\boldsymbol{\nu}) n^{(d+1) / 2}} d \exp \left(i a_{1} n^{(d+1) / 2} \phi(\boldsymbol{\nu})\right) .
$$

Integrating by parts on $Q_{\mathbf{k}}$ (note that $h$ has compact support) and using trivial bounds on $Q_{\mathrm{k}}^{c}$, we can conclude that

$$
\begin{aligned}
\left|J_{n, \mathbf{k}}\right| & \leq\left|\int \frac{\exp \left(i a_{1} n^{(d+1) / 2} \phi(\boldsymbol{\nu})\right)}{i \phi(\boldsymbol{\nu}) n^{(d+1) / 2}} h^{\prime}\left(a_{1}, \boldsymbol{\nu}\right) d a_{1}\right|+C \mathbf{P}(\{\phi(\boldsymbol{\nu}) \leq \varepsilon\}) \\
& \leq \frac{1}{\varepsilon n^{(d+1) / 2}} \int\left|h^{\prime}\left(a_{1}, \boldsymbol{\nu}\right)\right| d a_{1}+C \mathbf{P}(\{\phi(\boldsymbol{\nu}) \leq \varepsilon\})
\end{aligned}
$$

for small enough $\varepsilon$. But $h^{\prime}\left(a_{1}, \boldsymbol{\nu}\right)=\mathcal{O}\left(n^{d / 2}\right)$, hence the first term is $\mathcal{O}(1 / \sqrt{n})$. Therefore, first taking $n \rightarrow \infty$ and then taking $\varepsilon \rightarrow 0$ we have the required result.

Proposition 2.7.1 implies that as $n \rightarrow \infty$ the distribution of $n^{d / 2} \tilde{\boldsymbol{\Delta}}_{n}(\delta, K)$ converges to the distribution of

$$
\begin{equation*}
e^{-z^{2} / 2} \frac{\left|\mathfrak{a}_{d+1}-\mathfrak{a}_{1}\right|}{2 \sigma(\mathfrak{a}, \mathfrak{p}) \sqrt{\pi^{3}}} \sum_{\mathbf{m} \in \mathbb{Z}^{d} \backslash\{\mathbf{0}\}} \frac{\sin 2 \pi \theta(\mathbf{m})}{y(\mathbf{m})} e^{-4 \pi^{2} \mathbf{x}^{T} D_{\mathfrak{a}, \mathbf{p}} \mathbb{x}_{\left\{\delta<|y(\mathbf{m})|<K,|y(\mathbf{m})| \alpha\|\mathbf{x}(\mathbf{m})\|<2^{K}\right\}} . . . . . . .} \tag{2.48}
\end{equation*}
$$

Next we let $\delta \rightarrow 0$ and $K \rightarrow \infty$ in such a way that $\sqrt{\delta} 2^{K} \rightarrow 0$. Then,

$$
\mathbb{1}_{\left\{\delta<|y(\mathbf{m})|<K,|y(\mathbf{m})|^{\alpha}|\mathbf{x}(\mathbf{m})|<2^{K}\right\}} \rightarrow 1 .
$$

Thus, (2.48) converges to $\hat{\mathcal{X}}$ proving Theorem 2.1.2*.

### 2.8 Finite intervals.

The proofs of Theorems 2.1.3 and 2.1.4 are similar to the proofs of Theorems 2.1.1 and 2.1.2 so we just explain the necessary changes leaving the details to the readers.

Proof of Theorem 2.1.3. The random vector (2.10) can be approximated by $\left(\mathcal{Z}^{(1)}, \mathcal{Z}^{(2)}\right)$ where $\mathcal{Z}^{(i)}$ are defined as in (2.44) with $\mathbf{u}$ and $v$ replaced by

$$
\mathbf{u}^{(i)}=\sqrt{n} \mathbf{q}-z_{i} \sigma \boldsymbol{\omega} \text { and } v^{(i)}=\frac{n^{d / 2}}{\left|b_{d+1}\right|}\left(\sqrt{n} a_{1}-z_{i} \sigma\right)
$$

respectively. Define $\boldsymbol{\theta}^{(i)}$ as in Proposition 2.7.1 but $\mathbf{u}$ and $v$ replaced by $\mathbf{u}^{(i)}$ and $v^{(i)}$. To complete the proof we prove an analogue of Proposition 2.7.1. Namely that $\left((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a}), \boldsymbol{\theta}^{(1)}(n,(\mathbf{a}, \mathbf{p})), \boldsymbol{\theta}^{(2)}(n,(\mathbf{a}, \mathbf{p}))\right)$ converges to $\mathbf{P} \times \mu^{\prime}$ as $n \rightarrow \infty$ where $\mu^{\prime}$ is the Haar measure on $\left[S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z})\right] \times \mathbb{T}^{d} \times \mathbb{T}^{d}$.

As in the proof of Proposition 2.7.1 we prove that for individual terms in the Fourier series of a smooth function $f$ on $\left[S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z})\right] \times \mathbb{T}^{d} \times \mathbb{T}^{d}$

$$
\sum_{\left(\mathbf{k}_{1}, \mathbf{k}_{2}\right) \in \mathbb{Z}^{d} \times \mathbb{Z}^{d}} f_{\mathbf{k}_{1}, \mathbf{k}_{2}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) e^{2 \pi i\left[\mathbf{k}_{1}^{T} \boldsymbol{\theta}^{(1)}+\mathbf{k}_{2}^{T}\left(\boldsymbol{\theta}^{(1)}-\boldsymbol{\theta}^{(2)}\right)\right]}
$$

we have

$$
J_{n, \mathbf{k}_{1}, \mathbf{k}_{2}}:=\int_{\Omega} f_{\mathbf{k}_{1}, \mathbf{k}_{2}}((\mathbf{a}, \mathbf{p}), \mathcal{L}(n, \mathbf{a})) e^{2 \pi i\left[\mathbf{k}_{1}^{T} \boldsymbol{\theta}^{(1)}+\mathbf{k}_{2}^{T}\left(\boldsymbol{\theta}^{(1)}-\boldsymbol{\theta}^{(2)}\right)\right]} d \mathbf{P}
$$

$$
\xrightarrow{n \rightarrow \infty} \int_{\Omega \times \mathcal{M} \times \mathbb{T}^{d} \times \mathbb{T}^{d}} f_{\mathbf{k}_{1}, \mathbf{k}_{2}}((\mathbf{a}, \mathbf{p}), \mathcal{L}) e^{2 \pi i\left[\mathbf{k}_{1}^{T} \boldsymbol{\theta}_{1}+\mathbf{k}_{2}^{T}\left(\boldsymbol{\theta}_{1}-\boldsymbol{\theta}_{2}\right)\right]} d \mathbf{P} d \mathcal{L} d \boldsymbol{\theta}_{1} d \boldsymbol{\theta}_{2}
$$

The case $\mathbf{k}_{1}=\mathbf{k}_{2}=0$ follows from [39, Theorem 5.8]. Note that

$$
\left.\mathbf{k}_{2}^{T}\left(\boldsymbol{\theta}^{(1)}-\boldsymbol{\theta}^{(2)}\right)\right]=\left(z_{2}(n)-z_{1}(n)\right)\left(\frac{2 \pi n^{d / 2}}{\left|b_{d+1}\right|} \sum y_{j} k_{2, j}+\sum k_{2, j} \boldsymbol{\omega}^{T} \mathbf{x}_{j}\right) \sigma .
$$

If $\mathbf{k}_{1}=0$ choose appropriate local-coordinates in which $\sigma$ is a coordinate. Integrating by parts with respect to $\sigma=\sigma(\mathbf{a}, \mathbf{p})$ and using $\left|z_{1}(n)-z_{2}(n)\right| n^{d / 2} \rightarrow \infty$ we see that $J_{n, \mathbf{0}, \mathbf{k}_{2}} \rightarrow 0$ as $n \rightarrow \infty$.

If $\mathbf{k}_{1} \neq 0$ then using the same local coordinates $\left(a_{1}, \boldsymbol{\nu}\right)$ as in the proof of Proposition 2.7 .1 we can integrate by parts to conclude that $J_{n, \mathbf{k}_{1}, \mathbf{k}_{2}} \rightarrow 0$ as $n \rightarrow \infty$. The proof follows through because the leading term of $\mathbf{k}_{1}^{T} \boldsymbol{\theta}^{(1)}+\mathbf{k}_{2}^{T}\left(\boldsymbol{\theta}^{(1)}-\boldsymbol{\theta}^{(2)}\right)$ is still $n^{(d+1) / 2} \phi(\boldsymbol{\nu}) a_{1}$.

Proof of Theorem 2.1.4. To prove part (a) pick $\bar{\varepsilon}<\varepsilon$. Applying Theorem 2.1.1 we obtain that for almost every ( $\mathbf{a}, \mathbf{p}$ )

$$
\begin{gathered}
\mathbb{P}_{(\mathbf{a}, \mathbf{p})}\left(z_{1} \leq \frac{S_{n}}{\sigma \sqrt{n}} \leq z_{2}\right)=\mathcal{E}_{d-1}\left(z_{2}\right)-\mathcal{E}_{d-1}\left(z_{1}\right)+\mathcal{O}\left(n^{-(d-\bar{\varepsilon}) / 2}\right) \\
\quad=\mathfrak{n}\left(z_{1}\right) l_{n}+\mathcal{O}\left(l_{n}^{2}\right)+\mathcal{O}\left(l_{n} / \sqrt{n}\right)+\mathcal{O}\left(n^{-(d-\bar{\varepsilon}) / 2}\right)
\end{gathered}
$$

According to the assumptions of part (a) the first term is much larger than the remaining terms proving the result.

The proof of part (b) is similar except that we apply Theorem 2.1.3 instead of Theorem 2.1.1 so we only get convergence in probability.

To prove part (c) we first prove the following analogue of Theorem 2.1.3 in case $z_{2}=z_{1}+\frac{c\left|a_{d+1}-a_{1}\right|}{n^{d / 2} \sigma}$
$\frac{n^{d / 2}}{\Lambda(\mathbf{a}, \mathbf{p})}\left(e^{z_{1}^{2} / 2}\left[\mathcal{E}_{d}\left(z_{1}\right)-\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z_{1}\right)\right], e^{z_{2}^{2} / 2}\left[\mathcal{E}_{d}\left(z_{2}\right)-\mathbb{P}_{\mathbf{a}, \mathbf{p}}\left(\frac{S_{n}}{\sigma \sqrt{n}} \leq z_{2}\right)\right]\right)$
converges in law to a random vector $\left(\tilde{\mathcal{X}}_{1}, \tilde{\mathcal{X}}_{2}\right)(\mathcal{L}, \theta, c)$ where

$$
\left(\tilde{\mathcal{X}}_{1}, \tilde{\mathcal{X}}_{2}\right)(\mathcal{L}, \theta, c)=\sum_{\mathbf{m} \in \mathbb{Z}^{d} \backslash\{\mathbf{0}\}} \frac{e^{-4 \pi^{2}\|\mathbf{x}(\mathbf{m})\|^{2}}}{y(\mathbf{m})}(\sin \theta(\mathbf{m}), \sin (\theta(\mathbf{m})-c y(\mathbf{m}))) .
$$

Once this convergence is established the proof of part (c) is the same as the proof of part (b). The proof of convergence is similar to the proof of Theorem 2.1.3 except that $\boldsymbol{\theta}^{(1)}$ and $\boldsymbol{\theta}^{(2)}$ are now not independent. Namely using the same notation as in the proof of Theorem 2.1.3 we have that $\mathbf{u}^{(2)}=\mathbf{u}^{(1)}+o(1)$, while $v^{(2)}=v^{(1)}-c+o(1)$. Following the same argument as in the proof of Proposition 2.7.1 we obtain that $\left(\mathcal{L}(n, \mathbf{a}), \boldsymbol{\theta}^{(1)}(n, \mathbf{a}),\left[\boldsymbol{\theta}^{(2)}-\boldsymbol{\theta}^{(1)}\right](n, \mathbf{a})\right)$ converges as $n \rightarrow \infty$ to $\left(\mathcal{L}^{*}, \boldsymbol{\theta}^{*}, \hat{\boldsymbol{\theta}}^{*}\right)$ where $\left(\mathcal{L}^{*}, \boldsymbol{\theta}^{*}\right)$ is distributed according to the Haar measure on $S L_{d}(\mathbb{R}) / S L_{d}(\mathbb{Z}) \times \mathbb{T}^{d}$ and $\hat{\boldsymbol{\theta}}_{j}^{*}=\boldsymbol{\theta}_{j}^{*}-c y_{j}$. This justifies the formula for $\left(\tilde{\mathcal{X}}_{1}, \tilde{\mathcal{X}}_{2}\right)$.

# Chapter 3: Central Limit Theorem: Weakly Dependent Random Variables. 

### 3.1 Overview and main results.

Let $S_{N}=\sum_{n=1}^{N} X_{n}$ be a sum of random variables. We assume that there is a Banach space $\mathbb{B}$ and a family of bounded linear operators $\mathcal{L}_{t}: \mathbb{B} \rightarrow \mathbb{B}$ and vectors $v \in \mathbb{B}, \ell \in \mathbb{B}^{\prime}$ such that

$$
\begin{equation*}
\mathbb{E}\left(e^{i t S_{N}}\right)=\ell\left(\mathcal{L}_{t}^{N} v\right), t \in \mathbb{R} . \tag{3.1}
\end{equation*}
$$

We will make the following assumptions on the family $\mathcal{L}_{t}$.
(A1) $t \mapsto \mathcal{L}_{t}$ is continuous and there exists $s \in \mathbb{N}$ and $\delta>0$ such that $t \mapsto \mathcal{L}_{t}$ is $s$ times continuously differentiable for $|t| \leq \delta$.
(A2) 1 is an isolated and simple eigenvalue of $\mathcal{L}_{0}$, all other eigenvalues of $\mathcal{L}_{0}$ have absolute value less than 1 and its essential spectrum is contained strictly inside the disk of radius 1 (spectral gap).
(A3) For all $t \neq 0, \operatorname{sp}\left(\mathcal{L}_{t}\right) \subset\{|z|<1\}$.
(A4) There are positive real numbers $K, r_{1}, r_{2}$ and $N_{0}$ such that $\left\|\mathcal{L}_{t}^{N}\right\| \leq \frac{1}{N^{r_{2}}}$ for all $t$ satisfying $K \leq|t| \leq N^{r_{1}}$ and $N>N_{0}$.

## Remark 3.1.1.

1. In practice we would check (A3) by showing that when $t \neq 0$, the spectral radius of $\mathcal{L}_{t}$ is at most 1 and no eigenvalue of $\mathcal{L}_{t}$ is on the unit circle. Because the spectrum of a linear operator is a closed set this would imply that $\operatorname{sp}\left(\mathcal{L}_{t}\right)$ is contained in a closed disk strictly inside the unit disk.
2. Suppose (A4) holds. Let $N_{1}>N_{0}$ be such that $N_{1}^{\left(r_{1}-\epsilon\right) / r_{1}}>N_{0}$. Then, for all $N>N_{1}$,

$$
\begin{aligned}
\left\|\mathcal{L}_{t}^{N}\right\| & \leq\left\|\left(\mathcal{L}_{t}^{\left\lceil N^{\left(r_{1}-\epsilon\right) / r_{1}}\right\rceil}\right)^{N_{1}^{\epsilon / r_{1}}}\right\| \leq\left\|\left(\mathcal{L}_{t}^{\left\lceil N^{\left(r_{1}-\epsilon\right) / r_{1}}\right\rceil}\right)\right\|^{N_{1}^{\epsilon / r_{1}}} \\
& \leq \frac{1}{\left\lceil N^{\left(r_{1}-\epsilon\right) / r_{1}}\right\rceil^{r_{2} N_{1}^{\epsilon / r_{1}}}} \text { for } K \leq|t| \leq N^{r_{1}-\epsilon} \\
& \leq \frac{1}{N^{r_{2} K_{N_{1}}}}
\end{aligned}
$$

where $K_{N_{1}}=\frac{r_{1}-\epsilon}{r_{1}} N^{\epsilon / r_{1}}$. Therefore fixing $N_{1}$ large enough we can make $r_{2} K_{N_{1}}$ as large as we want. Hence, given (A4), by slightly reducing $r_{1}$, we may assume $r_{2}$ is sufficiently large.
3. Suppose (A1), (A2) and (A3) are satisfied with $s \geq 3$. Then, [24, Theorem 2.4] implies that there exists $A \in \mathbb{R}$ and $\sigma^{2} \geq 0$ such that

$$
\begin{equation*}
\frac{S_{N}-N A}{\sqrt{N}} \xrightarrow{d} \mathcal{N}\left(0, \sigma^{2}\right) \tag{3.2}
\end{equation*}
$$

Our interest is in $S_{N}$ that satisfies the CLT i.e. the case $\sigma^{2}>0$. Since in applications we specify conditions which guarantee this, in the following theorems we always assume that $\sigma^{2}>0$.

This is essentially an extension of Nagaev-Guivarc'h method. Some of the spectral assumptions in the theorem can be found in the proofs of decay of corre-
lations and the CLT using transfer operators. For example, see [24, 29, 37]. The key novelty here is the condition (A4) which guarantees a sufficient control over the characteristic function for intermediate values of $t$. This is analogous to the condition (1.3) in Theorem 1.1. In addition, parallels can be drawn between the moment condition in Theorem 1.1 with the condition $s=r+2$. The proof of the result is based on classical perturbation theory in [33], applicable due to (A1), (A2) and (A3), which provides the actual expansion and control of the error near 0, the BerryEsseen inequality (see (3.4) below) which reduces that error to a Fourier inversion integral over an interval of size $\mathcal{O}\left(n^{r / 2}\right)$ and the condition (A4).

Now we are in a position to state our first result on the existence of the classical Edgeworth expansion for random variables satisfying (A1) through (A4) which we refer to as weakly dependent random variables.

Theorem 3.1.1. Let $r \in \mathbb{N}$ with $r \geq 2$. Suppose (A1) through (A4) hold with $s=r+2$ and $r_{1}>\frac{r-1}{2}$. Then $S_{N}$ admits Edgeworth expansion of order $r$.

Next, we examine the error of the order 1 Edgeworth expansion in more detail. We first show that the order 1 expansion exists if (A1) through (A3) hold with $s=3$. Then, we show that the error of approximation can be improved if (A4) holds.

Theorem 3.1.2. Suppose (A1) through (A3) hold with $s \geq 3$. Then, the order 1 Edgeworth expansion exists.

Theorem 3.1.3. Suppose (A1) through (A4) hold with $s \geq 4$. Then,

$$
\mathbb{P}\left(\frac{S_{N}-N A}{\sqrt{N}} \leq z\right)=\mathfrak{N}(z)+\frac{P_{1}(z)}{N^{1 / 2}} \mathfrak{n}(z)+\mathcal{O}\left(\frac{1}{N^{q}}\right)
$$

where $q=\min \left\{1, \frac{1}{2}+r_{1}\right\}$.

As one would expect, more precise asymptotics than the usual $o\left(N^{-\frac{1}{2}}\right)$ are available when the characteristic function has better decay. The proof shows that the error depends mostly on the expansion of the characteristic function at 0 . This is an indication that the error in Theorem 3.1.2 cannot be improved more than by a factor of $\frac{1}{\sqrt{N}}$ even when $r_{1}$ is large.

In [9], analogous results are obtained for subshifts of finite type in the stationary case and an explicit description of the first order Edgeworth expansion is given. Here, we consider a wider class of (not necessarily stationary) sequences and give explicit descriptions of higher order Edgeworth polynomials by relating the coefficients to asymptotic moments. Also, we improve the condition

$$
H_{r}:\left|\mathbb{E}\left(e^{i t S_{N}}\right)\right| \leq K\left(1-\frac{c}{|t|^{\alpha}}\right)^{n}, \frac{\alpha(r-1)}{2}<1,|t|>K
$$

found in [9] by replacing it with (A4). In addition, this allows us to obtain better asymptotics for the first order expansion.

We also extend the results in [4] on the existence of weak Edgeworth expansions for i.i.d. random variables. In section 3.5.1, we compare their results with the ours.

Before we mention our results, we define the space $F_{k}^{m}$ of functions. Put

$$
C^{m}(f)=\max _{0 \leq j \leq m}\left\|f^{(j)}\right\|_{\mathrm{L}^{1}} \text { and } C_{k}(f)=\max _{0 \leq j \leq k}\left\|x^{j} f\right\|_{\mathrm{L}^{1}}
$$

Define

$$
C_{k}^{m}(f)=C^{m}(f)+C_{k}(f)
$$

We say $f \in F_{k}^{m}$ if $f$ is $m$ times continuously differentiable and $C_{k}^{m}(f)<\infty$.

Theorem 3.1.4. Suppose (A1) through (A4) hold with $s=r+2$. Choose $q \in \mathbb{N}$ such that $q>\frac{r+1}{2 r_{1}}$. Then, for $f \in F_{r+1}^{q+2}, S_{N}$ admits weak local Edgeworth expansion of order $r$.

Theorem 3.1.5. Suppose (A1) through (A4) hold with $s=r+2$. Choose $q \in \mathbb{N}$ such that $q>\frac{r+1}{2 r_{1}}$. Then, for $f \in F_{0}^{q+2}, S_{N}$ admits weak global Edgeworth expansion of order $r$.

In Theorem 3.1.4 and Theorem 3.1.5, $f$ is required to have at least three derivatives in order to guarantee the integrability of Fourier transforms of $f$ and its derivatives. In addition to (A1) through (A4), if we have, (A5) There exists $C, \alpha>0$ and $N_{1}$ such that $\left\|\mathcal{L}_{t}^{N}\right\| \leq \frac{C}{t^{\alpha}}$ for $|t|>N^{r_{1}}$ for $N>N_{1}$. then we can improve this assumption to $f$ having only one continuous derivative.

Theorem 3.1.4*. Suppose (A1) through (A5) hold with $s=r+2$ and $\alpha>\frac{r+1}{2 r_{1}}$ for sufficiently large $N$. Then, for $f \in F_{r+1}^{1}, S_{N}$ admits weak local Edgeworth expansion of order $r$.

Theorem 3.1.5*. Suppose (A1) through (A5) hold with $s=r+2$ and $\alpha>\frac{r+1}{2 r_{1}}$ for sufficiently large $N$. Then, for $f \in F_{0}^{1}, S_{N}$ admits weak global Edgeworth expansion of order $r$.

The proofs of these theorems are minor modifications of the proofs of the previous two theorems. This is described in remark 3.2.2 appearing after the proofs.

The next theorem gives sufficient conditions for the existence of the averaged Edgeworth expansion.

Theorem 3.1.6. Suppose (A1) through (A4) hold with $s=r+2$. Choose $q \in \mathbb{N}$
such that $q>\frac{r}{2 r_{1}}$. Then, $S_{N}$ admits averaged Edgeworth expansion of order $r$ for $f \in F_{0}^{q}$.

We note that for integer valued random variable assumptions (A3) and (A4) cannot hold since the characteristic function of $S_{N}$ is $2 \pi$-periodic. Therefore we replace (A3) by,
$\widetilde{(A 3)}$ When $t \notin 2 \pi \mathbb{Z}, \operatorname{sp}\left(\mathcal{L}_{t}\right) \subset\{|z|<1\}$ and when $t \in 2 \pi \mathbb{Z}, \operatorname{sp}\left(\mathcal{L}_{t}\right) \subset\{|z|<1\} \cup\{1\}$. Also, because of periodicity of the characteristic function, an assumption similar to (A4) is not required.

The following theorem provides conditions for the existence of asymptotic expansions for the LCLT for weakly dependent integer valued random variables. A similar result for $X_{n}$ 's that are $\mathbb{Z}^{d}$-valued, is obtained in [42]. Compare with Proposition 4.2 and 4.4 therein.

Theorem 3.1.7. Suppose $X_{n}$ are integer valued, $(A 1),(A 2)$ and $\widetilde{(A 3)}$ are satisfied with $s=r+2$. Then $S_{N}$ admits order r lattice Edgeworth expansion.

The layout of the rest of the chapter is as follows. In section 3.2 we prove the results mentioned earlier by constructing the Edgeworth polynomials using characteristic functions and concluding that they satisfy the required asymptotics. In section 3.3 we relate the coefficients of these polynomials to moments of $S_{N}$ and provide an algorithm to compute coefficients. A few applications of the Edgeworth expansions such as the Local Central Limit Theorem and Moderate Deviations, are discussed in section 3.4. In the last section we give examples of sequences of random variables for which our theory can be applied. First, we revisit the i.i.d. case and
recover previous results. Then, we focus on non-trivial examples like observations arising from piece-wise expanding maps of an interval, Markov chains with finitely many states and markov processes which are strongly ergodic.

### 3.2 Proofs of the main results.

Here we prove the results mentioned earlier. From now on we work in the setting described in section 3.1.

Proof of Theorem 3.1.1. We seek polynomials $P_{p}(x)$ with real coefficients such that

$$
\begin{equation*}
\mathbb{P}\left(\frac{S_{n}-n A}{\sqrt{n}} \leq x\right)-\mathfrak{N}(x)=\sum_{p=1}^{r} \frac{P_{p}(x)}{n^{p / 2}} \mathfrak{n}(x)+o\left(n^{-r / 2}\right) \tag{3.3}
\end{equation*}
$$

Once we have found suitable candidates for $P_{p}(x)$ we can apply the Berry-Esseen inequality,

$$
\begin{equation*}
\left|F_{n}(x)-\mathcal{E}_{r, N}(x)\right| \leq \frac{1}{\pi} \int_{-T}^{T}\left|\frac{\widehat{F}_{n}(t)-\widehat{\mathcal{E}}_{r, n}(t)}{t}\right| d t+\frac{C_{0}}{T} \tag{3.4}
\end{equation*}
$$

where

$$
F_{n}(x)=\mathbb{P}\left(\frac{S_{n}-n A}{\sqrt{n}} \leq x\right), \quad \mathcal{E}_{r, n}(x)=\mathfrak{N}(x)+\sum_{p=1}^{r} \frac{P_{p}(x)}{n^{p / 2}} \mathfrak{n}(x)
$$

and $C_{0}$ is independent of $T$. We refer the reader to [20, Chapter XVI.3] for a proof of (3.4). What follows is a formal derivation of $P_{p}(x)$. Later, we will use (3.4) along with other estimates to prove (3.3).

It follows from (A1), (A2) and classical perturbation theory (see [33, IV.3.6 and VII.1.8]) that there exist $\delta>0$ such that for $|t| \leq \delta, \mathcal{L}_{t}$ has a top eigenvalue $\mu(t)$ which is simple and the remainder of the spectrum is contained in a strictly
smaller disk. One can express $\mathcal{L}_{t}$ as

$$
\begin{equation*}
\mathcal{L}_{t}=\mu(t) \Pi_{t}+\Lambda_{t} \tag{3.5}
\end{equation*}
$$

where $\Pi_{t}$ is the eigenprojection to the top eigenspace of $\mathcal{L}_{t}$ and $\Lambda_{t}=\left(I-\Pi_{t}\right) \mathcal{L}_{t}$. Because $\Lambda_{t} \Pi_{t}=\Pi_{t} \Lambda_{t}=0$, iterating (3.5), we obtain

$$
\mathcal{L}_{t}^{n}=\mu^{n}(t) \Pi_{t}+\Lambda_{t}^{n} .
$$

Using (A3) and compactness, there exist $C$ (which does not depend on $n$ and $t$ ) and $0<r<1$ such that $\left\|\Lambda_{t}^{n}\right\| \leq C r^{n}$ for all $|t| \leq \delta$. By (3.1),

$$
\begin{equation*}
\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)=\mu\left(\frac{t}{\sqrt{n}}\right)^{n} \ell\left(\Pi_{t / \sqrt{n}} v\right)+\ell\left(\Lambda_{t / \sqrt{n}}^{n} v\right) \tag{3.6}
\end{equation*}
$$

Now, we focus on the first term of (3.6). Put

$$
\begin{equation*}
Z(t)=\ell\left(\Pi_{t} v\right) \tag{3.7}
\end{equation*}
$$

Then, substituting $t=0$ in (3.6) yields $1=Z(0)+\ell\left(\Lambda_{0}^{n} v\right)$. Also, we know that $\lim _{n \rightarrow \infty}\left\|\Lambda_{0}^{n} v\right\|=0$. This gives $\lim _{n \rightarrow \infty} \ell\left(\Lambda_{0}^{n} v\right)=0$. Therefore, $Z(0)=1$ and $Z(t) \neq 0$ when $|t|<\delta$. Also, this shows that $\ell\left(\Lambda_{0}^{n} v\right)=0$ for all $n$. Next, note that $t \mapsto \mu(t)$ and $t \mapsto \Pi_{t}$ are $r+2$ times continuously differentiable on $|t|<\delta$ (see [33, IV.3.6 and VII.1.8]). Therefore, $Z(t)$ is $r+2$ times continuously differentiable on $|t|<\delta$.

Now we are in a position to compute $P_{p}(x)$. To this end we make use of ideas in [20, Chapter XVI] (where the Edgeworth expansions for i.i.d. random variables are constructed) and [24] (where the CLT is proved using Nagaev-Guivarc'h method).

Consider the function $\psi$ such that,
$\log \mu\left(\frac{t}{\sqrt{n}}\right)=\frac{i A t}{\sqrt{n}}-\frac{\sigma^{2} t^{2}}{2 n}+\psi\left(\frac{t}{\sqrt{n}}\right) \Longleftrightarrow \mu^{n}\left(\frac{t}{\sqrt{n}}\right)=e^{\frac{i n A t}{\sqrt{n}}-\frac{\sigma^{2} t^{2}}{2}} \exp \left(n \psi\left(\frac{t}{\sqrt{n}}\right)\right)$.
where $A=\lim _{n \rightarrow \infty} \mathbb{E}\left(\frac{S_{n}}{n}\right)$ is the asymptotic mean and $\sigma^{2}=\lim _{n \rightarrow \infty} \mathbb{E}\left(\left[\frac{S_{n}-n A}{\sqrt{n}}\right]^{2}\right)$ is the asymptotic variance. (For details see section 3.3.)

By (3.6) we have,

$$
\begin{equation*}
\mathbb{E}\left(e^{i t \frac{S_{n-n}}{\sqrt{n}}}\right)=e^{-\frac{\sigma^{2} t^{2}}{2}} \exp \left(n \psi\left(\frac{t}{\sqrt{n}}\right)\right) Z\left(\frac{t}{\sqrt{n}}\right)+e^{-\frac{i n A t}{\sqrt{n}}} \ell\left(\Lambda_{\frac{t}{\sqrt{n}}}^{n} v\right) \tag{3.8}
\end{equation*}
$$

Notice that $\psi(0)=\psi^{\prime}(0)=0$ and $\psi(t)$ is $r+2$ times continuously differentiable. Now, denote by $t^{2} \psi_{r}(t)$ the order $(r+2)$ Taylor approximation of $\psi$. Then, $\psi_{r}$ is the unique polynomial such that $\psi(t)=t^{2} \psi_{r}(t)+o\left(|t|^{r+2}\right)$. Also, $\psi_{r}(0)=0$ and $\psi_{r}$ is a polynomial of degree $r$. In fact, we can write $\psi(t)=t^{2} \psi_{r}(t)+t^{r+2} \tilde{\psi}_{r}(t)$ where $\tilde{\psi}_{r}$ is continuous and $\tilde{\psi}_{r}(0)=0$. Thus,

$$
\exp \left(n \psi\left(\frac{t}{\sqrt{n}}\right)\right)=\exp \left(t^{2} \psi_{r}\left(\frac{t}{\sqrt{n}}\right)+\frac{1}{n^{r / 2}} t^{r+2} \tilde{\psi}_{r}\left(\frac{t}{\sqrt{n}}\right)\right)
$$

Denote by $Z_{r}(t)$ the order $-r$ Taylor expansion of $Z(t)-1$. Then, $Z_{r}(0)=0$ and $Z(t)=1+Z_{r}(t)+t^{r} \tilde{Z}_{r}(t)$ with twice continuously differentiable $\tilde{Z}_{r}(t)$ such that $\tilde{Z}_{r}(0)=0$. Then, to make the order $n^{-j / 2}$ terms explicit, we compute:

$$
\begin{aligned}
e^{\frac{\sigma^{2} t^{2}}{2}} & \mu^{n}\left(\frac{t}{\sqrt{n}}\right) Z\left(\frac{t}{\sqrt{n}}\right) \\
= & e^{\frac{\sigma^{2} t^{2}}{2}} \mu^{n}\left(\frac{t}{\sqrt{n}}\right) \exp \log Z\left(\frac{t}{\sqrt{n}}\right) \\
= & \exp \left(t^{2} \psi_{r}\left(\frac{t}{\sqrt{n}}\right)+\frac{1}{n^{r / 2}} t^{r+2} \tilde{\psi}_{r}\left(\frac{t}{\sqrt{n}}\right)\right. \\
& \left.-\sum_{k=1}^{r} \frac{(-1)^{k+1}}{k}\left[Z_{r}\left(\frac{t}{\sqrt{n}}\right)\right]^{k}-\frac{1}{n^{r / 2}} t^{r} \bar{Z}_{r}\left(\frac{t}{\sqrt{n}}\right)\right)
\end{aligned}
$$

$$
\begin{align*}
& =1+\sum_{m=1}^{r} \frac{1}{m!}\left[t^{2} \psi_{r}\left(\frac{t}{\sqrt{n}}\right)-\sum_{k=1}^{r} \frac{(-1)^{k+1}}{k}\left[Z_{r}\left(\frac{t}{\sqrt{n}}\right)\right]^{k}\right]^{m} \\
& \quad+\frac{1}{n^{r / 2}} t^{r+2} \tilde{\psi}_{r}\left(\frac{t}{\sqrt{n}}\right)-\frac{1}{n^{r / 2}} t^{r} \bar{Z}_{r}\left(\frac{t}{\sqrt{n}}\right)+t^{r+1} \mathcal{O}\left(n^{-\frac{r+1}{2}}\right) \\
& =\sum_{k=0}^{r} \frac{A_{k}(t)}{n^{k / 2}}+\frac{t^{r}}{n^{r / 2}} \varphi\left(\frac{t}{\sqrt{n}}\right)+t^{r+1} \mathcal{O}\left(n^{-\frac{r+1}{2}}\right) \tag{3.9}
\end{align*}
$$

where $A_{0} \equiv 1, \varphi(t)=t^{2} \tilde{\psi}_{r}(t)-\bar{Z}_{r}(t)$ is continuous and $\varphi(0)=0$. Here $\bar{Z}_{r}$ is the remainder of $\log Z(t)$ when approximated by powers of $Z_{r}$. Next write,

$$
\begin{equation*}
Q_{n}(t)=\sum_{k=1}^{r} \frac{A_{k}(t)}{n^{k / 2}} \tag{3.10}
\end{equation*}
$$

Notice that

$$
\begin{equation*}
A_{k} \text { and } k \text { have the same parity. } \tag{3.11}
\end{equation*}
$$

This can be seen directly from the construction, because we collect terms with the same power of $n^{-1 / 2}, \psi_{r}$ and $Z_{r}$ are a polynomial in $\frac{t}{\sqrt{n}}$ with no constant term and we take powers of $t^{2} \psi_{r}(t)$ and $Z_{r}(t)$, the resulting $A_{k}$ will contain terms of the form $c_{s} t^{2 s+k}$.

We claim that,

$$
\begin{align*}
\int_{|t|<\delta \sqrt{n}} & \left|\frac{\mu^{n}\left(\frac{t}{\sqrt{n}}\right) Z\left(\frac{t}{\sqrt{n}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}-e^{-\frac{t^{2} \sigma^{2}}{2}} Q_{n}(t)}{t}\right| d t  \tag{3.12}\\
& =\int_{|t|<\delta \sqrt{n}} e^{-\frac{t^{2} \sigma^{2}}{2}}\left|\frac{\exp \left[n \psi\left(\frac{t}{\sqrt{n}}\right)+\log Z\left(\frac{t}{\sqrt{n}}\right)\right]-1-Q_{n}(t)}{t}\right| d t \\
& =o\left(n^{-r / 2}\right)
\end{align*}
$$

We note that from the choice of $Q_{n}$,

$$
\frac{\exp \left[n \psi\left(\frac{t}{\sqrt{n}}\right)+\log Z\left(\frac{t}{\sqrt{n}}\right)\right]-1-Q_{n}(t)}{t}=\frac{1}{n^{r / 2}}\left(t^{r-1} \varphi\left(\frac{t}{\sqrt{n}}\right)+t^{r} \mathcal{O}\left(n^{-\frac{r+1}{2}}\right)\right)
$$

where $\varphi(t)=o(1)$ as $t \rightarrow 0$. As a result, for all $\varepsilon>0$ the integrand of (3.12) can be made smaller than $\frac{\varepsilon}{n^{r / 2}}\left(t^{r-1}+t^{r}\right) e^{-\frac{t^{2} \sigma^{2}}{2}}$ by choosing $\delta$ small enough. This proves the claim.

Even though the following derivation is only valid for $|t|<\delta \sqrt{n}$, once the polynomial function $Q_{n}(t)$ is obtained as above, we can consider it to be defined for all $t \in \mathbb{R}$.

Suppose $|t| \leq \delta$. From classical perturbation theory (see [33, Chapter IV] and $[29$, Section 7]) we have

$$
\begin{equation*}
\Lambda_{t}^{n}=\frac{1}{2 \pi i} \int_{\Gamma} z^{n}\left(z-\mathcal{L}_{t}\right)^{-1} d z \tag{3.13}
\end{equation*}
$$

where $\Gamma$ is the positively oriented circle centered at $z=0$ with radius $\varepsilon_{0}$. Here $\varepsilon_{0}$ is uniform in $t$ and $0<\varepsilon_{0}<1$. Now,

$$
\begin{aligned}
\Lambda_{t}^{n}-\Lambda_{0}^{n} & =\frac{1}{2 \pi i} \int_{\Gamma} z^{n}\left[\left(z-\mathcal{L}_{t}\right)^{-1}-\left(z-\mathcal{L}_{t}\right)^{-1}\right] d z \\
& =\frac{1}{2 \pi i} \int_{\Gamma} z^{n}\left[\left(z-\mathcal{L}_{0}\right)^{-1}\left(\mathcal{L}_{t}-\mathcal{L}_{0}\right)\left(z-\mathcal{L}_{t}\right)^{-1}\right] d z
\end{aligned}
$$

Because $\mathcal{L}_{t}-\mathcal{L}_{0}=\mathcal{O}(|t|)$ we have that $\frac{\Lambda_{t}^{n}-\Lambda_{0}^{n}}{|t|}=\mathcal{O}\left(\varepsilon_{0}^{n}\right) . \quad \ell \in \mathbb{B}^{\prime}$ and $\ell\left(\Lambda_{0}^{n} v\right)=0$ implies that

$$
\begin{aligned}
\int_{|t|<\delta \sqrt{n}}\left|\frac{e^{-\frac{i n A t}{\sqrt{n}}} \ell\left(\Lambda_{t / \sqrt{n}}^{n} v\right)}{t}\right| d t & =\int_{|t|<\delta \sqrt{n}}\left|\frac{e^{-\frac{i n A t}{\sqrt{n}}} \ell\left(\Lambda_{t / \sqrt{n}}^{n} v-\Lambda_{0}^{n} v\right)}{t}\right| d t \\
& \leq C \int_{|t|<\delta}\left|\frac{\Lambda_{t}^{n}-\Lambda_{0}^{n}}{t}\right| d t=\mathcal{O}\left(\varepsilon_{0}^{n}\right)
\end{aligned}
$$

This decays exponentially fast to 0 as $n \rightarrow \infty$. This allows us to control the second term in the RHS of (3.6). Combining this with (3.12) we can conclude that,

$$
\begin{equation*}
\int_{|t|<\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}-e^{-\frac{t^{2} \sigma^{2}}{2}} Q_{n}(t)}{t}\right| d t=o\left(n^{-r / 2}\right) . \tag{3.14}
\end{equation*}
$$

Observe that,

$$
(i t)^{k} e^{-\frac{\sigma^{2} t^{2}}{2}}=\frac{1}{\sqrt{2 \pi \sigma^{2}}} \frac{\widehat{d^{k}} e^{k} e^{-\frac{t^{2}}{2 \sigma^{2}}}}{=\frac{\widehat{d^{k}}}{d t^{k}} \mathfrak{n}(t)}
$$

where $\widehat{f}(x)=\int e^{-i t x} f(t) d t$ is the Fourier transform of $f$. Therefore,

$$
\begin{equation*}
R_{j}(t) \mathfrak{n}(t)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} A_{j}\left(-i \frac{d}{d t}\right)\left[e^{-\frac{t^{2}}{2 \sigma^{2}}}\right] \tag{3.15}
\end{equation*}
$$

Then, the required $P_{p}(x)$ for $p \geq 1$, can be found using the relation,

$$
\begin{equation*}
\mathfrak{n}(x) R_{p}(x)=\frac{d}{d x}\left[\mathfrak{n}(x) P_{p}(x)\right] . \tag{3.16}
\end{equation*}
$$

For more details, we refer the reader to [20, Chapter XVI.3,4].
Given $\varepsilon>0$, choose $B>\frac{C_{0}}{\varepsilon}$ where $C_{0}$ is as in (3.4). Let $r \in \mathbb{N}$. Then we choose polynomials $P_{p}(x)$ as described above. Then, from (3.4) it follows that,

$$
\begin{aligned}
\left|F_{n}(x)-\mathcal{E}_{r, n}(x)\right| & \leq \frac{1}{\pi} \int_{-B n^{r / 2}}^{B n^{r / 2}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}\left(1+Q_{n}(t)\right)}{t}\right| d t+\frac{C_{0}}{B n^{r / 2}} \\
& \leq I_{1}+I_{2}+I_{3}+\frac{\varepsilon}{n^{r / 2}}
\end{aligned}
$$

where

$$
\begin{gathered}
I_{1}=\frac{1}{\pi} \int_{|t|<\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}\left(1+Q_{n}(t)\right)}{t}\right| d t \\
I_{2}=\frac{1}{\pi} \int_{\delta \sqrt{n}<|t|<B n^{r / 2}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t \\
I_{3}=\frac{1}{\pi} \int_{|t|>\delta \sqrt{n}} e^{-\frac{t^{2} \sigma^{2}}{2}}\left|\frac{1+Q_{n}(t)}{t}\right| d t .
\end{gathered}
$$

From (3.12) we have that $I_{1}$ is $o\left(n^{-r / 2}\right)$. Because our choice of $\varepsilon>0$ is arbitrary the proof is complete, if $I_{2}$ and $I_{3}$ are also $o\left(n^{-r / 2}\right)$. These follow from (3.18), (3.19) and (3.17) below.

It is easy to see that,

$$
\begin{equation*}
\int_{|t|>\delta \sqrt{n}} e^{-\frac{t^{2} \sigma^{2}}{2}}\left|\frac{1+Q_{n}(t)}{t}\right| d t=\mathcal{O}\left(e^{-c n}\right) \tag{3.17}
\end{equation*}
$$

for some $c>0$. Thus, we only need to control,

$$
\begin{aligned}
I_{2} & =\int_{\delta \sqrt{n}<|t|<B n^{r / 2}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t \\
& =\int_{\delta \sqrt{n}<|t|<\bar{\delta} \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t+\int_{\bar{\delta} \sqrt{n}<|t|<B n^{r / 2}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t
\end{aligned}
$$

where $\bar{\delta}>\max \{\delta, K\}$ with $K$ as in (A4).
By (A3) the spectral radius of $\mathcal{L}_{t}$ has modulus strictly less than 1 . Because $t \mapsto \mathcal{L}_{t}$ is continuous, for all $p<q$, there exists $\gamma<1$ and $C>0$, such that $\left\|\mathcal{L}_{t}^{m}\right\| \leq C \gamma^{m}$ for all $p \leq|t| \leq q$ for sufficiently large $m$. Then using (3.1) for sufficiently large $n$ we have,

$$
\begin{equation*}
\int_{\delta \sqrt{n}<|t|<\bar{\delta} \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t \leq \frac{1}{\delta \sqrt{n}} \int_{\delta \sqrt{n}<|t|<\bar{\delta} \sqrt{n}}\left\|\mathcal{L}_{t / \sqrt{n}}^{n}\right\| d t \leq \frac{C \gamma^{n}}{\sqrt{n}} \tag{3.18}
\end{equation*}
$$

This shows that the integral converges to 0 faster than any inverse power of $\sqrt{n}$. Next for sufficiently large $n$,

$$
\begin{align*}
\int_{\bar{\delta} \sqrt{n}<|t|<B n^{r / 2}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t & \leq \frac{1}{\bar{\delta} \sqrt{n}} \int_{\bar{\delta} \sqrt{n}<|t|<B n^{r / 2}}\left|\ell\left(\mathcal{L}_{t / \sqrt{n}}^{n} v\right)\right| d t  \tag{3.19}\\
& \leq \frac{2 B n^{r / 2}}{\bar{\delta} n^{r_{2}+1 / 2}}\|\ell\|\|v\| \\
& =C n^{\frac{r-1}{2}-r_{2}}=o\left(n^{-r / 2}\right)
\end{align*}
$$

The second inequality is due to assumption (A4) i.e. $\left\|\mathcal{L}_{t / \sqrt{n}}^{n}\right\| \leq \frac{1}{n^{r_{2}}}$ where $r_{2}>\frac{r-1}{2}$ (we can assume $r_{2}>\frac{r-1}{2}$ for large $n$ due to Remark 3.1.1) and $K \leq \bar{\delta}<\frac{|t|}{\sqrt{n}}<B n^{\frac{r-1}{2}} \leq n^{r_{1}}$ for $n \in \mathbb{N}$ with $n^{r_{1}-\frac{r-1}{2}} \geq B$.

The proof of Theorem 3.1.2 follows the same idea. We include its proof for completion.

Proof of Theorem 3.1.2. Because (A1) through (A3) hold with $s \geq 3$, we have (3.9) where $\varphi$ is continuous, $\varphi(0)=0$ and $r=1$. Given $\varepsilon>0$, choose $B>\frac{C_{0}}{\varepsilon}$. Then,

$$
\begin{aligned}
\left|F_{n}(x)-\mathcal{E}_{1, n}(x)\right| & \leq \frac{1}{\pi} \int_{-B \sqrt{n}}^{B \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}\left(1+Q_{n}(t)\right)}{t}\right| d t+\frac{C_{0}}{B \sqrt{n}} \\
& \leq I_{1}+I_{2}+I_{3}+\frac{\varepsilon}{B \sqrt{n}} .
\end{aligned}
$$

Because, $\varphi(t)=o(1)$ as $t \rightarrow 0$ and

$$
\frac{\exp \left[n \psi\left(\frac{t}{\sqrt{n}}\right)+\log Z\left(\frac{t}{\sqrt{n}}\right)\right]-1-Q_{1}(t)}{t}=\frac{1}{\sqrt{n}} \varphi\left(\frac{t}{\sqrt{n}}\right)+t \mathcal{O}\left(\frac{1}{n}\right)
$$

we have that,

$$
I_{1}=\int_{|t|<\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}-e^{-\frac{t^{2} \sigma^{2}}{2}} Q_{1}(t)}{t}\right| d t=o\left(n^{-1 / 2}\right)
$$

Also, $I_{3}=\mathcal{O}\left(e^{-c n}\right)$. Finally, because of (A3) there is $\gamma<1$ such that,

$$
\int_{\delta \sqrt{n}<|t|<B \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t=\int_{\delta<|t|<B}\left|\frac{\mathbb{E}\left(e^{i t S_{n}}\right)}{t}\right| d t \leq C \sup _{\delta \leq|t| \leq B}\left\|\mathcal{L}_{t}^{n}\right\| \leq C \gamma^{n}
$$

Combining these estimates we have the result.

A slight modification of the previous proof gives us the proof of Theorem 3.1.3. Higher regularity assumption gives us better asymptotics near 0 and the assumption on the faster decay of the characteristic function gives us more control in the mid range.

Proof of Theorem 3.1.3. Because (A1) through (A4) hold with $s \geq 4$, we have (3.9) where $\varphi$ is $C^{1}, \varphi(0)=0$ and $r=1$. Then,

$$
\left|F_{n}(x)-\mathcal{E}_{1, n}(x)\right| \leq \frac{1}{\pi} \int_{-n^{1 / 2+r_{1}}}^{n^{1 / 2+r_{1}}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}\left(1+Q_{n}(t)\right)}{t}\right| d t+\frac{C_{0}}{n^{1 / 2+r_{1}}}
$$

$$
\leq I_{1}+I_{2}+I_{3}+\frac{C_{0}}{n^{1 / 2+r_{1}}}
$$

Because, $\varphi\left(\frac{t}{\sqrt{n}}\right) \sim \frac{t}{\sqrt{n}}$ near 0 and

$$
\frac{\exp \left[n \psi\left(\frac{t}{\sqrt{n}}\right)+\log Z\left(\frac{t}{\sqrt{n}}\right)\right]-1-Q_{1}(t)}{t}=\frac{1}{\sqrt{n}} \varphi\left(\frac{t}{\sqrt{n}}\right)+t \mathcal{O}\left(\frac{1}{n}\right)
$$

we have that,

$$
I_{1}=\int_{|t|<\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{t^{2} \sigma^{2}}{2}}-e^{-\frac{t^{2} \sigma^{2}}{2}} Q_{1}(t)}{t}\right| d t=\mathcal{O}\left(\frac{1}{n}\right)
$$

Also, $I_{3}=\mathcal{O}\left(e^{-c n}\right)$. As before, (3.18) holds for $\bar{\delta}>\max \{\delta, K\}$.

$$
\begin{aligned}
& \left\|\mathcal{L}_{t}^{n}\right\| \leq \frac{1}{n^{r_{2}}} \text { where } K \leq \bar{\delta}<|t|<n^{r_{1}} \\
& \quad \int_{\bar{\delta} \sqrt{n}<|t|<n^{1 / 2+r_{1}}}\left|\frac{\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)}{t}\right| d t=\int_{\bar{\delta}<|t|<n^{r_{1}}}\left|\frac{\mathbb{E}\left(e^{i t S_{n}}\right)}{t}\right| d t \leq C n^{r_{1}-r_{2}+\frac{1}{2}}
\end{aligned}
$$

Because $r_{2}$ can be made arbitrarily large by choosing $n$ large enough, $I_{2}=\mathcal{O}\left(\frac{1}{n}\right)$. Therefore,

$$
\left|F_{n}(x)-\mathcal{E}_{1, n}(x)\right|=\mathcal{O}\left(\frac{1}{n^{s}}\right)
$$

where $s=\min \left\{1, \frac{1}{2}+r_{1}\right\}$ and we have the required conclusion.
Remark 3.2.1. In the proof above, $I_{1}$ gives the contribution to the error from the expansion of the characteristic function near 0 . This dominates when $r_{1} \geq \frac{1}{2}$.

Weak forms of Edgeworth expansions are discussed in detail in [4]. We adapt the ideas found in [4] to our proofs of Theorems 3.1.4 and 3.1.5. One key difference is the requirement on $f$ to have two more derivatives than required in [4]. This compensates for the lack of control over the tail of the characteristic function of $S_{N}$. In fact, it is enough to assume $1+\alpha$ more derivatives. But to avoid technicalities
we stick to the stronger regularity assumption. In the i.i.d. case as shown in [4], a Diophantine assumption takes care of this. See section 3.5.1 for a detailed discussion.

Proof of Theorem 3.1.4. Recall that $\widehat{f}(t)=\int e^{-i t x} f(x) d x$ and pick $A$ as in (3.2). Then by Plancherel theorem,

$$
\begin{align*}
\mathbb{E}\left(f\left(S_{n}-n A\right)\right) & =\frac{1}{2 \pi} \int \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t  \tag{3.20}\\
\Longrightarrow \sqrt{n} \mathbb{E}\left(f\left(S_{n}-n A\right)\right) & =\frac{1}{2 \pi} \int \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t
\end{align*}
$$

We first estimate RHS away from 0 . Fix small $\delta>0$. (A particular $\delta$ is chosen later). Notice that for all $\delta \leq|t| \leq K$ (where $K$ as in (A4)), there exists $c_{0} \in(0,1)$ such that $\left\|\mathcal{L}_{t}^{n}\right\| \leq c_{0}^{n}$. Thus,

$$
\left|\int_{\delta<|t|<K} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right| \leq \int_{\delta<|t|<K}\left|\widehat{f}(t) \ell\left(\mathcal{L}_{t}^{n} v\right)\right| d t \leq C\|f\|_{1} c_{0}^{n}
$$

By Remark 3.1.1, for large $n$ we can assume $r_{2}>r_{1}+(r+1) / 2$. Therefore,

$$
\begin{aligned}
\left|\int_{K<|t|<n^{r_{1}}} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right| \leq\|f\|_{1}\|\ell\|\|v\| \int_{K<|t|<n^{r_{1}}}\left\|\mathcal{L}_{t}^{n}\right\| d t & \leq \frac{C\|f\|_{1}}{n^{r_{2}-r_{1}}} \\
& =\|f\|_{1} o\left(n^{-(r+1) / 2}\right)
\end{aligned}
$$

Because $f \in F_{r+1}^{q+2}$, we have that $t^{q} \widehat{f}(t)=(-i)^{q} \widehat{f^{(q)}}(t)$ and $\widehat{f^{(q)}}$ is integrable. In fact, $\left|\widehat{f^{(q)}}(t)\right| \leq \frac{C}{(1+|t|)^{2}}$. Note that we are using only the fact that $f$ is $q+2$ times continuously differentiable with integrable derivatives. Therefore for this to be true $f \in F_{0}^{q+2}$ is sufficient. Integrability of $\widehat{f^{(q)}}$ along with $q>\frac{r+1}{2 r_{1}}$ implies,

$$
\begin{equation*}
\left|\int_{|t|>n^{r_{1}}} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right| \leq \int_{|t|>n^{r_{1}}}|\widehat{f}(t)| d t \leq \int_{|t|>n^{r_{1}}}\left|\frac{\widehat{f^{(q)}}(t)}{t^{q}}\right| d t \tag{3.21}
\end{equation*}
$$

$$
\leq \frac{\left\|\widehat{f^{(q)}}\right\|_{1}}{n^{r_{1} q}}=\left\|\widehat{f^{(q)}}\right\|_{1} o\left(n^{-(r+1) / 2}\right)
$$

Therefore,

$$
\begin{equation*}
\left|\int_{|t|>\delta} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right|=o\left(n^{-(r+1) / 2}\right) . \tag{3.22}
\end{equation*}
$$

From (3.8), for $|t| \leq \delta \sqrt{n}$, we have,

$$
\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)=e^{-\frac{\sigma^{2} t^{2}}{2}} e^{t^{2} \mathcal{O}(\delta)}(1+\mathcal{O}(\delta))+\mathcal{O}\left(\epsilon_{0}^{n}\right)
$$

Thus, choosing small $\delta$, for large $n$ when $|t|<\delta \sqrt{n}$ there exist $c, C>0$ such that

$$
\left|\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)\right| \leq C e^{-c t^{2}}
$$

Then,

$$
\sqrt{D \log n}<|t|<\delta \sqrt{n} \Longrightarrow\left|\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)\right| \leq C e^{-c D \log n}=\frac{C}{n^{c D}}
$$

and

$$
\begin{aligned}
\left|\int_{\sqrt{\frac{D \log n}{n}}<|t|<\delta} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right| & =\left|\int_{\sqrt{D \log n}<|t|<\delta \sqrt{n}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) \frac{d t}{\sqrt{n}}\right| \\
& \leq \frac{C}{n^{c D}} \int_{\sqrt{\frac{D \log n}{n}}<|t|<\delta}|\widehat{f}(t)| d t=\frac{2 \delta C\|f\|_{1}}{n^{c D}} .
\end{aligned}
$$

Combining this with (3.22) and choosing $D$ such that, $c D>(r+1) / 2$ we have that,

$$
\begin{equation*}
\left|\int_{|t|>\sqrt{\frac{D \log n}{n}}} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right|=o\left(n^{-(r+1) / 2}\right) \tag{3.23}
\end{equation*}
$$

Next, suppose $|t|<\sqrt{\frac{D \log n}{n}}$. Then,

$$
\widehat{f}(t)=\sum_{j=0}^{r} \frac{\widehat{f}^{(j)}(0)}{j!} t^{j}+\frac{t^{r+1}}{(r+1)!} \widehat{f}^{(r+1)}(\epsilon(t))
$$

where $0 \leq|\epsilon(t)| \leq|t|$. Note that,

$$
\left|\widehat{f}^{(r+1)}(\epsilon(t))\right|=\left|\int x^{r+1} e^{-i \epsilon(t) x} f(x) d x\right| \leq \int\left|x^{r+1} f(x)\right| d x \leq C_{r+1}(f)
$$

Therefore,

$$
\begin{aligned}
\int_{|t|<\sqrt{D \log n}} & \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t \\
= & \sum_{j=0}^{r} \frac{\widehat{f}^{(j)}(0)}{j!n^{j / 2}} \int_{|t|<\sqrt{D \log n}} t^{j} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t \\
& \quad \frac{1}{n^{(r+1) / 2}} \frac{1}{(r+1)!} \int_{|t|<\sqrt{D \log n}} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) t^{r+1} \widehat{f}^{(r+1)}\left(\epsilon\left(\frac{t}{\sqrt{n}}\right)\right) d t
\end{aligned}
$$

where

$$
\left|\int_{|t|<\sqrt{D \log n}} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) t^{r+1} \widehat{f}^{(r+1)}\left(\epsilon\left(\frac{t}{\sqrt{n}}\right)\right) d t\right| \leq C_{r+1}(f) \int|t|^{r+1} e^{-c t^{2}} d t
$$

for large $n$. Hence,

$$
\begin{align*}
\int_{|t|<\sqrt{D \log n}} & \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t \\
& =\sum_{j=0}^{r} \frac{\widehat{f}^{(j)}(0)}{j!n^{j / 2}} \int_{|t|<\sqrt{D \log n}} t^{j} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t+C_{r+1}(f) \mathcal{O}\left(n^{-(r+1) / 2}\right) \tag{3.24}
\end{align*}
$$

Because $s=r+2$, from (3.9),

$$
\begin{align*}
e^{\frac{\sigma^{2} t^{2}}{2}} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) & =\exp \left(n \psi\left(\frac{t}{\sqrt{n}}\right)\right) Z\left(\frac{t}{\sqrt{n}}\right)+e^{-\frac{i n A t}{\sqrt{n}}+\frac{\sigma^{2} t^{2}}{2}} \ell\left(\Lambda_{t / \sqrt{n}}^{n} v\right) \\
& =\sum_{k=0}^{r} \frac{A_{k}(t)}{n^{k / 2}}+\frac{t^{r}}{n^{r / 2}} \varphi\left(\frac{t}{\sqrt{n}}\right)+\mathcal{O}\left(\frac{\log ^{(r+1) / 2}(n)}{n^{(r+1) / 2}}\right) \tag{3.25}
\end{align*}
$$

Substituting this in (3.24),

$$
\begin{align*}
& \int_{|t|<\sqrt{D \log n}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t  \tag{3.26}\\
& =\sum_{j=0}^{r} \frac{\widehat{f}^{(j)}(0)}{j!n^{j / 2}} \int_{|t|<\sqrt{D \log n}} t^{j} e^{-\sigma^{2} t^{2} / 2} \sum_{k=0}^{r} \frac{A_{k}(t)}{n^{k / 2}} d t+\mathcal{O}\left(\frac{\log ^{(r+1) / 2}(n)}{n^{(r+1) / 2}}\right) \\
& =\sum_{k=0}^{r} \sum_{j=0}^{r} \frac{\widehat{f}^{(j)}(0)}{j!n^{(k+j) / 2}} \int_{|t|<\sqrt{D \log n}} t^{j} A_{k}(t) e^{-\sigma^{2} t^{2} / 2} d t+o\left(n^{-r / 2}\right) .
\end{align*}
$$

Recall from (3.11) that $A_{k}$ and $k$ have the same parity. Therefore, if $k+j$ is odd then

$$
\int_{|t|<\sqrt{D \log n}} t^{j} A_{k}(t) e^{-\sigma^{2} t^{2} / 2} d t=0
$$

So only integral powers of $n^{-1}$ will remain in the expansion. Also there is $C$ that depends only on $r$ such that,

$$
\int_{|t| \geq \sqrt{D \log n}} t^{j} A_{k}(t) e^{-\sigma^{2} t^{2} / 2} d t \leq C \int_{|t| \geq \sqrt{D \log n}} t^{4 r} e^{-\sigma^{2} t^{2} / 2} d t \leq \frac{C}{e^{\sigma^{2} D \log (n) / 4}}=\frac{C}{n^{\sigma^{2} D / 4}}
$$

Choosing $D$ such that $2 \sigma^{2} D>(r+1) / 2$,

$$
\int_{\mathbb{R}} t^{j} A_{k}(t) e^{-\sigma^{2} t^{2} / 2} d t=\int_{|t| \leq \sqrt{D \log n}} t^{j} A_{k}(t) e^{-\sigma^{2} t^{2} / 2} d t+o\left(n^{-r / 2}\right)
$$

Therefore, fixing $D$ large, we can assume the integrals to be over the whole real line.
Now, define

$$
a_{k, j}=\int_{\mathbb{R}} t^{j} A_{k}(t) e^{-\sigma^{2} t^{2} / 2} d t
$$

and substitute

$$
\widehat{f}^{(j)}(0)=\int_{\mathbb{R}}(-i t)^{j} f(t) d t
$$

in (3.26) to obtain,

$$
\int_{|t|<\sqrt{D \log n}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t=\sum_{k=0}^{r} \sum_{j=0}^{r} a_{k, j} \frac{1}{j!n^{(k+j) / 2}} \int_{\mathbb{R}}(-i t)^{j} f(t) d t+o\left(n^{-r / 2}\right)
$$

$$
\begin{align*}
& =\sum_{p=0}^{r} \frac{1}{n^{p}} \int_{\mathbb{R}} f(t) \sum_{k+j=2 p} \frac{a_{k, j}}{j!}(-i t)^{j} d t+o\left(n^{-r / 2}\right)  \tag{3.27}\\
& =\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{n^{p}} \int_{\mathbb{R}} f(t) P_{p, l}(t) d t+o\left(n^{-r / 2}\right)
\end{align*}
$$

where

$$
\begin{equation*}
P_{p, l}(t)=\sum_{k+j=2 p} \frac{a_{k, j}}{j!}(-i t)^{j} \tag{3.28}
\end{equation*}
$$

The final simplification was done by absorbing the terms corresponding to higher powers of $n^{-1}$ into the error term. Note that $P_{p, l}$ is a polynomial of degree at most $2 p$ and that once we know $A_{0}, \ldots, A_{2 p}$ we can compute $P_{p, l}$.

Finally combining (3.27) and (3.23) substituting in (3.20) we obtain the required result as shown below.

$$
\begin{aligned}
\sqrt{n} \mathbb{E}\left(f\left(S_{n}-n A\right)\right)= & \frac{1}{2 \pi} \int_{|t|<\sqrt{D \log n}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t \\
& \quad+\frac{\sqrt{n}}{2 \pi} \int_{|t|>\sqrt{\frac{D \log n}{n}}} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t \\
= & \frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{n^{p}} \int_{\mathbb{R}} f(t) P_{p, l}(t) d t+o\left(n^{-r / 2}\right)+\sqrt{n} o\left(n^{-(r+1) / 2}\right) \\
= & \frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{n^{p}} \int_{\mathbb{R}} f(t) P_{p, l}(t) d t+o\left(n^{-r / 2}\right) .
\end{aligned}
$$

The proof of Theorem 3.1.5 uses the relation (3.25) derived in the previous proof. But we do not use the Taylor expansion of $\widehat{f}$, so differentiability of $\widehat{f}$ is not required. So the assumption on the decay of $f$ at infinity can be relaxed.

Proof of Theorem 3.1.5. Multiplying (3.25) by $\widehat{f}$ and integrating we obtain,

$$
\left.\begin{array}{rl}
\int_{|t|<\sqrt{D \log n}} & \widehat{f}\left(\frac{t}{\sqrt{n}}\right)
\end{array}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t .
$$

As in the proof of Theorem 3.1.4 the integrals above can be replaced by integrals over $\mathbb{R}$ without altering the order of the error because

$$
\int_{|t| \geq \sqrt{D \log n}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) A_{k}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t \leq\|f\|_{1} o\left(n^{-r / 2}\right)
$$

for $D$ such that $2 \sigma^{2} D>(r+1) / 2$. Therefore,

$$
\int_{|t|<\sqrt{D \log n}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t=\sum_{k=0}^{r} \frac{1}{n^{k / 2}} \int_{\mathbb{R}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) A_{k}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t+\|f\|_{1} o\left(n^{-r / 2}\right)
$$

We pick $R_{p}$ as in (3.15) and claim $P_{p, g}=R_{p}$.
Note that $\sqrt{n} f(t \sqrt{n}) \longleftrightarrow \widehat{f}(t / \sqrt{n})$. So by the Plancherel theorem,

$$
\int_{\mathbb{R}} \sqrt{n} f(t \sqrt{n}) R_{k}(t) \mathfrak{n}(t) d t=\frac{1}{2 \pi} \int_{\mathbb{R}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) A_{k}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t
$$

Thus,

$$
\begin{align*}
\frac{1}{2 \pi \sqrt{n}} \int_{|t|<\sqrt{D \log n}} & \widehat{f}\left(\frac{t}{\sqrt{n}}\right) \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) d t \\
& =\frac{1}{\sqrt{n}}\left(\sum_{p=0}^{r} \frac{1}{n^{p / 2}} \int_{\mathbb{R}} \sqrt{n} f(t \sqrt{n}) R_{p}(t) \mathfrak{n}(t) d t+\|f\|_{1} o\left(n^{-r / 2}\right)\right) \\
& =\sum_{p=0}^{r} \frac{1}{n^{p / 2}} \int_{\mathbb{R}} f(t \sqrt{n}) R_{p}(t) \mathfrak{n}(t) d t+\|f\|_{1} o\left(n^{-(r+1) / 2}\right) \tag{3.29}
\end{align*}
$$

Note that (3.23) holds because $f \in F_{0}^{q+2}$. Now, combining (3.29) with the estimate (3.23) completes the proof.

Remark 3.2.2. Proofs of both the Theorem 3.1.4* and Theorem 3.1.5* are almost identical except the estimate (3.21). In order to obtain the same asymptotics, the assumption on the integrability of $\widehat{f^{(q)}}$ can be replaced by (A5) and the fact that $|\widehat{f}(t)| \sim \frac{1}{t}$ for as $t \rightarrow \pm \infty$.

$$
\left|\int_{|t|>n^{r_{1}}} \widehat{f}(t) \mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right) d t\right| \leq C \int_{|t|>n^{r_{1}}}|\widehat{f}(t)|\left\|\mathcal{L}_{t}^{n}\right\| d t
$$

$$
\begin{aligned}
& \leq C\|f\|_{1} \int_{|t|>n^{r_{1}}} \frac{1}{t^{1+\alpha}} d t \\
& \leq \frac{C\|f\|_{1}}{n^{r_{1}(\alpha-\epsilon)}} \int \frac{1}{t^{1+\epsilon}} d t
\end{aligned}
$$

Since, $r_{1} \alpha>\frac{r+1}{2}$ choosing $\epsilon$ small enough we can make the expression $\|f\|_{1} o\left(n^{-(r+1) / 2}\right)$ as required.

Proof of Theorem 3.1.6. Select $A$ as in (3.2). Define $P_{p}$ by (3.15) and (3.16) and $\tilde{f}_{n}(x)=f(-\sqrt{n} x)$. Then the change of variables $-\frac{y}{\sqrt{n}} \rightarrow y$ yields,
$\int\left[\mathbb{P}\left(\frac{S_{n}-n A}{\sqrt{n}} \leq x+\frac{y}{\sqrt{n}}\right)-\mathfrak{N}\left(x+\frac{y}{\sqrt{n}}\right)-\mathcal{E}_{r, n}\left(x+\frac{y}{\sqrt{n}}\right)\right] f(y) d y=\sqrt{n} \Delta_{n} * \tilde{f}_{n}(x)$. where $\mathcal{E}_{r, n}(x)=\sum_{p=1}^{r} \frac{1}{n^{p / 2}} P_{p}(x) \mathfrak{n}(x)$.

Notice that $\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) \widehat{\tilde{f}}_{n} \in L^{1}$. Therefore,

$$
\left(F_{n} * \tilde{f}_{n}\right)^{\prime}(x)=\frac{1}{2 \pi} \int e^{-i t x} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right) \widehat{\tilde{f}}_{n}(t) d t
$$

Also,

$$
\left[\mathfrak{n}+\left(\sum_{p=1}^{r} \frac{1}{n^{p / 2}} R_{p} \mathfrak{n}\right)\right] * \tilde{f}_{n}(x)=\frac{1}{2 \pi} \int e^{-i t x} e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right) \widehat{\tilde{f}}_{n}(t) d t
$$

where $R_{p}$ 's are polynomials given by (3.15) and $Q_{n}(t)$ is given by (3.10). From these we conclude that,

$$
\begin{equation*}
\left(\Delta_{n} * \tilde{f}_{n}\right)^{\prime}(x)=\frac{1}{2 \pi} \int e^{-i t x}\left(\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right) \widehat{\tilde{f}}_{n}(t) d t\right. \tag{3.30}
\end{equation*}
$$

We claim that,

$$
\begin{equation*}
\left(\Delta_{n} * \tilde{f}_{n}\right)(x)=\frac{1}{2 \pi} \int e^{-i t x} \frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right)}{-i t} \widehat{\tilde{f}}_{n}(t) d t \tag{3.31}
\end{equation*}
$$

Indeed, if the right side of (3.31) converges absolutely, then Riemann-Lebesgue Lemma gives us that it converges 0 as $|x| \rightarrow \infty$. Differentiating (3.31) we obtain (3.30). Thus the two sides in (3.31) can differ only by a constant. Since both are 0 at $\pm \infty$, this constant is 0 and (3.31) holds.

Now, we are left with the task of showing that the right side of (3.31) converges absolutely. From the definition of $\tilde{f}_{n}$ it follows that, $\widehat{\tilde{f}}_{n}(t)=\frac{1}{\sqrt{n}} \widehat{f}\left(-\frac{t}{\sqrt{n}}\right)$. Combining this with (3.14), we have that,

$$
\begin{aligned}
\mid \int_{|t|<\delta \sqrt{n}} e^{-i t x} & \left.\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right)}{-i t} \widehat{\tilde{f}}_{n}(t) d t \right\rvert\, \\
& \leq \int_{|t|<\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right)}{t} \widehat{\tilde{f}}_{n}(t)\right| d t \\
& \leq \frac{\|f\|_{1}}{\sqrt{n}} \int_{|t|<\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right)}{t}\right| d t \\
& =\|f\|_{1} o\left(n^{-(r+1) / 2}\right) .
\end{aligned}
$$

Note that,

$$
\begin{aligned}
\mid \int_{|t|>\delta \sqrt{n}} e^{-i t x} & \left.\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right)}{-i t} \widehat{\tilde{f}}_{n}(t) d t \right\rvert\, \\
& \leq \int_{|t|>\delta \sqrt{n}}\left|\frac{\mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)-e^{-\frac{\sigma^{2} t^{2}}{2}}\left(1+Q_{n}(t)\right)}{t} \widehat{f}\left(-\frac{t}{\sqrt{n}}\right)\right| d t \\
& \leq \frac{1}{\sqrt{n}} \int_{|t|>\delta}\left|\frac{\mathbb{E}\left(e^{-i t\left(S_{n}-n A\right)}\right)-e^{-\frac{n^{2} \sigma^{2} t^{2}}{2}}\left(1+Q_{n}(-\sqrt{n} t)\right)}{t} \widehat{f}(t)\right| d t \\
& \leq \frac{1}{\sqrt{n}} \int_{|t|>\delta}\left|\frac{\mathbb{E}\left(e^{-i t\left(S_{n}-n A\right)}\right)}{t} \widehat{f}(t)\right| d t+\mathcal{O}\left(e^{-c n^{2}}\right) .
\end{aligned}
$$

Put,

$$
J_{n}=\frac{1}{\sqrt{n}} \int_{|t|>\delta}\left|\frac{\mathbb{E}\left(e^{-i t\left(S_{n}-n A\right)}\right)}{t} \widehat{f}(t)\right| d t
$$

We claim $J_{n}=o\left(n^{-(r+1) / 2}\right)$. This proves that (3.31) converges absolutely as required.

To conclude the asymptotics of $J_{n}$, choose $\bar{\delta}>\max \{\delta, K\}$ where $K$ as in (A4). From (A3) there exists $\gamma<1$ such that $\left\|\mathcal{L}_{t}^{n}\right\| \leq \gamma^{n}$ for all $\delta \leq|t| \leq \bar{\delta}$ for sufficiently large $n$. Then, using (3.1) for sufficiently large $n$ we have,

$$
\frac{1}{\sqrt{n}} \int_{\delta<|t|<\bar{\delta}}\left|\frac{\mathbb{E}\left(e^{-i t\left(S_{n}-n A\right)}\right)}{t} \widehat{f}(t)\right| d t \leq \frac{C\|f\|_{1}}{\delta \sqrt{n}} \int_{\delta<|t|<\bar{\delta}}\left\|\mathcal{L}_{t}^{n}\right\| d t=\mathcal{O}\left(\gamma^{n}\right)
$$

Next, for $K \leq \bar{\delta} \leq|t| \leq n^{r_{1}},\left\|\mathcal{L}_{t}^{n}\right\| \leq \frac{1}{n^{r_{2}}}$. Hence, for $n$ sufficiently large so that $r_{2}>\frac{r}{2}$,

$$
\begin{aligned}
\frac{1}{\sqrt{n}} \int_{\bar{\delta}<|t|<n^{r_{1}}}\left|\frac{\mathbb{E}\left(e^{-i t\left(S_{n}-n A\right)}\right)}{t} \widehat{f}(t)\right| d t & \leq \frac{C}{\delta \sqrt{n}} \int_{\bar{\delta}<|t|<n^{r_{1}}}\left\|\mathcal{L}_{t}^{n}\right\||\widehat{f}(t)| d t \\
& \leq \frac{C\|\mid \widehat{f}\|_{1}}{n^{r_{2}+1 / 2}}=o\left(n^{-(r+1) / 2}\right)
\end{aligned}
$$

Since $q>\frac{r}{2 r_{1}}$, we have that,

$$
\begin{aligned}
\frac{1}{\sqrt{n}} \int_{|t|>n^{r_{1}}}\left|\frac{\mathbb{E}\left(e^{-i t\left(S_{n}-n A\right)}\right)}{t} \widehat{f}(t)\right| d t & \leq \frac{\left\|f^{(q)}\right\|_{1}}{\sqrt{n}} \int_{|t|>n^{r_{1}}} \frac{1}{|t|^{q+1}} d t \quad \leq \frac{C\left\|f^{(q)}\right\|_{1}}{n^{q r_{1}+1 / 2}} \\
& =o\left(n^{-(r+1) / 2}\right)
\end{aligned}
$$

Combining the above estimates, $J_{n}=C^{q}(f) o\left(n^{-(r+1) / 2}\right)$.
This completes the proof that $\left(\Delta_{n} * \tilde{f}_{n}\right)(x)=o\left(n^{-(r+1) / 2}\right)$. Hence,

$$
\begin{aligned}
\int\left[\mathbb { P } \left(\frac{S_{n}-n A}{\sqrt{n}}\right.\right. & \left.\left.\left.\leq x+\frac{y}{\sqrt{n}}\right)-\mathfrak{N}\left(x+\frac{y}{\sqrt{n}}\right)\right)\right] f(y) d y \\
& =\int \mathcal{E}_{r, n}\left(x+\frac{y}{\sqrt{n}}\right) f(y) d y+\sqrt{n} \Delta_{n} * \tilde{f}_{n}(x) \\
& =\sum_{p=1}^{r} \frac{1}{n^{p / 2}} \int P_{p}\left(x+\frac{y}{\sqrt{n}}\right) \mathfrak{n}(x) f(y) d y+C^{q}(f) o\left(n^{-r / 2}\right)
\end{aligned}
$$

as required.

In the lattice case, periodicity allows us to simplify the proof significantly although the idea behind the proof is similar to the previous proofs.

Proof of Theorem 3.1.7. Under assumptions (A1) and (A2) we have the CLT for $S_{n}$. Put $A$ as in (3.2). We observe that,

$$
2 \pi \mathbb{P}\left(S_{n}=k\right)=\int_{-\pi}^{\pi} e^{-i t k} \mathbb{E}\left(e^{i t S_{n}}\right) d t=\int_{-\pi}^{\pi} e^{-i t k} \ell\left(\mathcal{L}_{t}^{n} v\right) d t
$$

After changing variables and using (3.6), (3.7) we have,

$$
\begin{equation*}
2 \pi \sqrt{n} \mathbb{P}\left(S_{n}=k\right)=\int_{-\pi \sqrt{n}}^{\pi \sqrt{n}} e^{-\frac{i t k}{\sqrt{n}}} \mu\left(\frac{t}{\sqrt{n}}\right)^{n} Z\left(\frac{t}{\sqrt{n}}\right) d t+\int_{-\pi \sqrt{n}}^{\pi \sqrt{n}} e^{-\frac{i t k}{\sqrt{n}}} \ell\left(\Lambda_{t / \sqrt{n}}^{n} v\right) d t \tag{3.32}
\end{equation*}
$$

By $\widetilde{(\mathrm{A} 3)}$ there exists $C>0$ and $r \in(0,1)$ (both independent of $t$ ) such that $\left|\ell\left(\Lambda_{t}^{n} v\right)\right| \leq C r^{n}$ for all $t \in[-\pi, \pi]$. Therefore the second term of (3.32) decays exponentially fast to 0 as $n \rightarrow \infty$.

Now, we focus on the first term. Using the same strategy as in the proof of Theorem 3.1.1 we have,

$$
\begin{equation*}
\mu\left(\frac{t}{\sqrt{n}}\right)^{n} Z\left(\frac{t}{\sqrt{n}}\right)=e^{\frac{i n A t}{\sqrt{n}}-\frac{\sigma^{2} t^{2}}{2}}\left[1+Q_{n}(t)+o\left(n^{-r / 2}\right)\right] \tag{3.33}
\end{equation*}
$$

where $Q_{n}(t)$ is as in (3.10). Define $R_{j}$ as in (3.15).

$$
\begin{aligned}
& 2 \pi \sqrt{n} \mathbb{P}\left(S_{n}\right.=k)-2 \pi\left\{\frac{1}{\sqrt{2 \pi}} e^{-\frac{(k-n A)^{2}}{2 \sigma^{2} n}}\left(1+\sum_{j=1}^{r} \frac{\left(R_{p}(k-n A) / \sqrt{n}\right)}{n^{j / 2}}\right)\right\} \\
&=\int_{-\pi \sqrt{n}}^{\pi \sqrt{n}} e^{-\frac{i t k}{\sqrt{n}}} \mu\left(\frac{t}{\sqrt{n}}\right)^{n} Z\left(\frac{t}{\sqrt{n}}\right) d t \\
& \quad-\int_{-\infty}^{\infty} e^{-\frac{i t(k-n A)}{\sqrt{n}}} e^{-\sigma^{2} t^{2} / 2} d t-\int_{-\infty}^{\infty} e^{-\frac{i t k}{\sqrt{n}}} e^{-\frac{\sigma^{2} t^{2}}{2}} Q_{n}(t) d t+o\left(n^{-r / 2}\right) .
\end{aligned}
$$

We estimate the RHS by estimating the three integrals given below,

$$
\begin{aligned}
& I_{1}=\int_{-\delta \sqrt{n}}^{\delta \sqrt{n}} e^{-\frac{i t k}{\sqrt{n}}} \mu\left(\frac{t}{\sqrt{n}}\right)^{n} Z\left(\frac{t}{\sqrt{n}}\right)-e^{-\frac{i t(k-n A)}{\sqrt{n}}} e^{-\frac{\sigma^{2} t^{2}}{2}}\left[1+Q_{n}(t)\right] d t \\
& I_{2}=\int_{\delta \sqrt{n}<|t|<\pi \sqrt{n}} e^{-\frac{i t k}{\sqrt{n}}} \mu\left(\frac{t}{\sqrt{n}}\right)^{n} Z\left(\frac{t}{\sqrt{n}}\right) d t
\end{aligned}
$$

$$
I_{3}=\int_{|t|>\delta \sqrt{n}} e^{-\frac{i t(k-n A)}{\sqrt{n}}} e^{-\frac{\sigma^{2} t^{2}}{2}}\left[1+Q_{n}(t)\right] d t .
$$

Clearly, $\left|I_{3}\right|$ decays to 0 exponentially fast as $n \rightarrow \infty$. Also, $|\mu(2 \pi)|=1$ and $|\mu(t)| \in(0,1)$ for $0<|t|<2 \pi$. Therefore, there exists $\epsilon>0$ such that $|\mu(t)|<\epsilon$ on $\delta \leq|t| \leq \pi$. Put $M=\max _{\delta \leq|t| \leq \pi}|Z(t)|$. Then,

$$
\left|I_{2}\right| \leq M \sqrt{n} \int_{\epsilon<|t|<\pi}|\mu(t)|^{n} d t \leq 2 M(\pi-\delta) \sqrt{n} \epsilon^{n}
$$

Hence, $\left|I_{2}\right|$ decays to 0 exponentially fast as $n \rightarrow \infty$. From (3.33), we have that

$$
e^{-\frac{i t k}{\sqrt{n}}}\left[\mu\left(\frac{t}{\sqrt{n}}\right)^{n} Z\left(\frac{t}{\sqrt{n}}\right)-e^{\frac{i n A t}{\sqrt{n}}} e^{-\frac{\sigma^{2} t^{2}}{2}}\left[1+Q_{n}(t)\right]\right]=e^{-\frac{\sigma^{2} t^{2}}{2}} o\left(n^{-r / 2}\right)
$$

This implies $\left|I_{1}\right|=o\left(n^{-r / 2}\right)$. Combining these estimates we have the required result.

### 3.3 Computing coefficients.

Since $\int_{|t|>\delta} \mathbb{E}\left(e^{i t S_{n}}\right) d t$ decays sufficiently fast, the Edgeworth expansion, and hence its coefficients, depend only on the Taylor expansion of $\mathbb{E}\left(e^{i t S_{n}}\right)$ about 0 . Here we relate the coefficients of Edgeworth polynomials to the asymptotics of moments of $S_{n}$ by relating them to derivatives of $\mu(t)$ and $Z(t)$ at 0 .

Suppose (A1) through (A4) are satisfied with $s=r+2$. Recall (3.6):

$$
\begin{equation*}
\mathbb{E}\left(e^{i t S_{n}}\right)=\mu(t)^{n} \ell\left(\Pi_{t} v\right)+\ell\left(\Lambda_{t}^{n} v\right) \tag{3.34}
\end{equation*}
$$

Put $Z(t)=\ell\left(\Pi_{t} v\right)$ as before. Also write $U_{n}(t)=\ell\left(\Lambda_{t}^{n} v\right)$. We already know that $\mu(t), Z(t)$ and $U(t)$ are $r+2$ times continuously differentiable. Using (3.13) one can
show further that the derivatives of $U_{n}(t)$ satisfy:

$$
\sup _{|t| \leq \delta}\left\|U_{n}^{(k)}\right\| \leq C \varepsilon_{0}^{n}
$$

for all $n$ and for all $1 \leq k \leq r+2$.
Taking the first derivative of (3.34) at $t=0$ we have:

$$
i \mathbb{E}\left(S_{n}\right)=n \mu^{\prime}(0)+Z^{\prime}(0)+U_{n}^{\prime}(0) \Longrightarrow \lim _{n \rightarrow \infty} i \mathbb{E}\left(\frac{S_{n}}{n}\right)=\mu^{\prime}(0)
$$

In fact, using the Taylor expansion of $\log \mu(t)$ and above limit one can conclude that the number $A$ we used in the statement of the CLT in (3.2), is given by

$$
A=\lim _{n \rightarrow \infty} \mathbb{E}\left(\frac{S_{n}}{n}\right)
$$

Therefore one can rewrite (3.6) as

$$
\begin{equation*}
\mathbb{E}\left(e^{i t\left(S_{n}-n A\right)}\right)=e^{-n t \mu^{\prime}(0)} \mu(t)^{n} Z(t)+\bar{U}_{n}(t) \tag{3.35}
\end{equation*}
$$

where $\bar{U}_{n}(t)=e^{-n t \mu^{\prime}(0)} U_{n}(t)$. Also note that its derivatives satisfy $\left\|\bar{U}_{n}^{(k)}\right\|_{\infty}=\mathcal{O}\left(\varepsilon_{0}^{n}\right)$ for all $1 \leq k \leq r+2$.

From (3.35), it follows that moments of $S_{n}-n A$ can be expanded in powers of $n$ with coefficients depending on derivatives of $\mu$ and $Z$ at 0 . However, only powers of $n$ upto order $k / 2$ will appear. We prove this fact below.

Lemma 3.3.1. Let $1 \leq k \leq r+2$. Then for large $n$,

$$
\begin{equation*}
\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)=\sum_{j=0}^{\lfloor k / 2\rfloor} a_{k, j} n^{j}+\mathcal{O}\left(\epsilon_{0}^{n}\right) . \tag{3.36}
\end{equation*}
$$

Proof. We first note that taking the $k$ th derivative of (3.35) at $t=0$,

$$
i^{k} \mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)=\left.\frac{d^{k}}{d t^{k}}\right|_{t=0}\left[e^{-n t \mu^{\prime}(0)} \mu(t)^{n} Z(t)\right]+\bar{U}^{(k)}(0)
$$

$$
=\left.\frac{d^{k}}{d t^{k}}\right|_{t=0}\left[e^{-n t \mu^{\prime}(0)} \mu(t)^{n} Z(t)\right]+\mathcal{O}\left(\epsilon_{0}^{n}\right)
$$

Observe that all the derivatives of $e^{-n t \mu^{\prime}(0)} \mu(t)^{n} Z(t)$ will only have positive integral powers of $n$ (possibly) up to order $k$. Therefore, $\left.\frac{d^{k}}{d t^{k}}\right|_{t=0}\left[e^{-n t \mu^{\prime}(0)} \mu(t)^{n} Z(t)\right]=$ $\sum_{j=0}^{k} a_{k, j} n^{j}$. We claim that for $j>k / 2, a_{k, j}=0$. This claim proves the result.

We notice that the first derivative of $e^{-t \mu^{\prime}(0)} \mu(t)$ at $t=0$ is 0 . Thus we prove the more general claim that if $g(0)=1$ and $g^{\prime}(0)=0$ then $\left.\frac{d^{k}}{d t^{k}}\right|_{t=0}\left[g(t)^{n} Z(t)\right]$ has no terms with powers of $n$ greater than $k / 2$. From the Leibniz rule,

$$
\left.\frac{d^{k}}{d t^{k}}\right|_{t=0}\left[g(t)^{n} Z(t)\right]=\left.\sum_{l=0}^{k}\binom{k}{l} Z^{(k-l)}(0) \frac{d^{l}}{d t^{l}}\right|_{t=0}\left[g(t)^{n}\right] .
$$

Therefore it is enough to prove that $\left.\frac{d^{l}}{d t^{l}}\right|_{t=0}\left[g(t)^{n}\right]$ has no powers of $n$ greater than $l / 2$.

To this end we use the order $l$ Taylor expansion of $g(t)$ about $t=0$. Since $g^{\prime}(0)=0$ and $g$ is $r+2$ times continuously differentiable for $l \leq r+2$ there exists $\phi(t)$ continuous such that,

$$
\begin{aligned}
g(t) & =1+a_{2} t^{2}+\cdots+a_{l} t^{l}+t^{l+1} \phi(t) \\
\Longrightarrow g(t)^{n} & =\sum_{k_{0}+k_{2}+\cdots+k_{l+1}=n} \frac{n!}{k_{0}!k_{2}!\ldots k_{l+1}!}\left(a_{2} t^{2}\right)^{k_{2}} \ldots t^{(l+1) k_{l+1}} \phi(t)^{k_{l+1}} \\
& =\sum_{k_{0}+k_{2}+\cdots+k_{l+1}=n} \frac{C_{k_{0} k_{2} \ldots k_{l+1} n!}^{k_{0}!k_{2}!\ldots k_{l+1}!} t^{2 k_{2}+\cdots+(l+1) k_{l+1}} \phi(t)^{k_{l+1}} .}{}
\end{aligned}
$$

After combining and rearranging terms according to powers of $t$, we can obtain the order $l$ Taylor expansion of $g(t)^{n}$. Notice that if $k_{l+1} \geq 1$ then $2 k_{2}+\cdots+(l+$ 1) $k_{l+1} \geq l+1$. Terms with $k_{l+1} \geq 1$ are part of the error term of the order $l$ Taylor expansion of $g(t)^{n}$. Since our focus is on the derivative at $t=0$, the
only terms that matter are terms with $k_{l+1}=0$ and $2 k_{2}+\cdots+l k_{l}=l$. This implies that $k_{2}+\cdots+k_{l} \leq \frac{l}{2}$. Because $k_{i}$ 's are non-negative integers, this means $k_{2}+\cdots+k_{l} \leq\left\lfloor\frac{l}{2}\right\rfloor$. Hence, $k_{0} \geq n-\left\lfloor\frac{l}{2}\right\rfloor$.

This analysis shows that the largest contribution to $\left.\frac{d^{l}}{d t^{l}}\right|_{t=0}\left[g(t)^{n}\right]$ comes from the term,

$$
\frac{C_{\left(n-\left\lfloor\frac{l}{2}\right\rfloor\right), 1, \ldots, 1,0, \ldots, 0} n!}{\left(n-\left\lfloor\frac{l}{2}\right\rfloor\right)!} t^{l}
$$

whose $k$ th derivative at 0 is,

$$
\frac{C_{\left(n-\left\lfloor\frac{l}{2}\right\rfloor\right), 1, \ldots, 1,0, \ldots, 0} l!n!}{\left(n-\left\lfloor\frac{l}{2}\right\rfloor\right)!}=C_{\left(n-\left\lfloor\frac{l}{2}\right\rfloor\right), 1, \ldots, 1,0, \ldots, 0} l!n \ldots\left(n-\left\lfloor\frac{l}{2}\right\rfloor+1\right)=\mathcal{O}\left(n^{\left.n \frac{l}{2}\right\rfloor}\right)
$$

Therefore,

$$
\left.\frac{d^{l}}{d t^{l}}\right|_{t=0}\left[g(t)^{n}\right]=\mathcal{O}\left(n^{\left\lfloor\frac{l}{2}\right\rfloor}\right)
$$

It is immediate from the proof that the coefficients $a_{k, j}$ are determined by the derivatives of $\mu(t)$ and $Z(t)$ near 0 . For example, the constant term $a_{k, 0}=$ $(-i)^{k} Z^{(k)}(0)$. This follows from these three facts. The expansion (3.36) is the $k$ th derivative of the product of the three functions $e^{-n t \mu^{\prime}(0)}, \mu(t)^{n}$ and $Z(t)$ at $t=0$. All derivatives of $\mu(t)^{n}$ and $e^{-n t \mu^{\prime}(0)}$ at $t=0$ contain powers of $n$ and thus, $a_{k, 0}$ corresponds to the term $Z(t)$ being differentiated $k$ times in the Leibneiz rule. Both $e^{-n t \mu^{\prime}(0)}$ and $\mu(t)^{n}$ are 1 at $t=0$. We will see later that the other coefficients $a_{k, j}$ are combinations of $\mu^{\prime}(0)=i A$, higher order derivatives of $\mu$ at 0 upto order $k$ and derivatives of $Z$ at 0 upto order $k-1$.

As a corollary to Lemma 3.3.1, we conclude that asymptotic moments of orders upto $r+2$ exist. These provide us an alternative way to describe $a_{k, j}$.

Corollary 3.3.2. For all $1 \leq m \leq r+2$ and $0 \leq j \leq \frac{m}{2}$,

$$
a_{m, j}=\lim _{n \rightarrow \infty} \frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{m}\right)-n^{j+1} a_{m, j+1}-\cdots-n^{\left\lfloor\frac{m}{2}\right\rfloor} a_{m,\left\lfloor\frac{m}{2}\right\rfloor}}{n^{j}} .
$$

Proof. When $m=1, \mathbb{E}\left(\left[S_{n}-n A\right]\right)=a_{1,0}+\mathcal{O}\left(\epsilon_{0}^{n}\right)$ and it is immediate that $a_{1,0}=$ $\lim _{n \rightarrow \infty} \mathbb{E}\left(\left[S_{n}-n A\right]\right)$. For arbitrary $k$ we have,

$$
\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)=a_{k,\lfloor k / 2\rfloor} n^{\lfloor k / 2\rfloor}+a_{k,\lfloor k / 2\rfloor-1} n^{\lfloor k / 2\rfloor-1}+\cdots+a_{k, 0}+\mathcal{O}\left(\epsilon_{0}^{n}\right)
$$

and dividing by $n$ we obtain,

$$
\frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)}{n^{\lfloor k / 2\rfloor}}=a_{k,\lfloor k / 2\rfloor}+\mathcal{O}\left(\frac{1}{n}\right)
$$

Now, it is immediate that,

$$
a_{k,\lfloor k / 2\rfloor}=\lim _{n \rightarrow \infty} \frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)}{n^{\lfloor k / 2\rfloor}} .
$$

Having computed $a_{k, j}$, for $r \leq j \leq\left\lfloor\frac{k}{2}\right\rfloor$, we can write,

$$
\left.\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)-a_{k,\lfloor k / 2\rfloor}\right\rfloor^{\lfloor k / 2\rfloor}-\cdots-a_{k, r} n^{r}=a_{k, r-1} n^{r-1}+\cdots+a_{k, 0}+\mathcal{O}\left(\epsilon_{0}^{n}\right) .
$$

Dividing by $n^{r-1}$, we obtain,

$$
\frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)-n^{r} a_{k, r}-\cdots-n^{\lfloor k / 2\rfloor} a_{k,\lfloor k / 2\rfloor}}{n^{r-1}}=a_{k, r-1}+\mathcal{O}\left(\frac{1}{n}\right) .
$$

Now, we can compute $a_{m+1, r-1}$,

$$
a_{k, r-1}=\lim _{n \rightarrow \infty} \frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{k}\right)-n^{r} a_{k, r}-\cdots-n^{\lfloor k / 2\rfloor} a_{k,\lfloor k / 2\rfloor}}{n^{r-1}} .
$$

This proves the Corollary for arbitrary $k \in\{1, \ldots, r+2\}$.

Because the coefficients of polynomials $A_{p}(t)$ (see (3.10)) are combinations of derivatives of $\mu(t)$ and $Z(t)$ at $t=0$, we can write them explicitly in terms of $a_{k, j}$, and hence, by applying Corollary 3.3.2, the coefficients of Edgeworth polynomials can be expressed in terms of moments of $S_{n}$. Next, we will introduce a recursive algorithm to do this and illustrate the process by computing the first and second Edgeworth polynomials.

Taking the first derivative of (3.35) at $t=0$,

$$
i \mathbb{E}\left(\left[S_{n}-n A\right]\right)=Z^{\prime}(0)+\bar{U}_{n}^{\prime}(0)
$$

Then,

$$
a_{1,0}=\lim _{n \rightarrow \infty} \mathbb{E}\left(\left[S_{n}-n A\right]\right)=-i Z^{\prime}(0)
$$

Next, taking the second derivative of (3.35) at $t=0$ we have,

$$
i^{2} \mathbb{E}\left(\left[S_{n}-n A\right]^{2}\right)=n\left[\mu^{\prime \prime}(0)-\mu^{\prime}(0)^{2}\right]+Z^{\prime \prime}(0)+\bar{U}_{n}^{\prime \prime}(0)
$$

Therefore, dividing by $n$ and taking the limit we have,

$$
\begin{equation*}
a_{2,1}=\sigma^{2}=\lim _{n \rightarrow \infty} \mathbb{E}\left(\left[\frac{S_{n}-n A}{\sqrt{n}}\right]^{2}\right)=\mu^{\prime}(0)^{2}-\mu^{\prime \prime}(0) \tag{3.37}
\end{equation*}
$$

Once we have found $a_{2,1}$ we can find

$$
a_{2,0}=\lim _{n \rightarrow \infty}\left(\mathbb{E}\left(\left[S_{n}-n A\right]^{2}\right)-n \sigma^{2}\right)=-Z^{\prime \prime}(0)
$$

We can repeat this procedure iteratively. For example, after we compute the 3rd derivative of (3.35) at $t=0$ :

$$
i^{3} \mathbb{E}\left(\left[S_{n}-n A\right]^{3}\right)=Z^{(3)}(0)+n \mu^{\prime}(0)\left[2 \mu^{\prime}(0)^{2}-3 \mu^{\prime \prime}(0)\right]+n \mu^{(3)}(0)
$$

$$
+3 n Z^{\prime}(0)\left[\mu^{\prime}(0)^{2}-\mu^{\prime \prime}(0)\right]+\bar{U}_{n}^{(3)}(0)
$$

we get that,

$$
\begin{aligned}
a_{3,1}=\lim _{n \rightarrow \infty} \frac{1}{n} \mathbb{E}\left(\left[S_{n}-n A\right]^{3}\right) & =-A\left(3 \sigma^{2}+A^{2}\right)+i \mu^{(3)}(0)-3 i \sigma^{2} Z^{\prime}(0) \\
& =-A\left(3 \sigma^{2}+A^{2}\right)+i \mu^{(3)}(0)+3 \sigma^{2} a_{1,0}
\end{aligned}
$$

This gives us $\mu^{(3)}(0)$ and $Z^{(3)}(0)$ in terms of asymptotics of moments of $S_{n}$ :

$$
\begin{aligned}
& i \mu^{(3)}(0)=a_{3,1}+A\left(3 \sigma^{2}+A^{2}\right)-3 \sigma^{2} a_{1,0} \\
& i Z^{(3)}(0)=\lim _{n \rightarrow \infty}\left(\mathbb{E}\left(\left[S_{n}-n A\right]^{3}\right)-n a_{3,1}\right)
\end{aligned}
$$

Given that we have all the coefficients $a_{k, j}, 1 \leq k \leq m$ computed and $\mu^{(k)}(0), Z^{(k)}(0)$ for $1 \leq k \leq m$ expressed in terms of the former, we can compute $a_{m+1, j}$ and express $\mu^{(m+1)}(0), Z^{(m+1)}(0)$ in terms of $a_{k, j}, 1 \leq k \leq m+1$.

To see this note that $\mu^{(m+1)}(0)$ appears only as a result of $\mu^{n}(t)$ being differentiated $m+1$ times. So, $\mu^{(m+1)}(0)$ only appears in derivatives of order $m+1$ and higher. It is also easy to see that it appears in the form $n \mu^{(m+1)}(0)$ in the $(m+1)$ th derivative of (3.35). Thus, it is a part of $a_{m+1,1}$ and all the other terms in $a_{m+1,1}$ are products of $\mu^{(k)}(0), Z^{(k)}(0)$ for $1 \leq k \leq m$ whose orders add upto $m+1$ and hence they are products of $a_{k, j}, 1 \leq k \leq m$.

Also, $Z^{m+1}(0)$ appears only in $a_{m+1,0}$. This is because $Z^{m+1}(0)$ appears only as a result of $Z(t)$ being differentiated $m+1$ times. Thus, it appears only in derivatives of (3.35) of order $m+1$ or higher. In the $(m+1)$ th derivative of (3.35), there is only one term containing $Z^{(m+1)}(t)$ and it is $e^{-n t \mu^{\prime}(0)} \mu(t)^{n} Z^{m+1}(t)$. So $a_{m+1,0}=(-i)^{m+1} Z^{m+1}(0)$.

Using Corollary 3.3.2, we have,

$$
a_{m+1,\left\lfloor\frac{m+1}{2}\right\rfloor}=\lim _{n \rightarrow \infty} \frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{m+1}\right)}{n^{\left\lfloor\frac{m+1}{2}\right\rfloor}} .
$$

Having computed $a_{m+1, j}$, for $r \leq j \leq\left\lfloor\frac{m+1}{2}\right\rfloor$, we compute $a_{m+1, r-1}$ :

$$
a_{m+1, r-1}=\lim _{n \rightarrow \infty} \frac{\mathbb{E}\left(\left[S_{n}-n A\right]^{m+1}\right)-n^{r} a_{m+1, r}-\cdots-n^{\left\lfloor\frac{m+1}{2}\right\rfloor} a_{m+1,\left\lfloor\frac{m+1}{2}\right\rfloor}}{n^{r-1}} .
$$

This gives us $Z^{(m+1)}(0)=i^{m+1} a_{m+1,0}$ and $\mu^{m+1}(0)$ in terms of $a_{m+1,1}$ and $a_{k, j}$, $1 \leq k \leq m$ i.e. explicitly in terms of moments of $S_{n}$. Proceeding inductively we can compute all the derivatives upto order $r$ of $\mu(t)$ and $Z(t)$ at $t=0$ in this manner by taking derivatives up to order $r$ of (3.35) at $t=0$. This is possible because our assumptions guarantee the existence of the first $r+2$ derivatives of (3.35) near $t=0$.

Remark 3.3.1. This representation of $\mu^{(k)}(0)$ and $Z^{(k)}(0)$ in terms of $a_{k, j}$ is not unique. However, it is convenient to choose the $a_{k, j}$ 's with the lowest possible indices. The inductive procedure explained above yields exactly this representation.

We will illustrate how the first and the second order Edgeworth expansion can be computed explicitly once we have $\mu^{(4)}(0), \mu^{(3)}(0), Z^{\prime \prime}(0)$ and $Z^{\prime}(0)$ in terms of asymptotic moments of $S_{n}$. Because $A_{0}(t)=1$ we have $R_{0}(t)=1$. From the derivation of (3.9) we have,

$$
\begin{aligned}
A_{1}(t)=(\log \mu)^{(3)}(0) \frac{t^{3}}{6}-Z^{\prime}(0) t & =\left(\mu^{(3)}(0)-3 \mu^{\prime \prime}(0) \mu^{\prime}(0)+2 \mu^{\prime}(0)^{3}\right) \frac{t^{3}}{6}-Z^{\prime}(0) t \\
& =\left(\mu^{(3)}(0)+i A\left(3 \sigma^{2}+A^{2}\right)\right) \frac{t^{3}}{6}-Z^{\prime}(0) t \\
& =\left(a_{3,1}-3 \sigma^{2} a_{1,0}\right) \frac{(i t)^{3}}{6}-a_{1,0}(i t)
\end{aligned}
$$

After taking the inverse Fourier transform as shown in (3.15) we have,

$$
R_{1}(x)=\frac{\left(a_{3,1}-3 \sigma^{2} a_{1,0}\right)}{6 \sigma^{6}} x\left(3 \sigma^{2}-x^{2}\right)+\frac{a_{1,0}}{\sigma^{2}} x
$$

Using (3.16) we obtain the first Edgeworth polynomial,

$$
P_{1}(x)=\frac{\left(a_{3,1}-3 \sigma^{2} a_{1,0}\right)}{6 \sigma^{4}}\left(\sigma^{2}-x^{2}\right)-\frac{a_{1,0}}{\sigma} .
$$

Similar calculations give us,

$$
\begin{aligned}
A_{2}(t)=\left(a_{3,1}\right. & \left.+3 \sigma^{2} a_{1,0}\right)^{2} \frac{(i t)^{6}}{72}+\left[A^{2}\left(6 \sigma^{2}+A^{4}\right)+4 a_{3,1}\left(A-2 a_{1,0}\right)\right. \\
& \left.-3 \sigma^{2}\left(2 a_{2,0}-4 A a_{1,0}+\sigma^{2}\right)+a_{4,1}\right] \frac{(i t)^{4}}{24}+\left(2 a_{1,0}^{2}-a_{2,0}\right) \frac{(i t)^{2}}{2}
\end{aligned}
$$

From (3.15) and (3.16) we have,

$$
\begin{aligned}
& R_{2}(t)=\left(a_{3,1}+3 \sigma^{2} a_{1,0}\right)^{2} \frac{x^{6}-15 \sigma^{2} x^{4}+45 \sigma^{4} x^{2}-15 \sigma^{6}}{72 \sigma^{12}} \\
& +\left[A^{2}\left(6 \sigma^{2}+A^{4}\right)+4 a_{3,1}\left(A-2 a_{1,0}\right)-3 \sigma^{2}\left(2 a_{2,0}-4 A a_{1,0}+\sigma^{2}\right)+a_{4,1}\right] \\
& \times \frac{\left(x^{4}-6 \sigma^{2} x^{2}+3 \sigma^{2}\right)}{24 \sigma^{8}}+\left(2 a_{1,0}^{2}-a_{2,0}\right) \frac{\left(x^{2}-\sigma^{2}\right)}{2 \sigma^{4}}, \\
& P_{2}(t)=\left(a_{3,1}+3 \sigma^{2} a_{1,0}\right)^{2} \frac{x\left(15 \sigma^{2}-10 \sigma^{2} x^{2}+x^{6}\right)}{72 \sigma^{10}} \\
& +\left[A^{2}\left(6 \sigma^{2}+A^{4}\right)+4 a_{3,1}\left(A-2 a_{1,0}\right)-3 \sigma^{2}\left(2 a_{2,0}-4 A a_{1,0}+\sigma^{2}\right)+a_{4,1}\right] \\
& \times \frac{x\left(3 \sigma^{2}-x^{2}\right)}{24 \sigma^{6}}+\left(2 a_{1,0}^{2}-a_{2,0}\right) \frac{x}{2 \sigma^{2}} .
\end{aligned}
$$

Remark 3.3.2. Once we have $R_{p}$ for $p \in \mathbb{N}_{0}$ and $P_{p}$ for $p \in \mathbb{N}$, the polynomials $P_{p, g}, P_{p, d}$ and $P_{p, a}$ are given by $P_{p, g}=P_{p, d}=R_{p}$ and $P_{p, a}=P_{p}$. These relations were obtained in the proofs in section 3.2.

Also, one can compute $P_{p, l}$ using (3.28):

$$
P_{p, l}(x)=\sum_{l+j=2 p} \frac{(-i x)^{j}}{j!} \int t^{j} A_{l}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t
$$

For example,

$$
\begin{gathered}
P_{0, l}(x)=\int A_{0}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t=\sqrt{\frac{2 \pi}{\sigma^{2}}} \\
P_{1, l}(x)=\int A_{2}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t-i x \int t A_{1}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t-\frac{x^{2}}{2} \int t^{2} A_{0}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t \\
\frac{P_{1, l}(x)}{\sqrt{2 \pi}}=\left(a_{3,1}+3 \sigma^{2} a_{1,0}\right)^{2} \frac{5}{24 \sigma^{7}} \\
\\
+\left[A^{2}\left(6 \sigma^{2}+A^{4}\right)+4 a_{3,1}\left(A-2 a_{1,0}\right)-3 \sigma^{2}\left(2 a_{2,0}-4 A a_{1,0}+\sigma^{2}\right)+a_{4,1}\right] \frac{1}{8 \sigma^{5}} \\
\\
\quad-\left(2 a_{1,0}^{2}-a_{2,0}\right) \frac{1}{2 \sigma^{6}}-\left(\left(a_{3,1}-3 \sigma^{2} a_{1,0}\right) \frac{1}{\sigma^{5}}+\frac{2 a_{1,0}}{\sigma^{3}}\right) \frac{x}{2}-\frac{x^{2}}{2 \sigma^{3}}
\end{gathered}
$$

Higher order Edgeworth polynomials can be computed similarly.
We can compare our results with the centered i.i.d. case. Then, we have that $A=0, a_{1,0}=0$ because the sequence is stationary. Also, $a_{3,1}=\lim _{n \rightarrow \infty} \frac{1}{n} \mathbb{E}\left(\left[S_{n}-\right.\right.$ $\left.n A]^{3}\right)=\mathbb{E}\left(\left(X_{1}-A\right)^{3}\right), a_{2,0}=0$ and $a_{4,1}=\mathbb{E}\left(X_{1}^{4}\right)$. So, the above polynomials reduce to,

$$
\begin{aligned}
A_{1}(t) & =\frac{\mathbb{E}\left(X_{1}^{3}\right)}{6}(i t)^{3}, \quad R_{1}(x)=\frac{\mathbb{E}\left(X_{1}^{3}\right)}{6 \sigma^{6}} x\left(3 \sigma^{2}-x^{2}\right), \quad P_{1}(x)=\frac{\mathbb{E}\left(X_{1}^{3}\right)}{6 \sigma^{4}}\left(\sigma^{2}-x^{2}\right) \\
A_{2}(t) & =\mathbb{E}\left(X_{1}^{3}\right)^{2} \frac{(i t)^{6}}{72}+\left(\mathbb{E}\left(X_{1}^{4}\right)-3 \sigma^{4}\right) \frac{(i t)^{4}}{24} \\
\frac{P_{0, l}(x)}{\sqrt{2 \pi}} & =\frac{1}{\sigma}, \frac{P_{1, l}(x)}{\sqrt{2 \pi}}=\frac{\mathbb{E}\left(X_{1}^{3}\right)^{2}}{\sigma^{7}} \frac{5}{24}+\left(\frac{\mathbb{E}\left(X_{1}^{4}\right)}{\sigma^{5}}-\frac{3}{\sigma}\right) \frac{1}{8}-\frac{\mathbb{E}\left(X_{1}^{3}\right)}{\sigma^{5}} \frac{x}{2}-\frac{1}{\sigma^{3}} \frac{x^{2}}{2}
\end{aligned}
$$

These agree with the polynomials found in [20, Chapter XVI] (to see this one has to replace $x$ by $x / \sigma$ to make up for not normalizing by $\sigma$ here) and [4]. The polynomials $Q_{k}$ found in the latter are related to $P_{k, l}$ by $Q_{k}(x)=\frac{1}{2 \pi} P_{k, l}(x)$.

It is also easy to see that these agree with previous work on non-i.i.d. examples. In both [9, 29] only the first order Edgeworth polynomial is given explicitly. In [9], because the sequence is stationary and centered, we can take $A=0$ and $a_{1,0}=0$.

Also, the pressure $P(t)$ given there, corresponds to $\log \mu(t)$ here. So we recover $A_{1}(t)=P^{\prime \prime \prime}(0) \frac{(i t)^{3}}{6}$ in [9, Theorem 3]. In [29], sequence is centered but not assumed to be stationary. So $A=0$ and $a_{1,0} \neq 0$ and the asymptotic bias appears in the expansion and $A_{1}(t)=i \mu^{(3)}(0) \frac{(i t)^{3}}{6}-a_{1,0}(i t)$ which agrees with [29, Theorem 8.1]. This dependence on initial distribution corresponds to presence of $\ell$ in (3.1).

### 3.4 Applications.

### 3.4.1 Local Limit Theorem.

Existence of the Edgeworth expansion allows us to derive Local Limit Theorems (LLTs). For example see [16, Theorem 4]. Also, as direct consequences of weak global Edgeworth expansions, an LCLT comparable to the one given in [27, Chapter II], holds. In fact, a stronger version of LCLT holds true in special cases.

To make the notation simpler, we assume that the asymptotic mean of $S_{N}$ is 0. That is $A=\lim _{N \rightarrow \infty} \mathbb{E}\left(\frac{S_{N}}{N}\right)=0$.

Proposition 3.4.1. Suppose that $S_{N}$ satisfies the weak global Edgeworth expansion of order 0 for an integrable function $f \in(\mathcal{F},\|\cdot\|)$ where $\|\cdot\|$ is translation invariant. Further, assume that $|x f(x)|$ is integrable. Then,

$$
\begin{equation*}
\sqrt{N} \mathbb{E}\left(f\left(S_{N}-u\right)\right)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} e^{-\frac{u^{2}}{2 N \sigma^{2}}} \int f(x) d x+o(1) \tag{3.38}
\end{equation*}
$$

uniformly for $u \in \mathbb{R}$.

Proof. After the change of variables $z \sqrt{N} \rightarrow z$ in the RHS of the weak global

Edgeworth expansion,

$$
\begin{aligned}
\sqrt{N} \mathbb{E} & \left(f\left(S_{N}-u\right)\right) \\
& =\int \mathfrak{n}\left(\frac{z}{\sqrt{N}}\right) f(z-u) d z+\|f\| o(1) \\
& =\int\left[\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right)+(z-u) \mathfrak{n}^{\prime}\left(\frac{z_{u}}{\sqrt{N}}\right)\right] f(z-u) d z+\|f\| o(1) \\
& =\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right) \int f(z-u) d z+\frac{C}{N} \int(z-u) \mathfrak{n}\left(\frac{z_{u}}{\sqrt{N}}\right) f(z-u) d z+\|f\| o(1)
\end{aligned}
$$

Here $z_{u}$ is between $u$ and $z$ and depends continuously on $u$.
Notice that,

$$
\left|\int(z-u) \mathfrak{n}\left(\frac{z_{u}}{\sqrt{N}}\right) f(z-u) d z\right| \leq \int|(z-u) f(z-u)| d z \leq\|x f\|_{1}
$$

Therefore, after a change of variables $z-u \rightarrow z$ in the RHS,

$$
\sqrt{N} \mathbb{E}\left(f\left(S_{N}-u\right)\right)=\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right) \int f(z) d z+\max \left\{\|x f\|_{1},\|f\|\right\} o(1)
$$

as required.

In particular, the result holds for $\mathcal{F}=F_{0}^{1}$. If the order 0 weak global Edgeworth expansion holds for all $f \in F_{0}^{1}$, then we have the following corollary. We note that this is indeed the case for faster decaying $\left|\mathbb{E}\left(e^{i t S_{N}}\right)\right|$ as in Markov chains and piecewise expanding maps described in sections 3.5.3.1, 3.5.3.2 and 3.5.4.

Corollary 3.4.2. Suppose that $S_{N}$ admits the weak global Edgeworth expansion of order 0 for all $f \in F_{0}^{1}$. Then, for all $a<b$,

$$
\frac{\sqrt{N}}{(b-a)} \mathbb{P}\left(S_{N} \in(u+a, u+b)\right)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} e^{-\frac{u^{2}}{2 N \sigma^{2}}}+o(1)
$$

uniformly in $u \in \mathbb{R}$.

Proof. Fix $a<b$. It is elementary to see that there exists a sequence $f_{k} \in F_{0}^{1}$ with compact support such that $f_{k} \rightarrow 1_{(u+a, u+b)}$ point-wise and $f_{k}$ 's are uniformly bounded in $F_{1}^{1}$. This bound can be chosen uniformly in $u$, call it $C$.

Therefore, from the proof of Proposition 3.4.1, we have,

$$
\sqrt{N} \mathbb{E}\left(f_{k}\left(S_{N}-u\right)\right)=\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right) \int f_{k}(z) d z+C_{1}^{1}\left(f_{k}\right) o(1)
$$

Because $0 \leq C_{1}^{1}\left(f_{k}\right) \leq C$, taking the limit as $k \rightarrow \infty$ we conclude,

$$
\sqrt{N} \mathbb{P}\left(S_{N} \in(u+a, u+b)\right)=\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right) \int_{u+a}^{u+b} 1 d z+C o(1)
$$

and the result follows.

In fact, $u$ in the previous theorem need not be fixed. For example, for a sequence $u_{N}$ with $\frac{u_{N}}{\sqrt{N}} \rightarrow u$, we have the following:

Corollary 3.4.3. Suppose that $S_{N}$ admits the weak global Edgeworth expansion of order 0 for all $f \in F_{0}^{1}$. Let $u_{N}$ be a sequence such that $\lim _{N \rightarrow \infty} \frac{u_{N}}{\sqrt{N}}=u$. Then, for all $a<b$,

$$
\lim _{N \rightarrow \infty} \frac{\sqrt{N}}{(b-a)} \mathbb{P}\left(S_{N} \in\left(u_{N}+a, u_{N}+b\right)\right)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} e^{-\frac{u^{2}}{2 \sigma^{2}}}
$$

Now, we state the stronger version of LCLT in which we allow intervals to shrink.

Definition 8. Given a sequence $\epsilon_{N}$ in $\mathbb{R}^{+}$with $\epsilon_{N} \rightarrow 0$ as $N \rightarrow \infty$, we say that $S_{N}$ admits an LCLT for $\epsilon_{N}$ if we have,

$$
\frac{\sqrt{N}}{2 \epsilon_{N}} \mathbb{P}\left(S_{N} \in\left(u-\epsilon_{N}, u+\epsilon_{N}\right)\right)=\frac{1}{\sqrt{2 \pi \sigma^{2}}} e^{-\frac{u^{2}}{2 N \sigma^{2}}}+o(1)
$$

uniformly in $u \in \mathbb{R}$.

The next proposition gives a existence of weak global Edgeworth expansions as a sufficient condition for $S_{N}$ to admit a LCLT for a sequence $\epsilon_{N}$. Notice that existence of higher order expansions allow $\epsilon_{N}$ to decay faster. In case expansions of all orders exist, $\epsilon_{N}$ can decay at any subexponential rate.

Proposition 3.4.4. Suppose that $S_{N}$ satisfies the weak global Edgeworth expansion of order $r(\geq 1)$ for all $f \in F_{0}^{1}$. Let $\epsilon_{N}$ be a sequence of positive real numbers such that $\epsilon_{N} \rightarrow 0$ and $\epsilon_{N} N^{r / 2} \rightarrow \infty$ as $N \rightarrow \infty$. Then, $S_{N}$ admits an LCLT for $\epsilon_{N}$.

Proof. WLOG assume $\epsilon_{N}<1$ for all $N$. As in the previous proof, there exists a sequence $f_{k} \in F_{0}^{1}$ with compact support such that $f_{k} \rightarrow 1_{\left(u-\epsilon_{N}, u+\epsilon_{N}\right)}$ point-wise and $f_{k}$ 's are uniformly bounded in $F_{0}^{1}$. This bound can be chosen uniformly in $N$ and $u$, call it $C$.

Let $N \in \mathbb{N}$. Note that for all $k$,

$$
\mathbb{E}\left(f_{k}\left(S_{N}\right)\right)=\sum_{p=0}^{r} \frac{1}{N^{\frac{p}{2}}} \int P_{p, g}(z) \mathfrak{n}(z) f_{k}(z \sqrt{N}) d z+C_{0}^{1}\left(f_{k}\right) o\left(N^{-(r+1) / 2}\right)
$$

By taking the limit as $k \rightarrow \infty$ and using the fact $0 \leq C_{0}^{1}\left(f_{k}\right) \leq C$, we conclude,

$$
\mathbb{P}\left(S_{N} \in\left(u-\epsilon_{N}, u+\epsilon_{N}\right)\right)=\sum_{p=0}^{r} \frac{1}{N^{\frac{p}{2}}} \int_{\frac{u-\epsilon_{N}}{\sqrt{N}}}^{\frac{u+\epsilon_{N}}{\sqrt{N}}} P_{p, g}(z) \mathfrak{n}(z) d z+C o\left(N^{-(r+1) / 2}\right)
$$

After a change of variables $z \rightarrow \frac{z}{\sqrt{N}}$ in the $p=0$ term and divide the whole equation by $2 \epsilon_{N}$ to get,

$$
\begin{aligned}
& \frac{\sqrt{N}}{2 \epsilon_{N}} \mathbb{P}\left(S_{N} \in\left(u-\epsilon_{N}, u+\epsilon_{N}\right)\right) \\
& =\frac{1}{2 \epsilon_{N}} \int 1_{J_{N}}(z-u) \mathfrak{n}\left(\frac{z}{\sqrt{N}}\right) d z+\sum_{p=1}^{r} \frac{\sqrt{N}}{2 \epsilon_{N} N^{\frac{p}{2}}} \int_{\frac{u-\epsilon_{N}}{\sqrt{N}}}^{\frac{u+\epsilon_{N}}{\sqrt{N}}} P_{p, g}(z) \mathfrak{n}(z) d z+C o\left(\frac{1}{\epsilon_{N} N^{r / 2}}\right)
\end{aligned}
$$

where $J_{N}=\left(-\epsilon_{N}, \epsilon_{N}\right)$.

Note that for $p \geq 1$, there exists $C_{p}$ such that $\left|P_{p, g}(z) \mathfrak{n}(z)\right|<C_{p}$. Therefore,

$$
\left|\frac{\sqrt{N}}{2 \epsilon_{N} N^{\frac{p}{2}}} \int_{\frac{u-\epsilon_{N}}{\sqrt{N}}}^{\frac{u+\epsilon_{N}}{\sqrt{N}}} P_{p, g}(z) \mathfrak{n}(z) d z\right| \leq \frac{C_{p} \sqrt{N}}{2 \epsilon_{N} N^{\frac{p}{2}}} \int_{\frac{u-\epsilon_{N}}{\sqrt{N}}}^{\frac{u+\epsilon_{N}}{\sqrt{N}}} 1 d z \leq \frac{C_{p}}{N^{p / 2}}=o(1)
$$

Also, as in the proof of Proposition 3.4.1,

$$
\begin{aligned}
\frac{1}{2 \epsilon_{N}} \int 1_{J_{N}}(z-u) \mathfrak{n}\left(\frac{z}{\sqrt{N}}\right) d z=\frac{1}{2 \epsilon_{N}} \mathfrak{n}\left(\frac{u}{\sqrt{N}}\right) & \int_{u-\epsilon_{N}}^{u+\epsilon_{N}} 1 d z \\
& +\frac{C}{2 \epsilon_{N} N} \int_{u-\epsilon_{N}}^{u+\epsilon_{N}}(z-u) \mathfrak{n}\left(\frac{z_{u}}{\sqrt{N}}\right) d z
\end{aligned}
$$

Note that,

$$
\left|\frac{C}{2 \epsilon_{N} N} \int_{u-\epsilon_{N}}^{u+\epsilon_{N}}(z-u) \mathfrak{n}\left(\frac{z_{u}}{\sqrt{N}}\right) d z\right| \leq \frac{C}{2 \epsilon_{N} N} \int_{u-\epsilon_{N}}^{u+\epsilon_{N}}|z-u| d z=\frac{C \epsilon_{N}}{2 N}
$$

Therefore,

$$
\frac{1}{2 \epsilon_{N}} \int 1_{J_{N}}(z-u) \mathfrak{n}\left(\frac{z}{\sqrt{N}}\right) d z=\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right)+o(1)
$$

Combining these estimates with $\epsilon_{N} N^{r / 2} \rightarrow \infty$ we have that,

$$
\frac{\sqrt{N}}{2 \epsilon_{N}} \mathbb{P}\left(S_{N} \in\left(u-\epsilon_{N}, u+\epsilon_{N}\right)\right)=\mathfrak{n}\left(\frac{u}{\sqrt{N}}\right)+o(1)
$$

and it is straightforward from the proof that this is uniform.

Remark 3.4.1. We note that this result implies [16, Theorem 4] because existence of classical Edgeworth expansions imply the existence of the weak global Edgeworth expansion and this result is uniform in $u$.

### 3.4.2 Moderate Deviations.

While the CLT describes the typical behaviour or ordinary deviations from the mean provided by the law of large numbers, it is not sufficient to understand prop-
erties of distribution of $X_{n}$ completely. Therefore, the study of excessive deviations is important.

For example, deviations of order $n$ are called large deviations. An exponential moment condition is required for a large deviation principle to hold, even for the i.i.d. case. However, when deviations are of order $\sqrt{n \log n}$ (moderate deviations) this is not the case. We show here that a moderate deviation principle holds for $S_{N}$ under a weaker assumption than the exponential moment assumption.

It is also worth noting that moderate deviations have numerous applications in areas like statistical physics and risk analysis. For example, moderate deviations are greatly involved in the computation of Bayes risk efficiency. See [44] for details.

Proposition 3.4.5. Suppose $S_{N}$ admits the order $r$ Edgeworth expansion. Then for all $c \in(0, r)$, when $1 \leq x \leq \sqrt{c \sigma^{2} \ln N}$,

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \frac{1-\mathbb{P}\left(\frac{S_{N}-A N}{\sqrt{N}} \leq x\right)}{1-\mathfrak{N}(x)}=1 \tag{3.39}
\end{equation*}
$$

Proof. Note that,

$$
\begin{aligned}
1-\mathfrak{N}(x)-\left[1-\mathbb{P}\left(\frac{S_{N}-A N}{\sqrt{N}} \leq x\right)\right] & =\mathbb{P}\left(\frac{S_{N}-A N}{\sqrt{N}} \leq x\right)-\mathfrak{N}(x) \\
& =\sum_{p=1}^{r} \frac{P_{p}(x)}{N^{p / 2}} \mathfrak{n}(x)+o\left(N^{-r / 2}\right)
\end{aligned}
$$

uniformly in $x$. So it is enough to show that for $1 \leq x \leq \sqrt{c \sigma^{2} \ln N}$,

$$
\lim _{N \rightarrow \infty} \frac{P_{p}(x) \mathfrak{n}(x)}{N^{p / 2}(1-\mathfrak{N}(x))}=0 \text { and } \frac{N^{-r / 2}}{1-\mathfrak{N}(x)}=o(1)
$$

Note that for $x \geq 1$,

$$
1-\mathfrak{N}(x)=\frac{\sigma^{2} \mathfrak{n}(x)}{x}+\mathcal{O}\left(\frac{\mathfrak{n}(x)}{x^{3}}\right)
$$

Thus,

$$
\begin{aligned}
\frac{N^{-r / 2}}{1-\mathfrak{N}(x)} \leq \frac{N^{-r / 2}}{1-\mathfrak{N}\left(\sqrt{c \sigma^{2} \ln N}\right)} & =\mathcal{O}\left(\sqrt{\ln N} \frac{N^{-r / 2}}{e^{-\frac{c}{2} \ln N}}\right) \\
& =\mathcal{O}\left(\frac{\ln N}{N^{(r-c) / 2}}\right)
\end{aligned}
$$

Say $P_{p}(x)$ is of degree $q$. Then for some $C$ and $K$,

$$
\begin{aligned}
\left|\frac{P_{p}(x) \mathfrak{n}(x)}{N^{p / 2}(1-\mathfrak{N}(x))}\right| \leq C \frac{\left(x^{q}+K\right) \mathfrak{n}(x)}{N^{p / 2}(1-\mathfrak{N}(x))} & =C \frac{\left(x^{q}+K\right)}{N^{p / 2}} x\left(1+\mathcal{O}\left(\frac{1}{x^{2}}\right)\right) \\
& \leq C \frac{(\ln N)^{q+1}}{N^{p / 2}} \rightarrow 0 \text { as } N \rightarrow \infty
\end{aligned}
$$

This completes the proof of (3.39).

Proposition 3.4.5 is a generalization of the results on moderate deviations found in [43] to the non-i.i.d. case along with improvements on the moment condition. It should be noted that [4] contains an improvement of the moment condition for the i.i.d. case. But the proof we present here is different from the proof presented in [4].

As an immediate corollary to the above theorem, we can state the following first order asymptotic for probability of moderate deviations.

Corollary 3.4.6. Assume $S_{N}$ admits the order $r$ Edgeworth expansion. Then for all $c \in(0, r)$,

$$
\mathbb{P}\left(S_{N} \geq A N+\sqrt{c \sigma^{2} N \ln N}\right) \sim \frac{1}{\sqrt{2 \pi c}} \frac{1}{\sqrt{N^{c} \ln N}}
$$

### 3.5 Examples

Here we give several examples of systems satisfying assumptions (A1)-(A4).

### 3.5.1 Independent variables.

Let $X_{n}$ be i.i.d. with $r+2$ moments. In this case we can take $\mathbb{B}=\mathbb{R}$, and define $\mathcal{L}_{t} v=\mathbb{E}\left(e^{i t X_{1}} v\right)=\phi(t) v$ where $\phi$ is the characteristic function of $X_{1}$. Here we have taken $\ell=1$. Put $v=1$. Then, the independence of the random variables gives us, $\mathcal{L}_{t}^{n} 1=\mathbb{E}\left(e^{i t S_{n}}\right)=\phi(t)^{n}$. Also, the moment condition implies $t \rightarrow \phi(t)$ is $C^{r+2}$. This means (A1) is satisfied. (A2) is clear.

Suppose $X_{1}$ is $l$-Diophantine. That is there exists $C>0$ and $t_{0}>0$ such that for all $|t|>t_{0},|\phi(t)|<1-\frac{C}{|t|^{\mid}}$. Then $|\phi(t)| \leq e^{-\frac{C}{|t|^{2}}}$. So $|\phi(t)|<1$ for all $t \neq 0$. So we have (A3). Also, this implies that $X_{1}$ is non-lattice. An easy computation shows that when $r_{1}<\frac{1}{l}$, there exists $r_{2}$ such that $t_{0}<|t|<n^{r_{1}} \Longrightarrow|\phi(t)|^{n} \leq n^{-r_{2}}$. In fact, $|\phi(t)|^{n} \leq e^{-c n^{\alpha}}$ where $\alpha=1-r_{1} l>0$. So, (A4) is satisfied with $r_{1}<\frac{1}{l}$.

When $l=0$ we see that (A4) is satisfied with $r_{1}>\frac{r-1}{2}$. Hence, by Theorem 3.1.1 order $r$ Edgeworth expansion for $S_{n}$ exists. This is exactly the classical result of Cramér because the condition: $\limsup _{|t| \rightarrow \infty}|\phi(t)|<1$ corresponds to $l=0$.

Choose $q>\frac{r+1}{2 r_{1}}>\frac{(r+1) l}{2}$. Then, by Theorem 3.1.4 and Theorem 3.1.5 we have that $S_{n}$ admits weak global expansion for $f \in F_{0}^{q+2}$ and weak local expansion for $f \in F_{r+1}^{q+2}$. These are similar to the results appearing in [4] but slightly weaker because we require one more derivative: $q+2>2+\frac{(r+1) l}{2}$ as opposed to $1+$ $\frac{(r+1) l}{2}$. This is because we do not use the optimal conditions for the integrability of the Fourier transform. If we required $f \in F_{r}^{q+1}$ and $f^{(q+1)}$ to be $\alpha$-Hölder for small $\alpha$, then the proof would still hold true and we could recover the results in [4].

### 3.5.2 Finite state Markov chains.

Here we present a non-trivial example for which the weak Edgeworth expansions exist but the strong expansion does not exist.

Consider the Markov chain $x_{n}$ with states $S=\{1, \ldots, d\}$ whose transition probability matrix $P=\left(p_{j k}\right)_{d \times d}$ is positive. Then, by the Perron-Forbenius theorem, 1 is a simple eigenvalue of $P$ and all other eigenvalues are strictly contained inside the unit disk. Suppose $\mathbf{h}=\left(h_{j k}\right)_{d \times d} \in \mathrm{M}(d, \mathbb{R})$ and that there does not exist constants $c, r$ and a $d$-vector $H$ such that

$$
r h_{j k}=c+H(k)-H(j) \quad \bmod \quad 2 \pi
$$

for all $j, k$. Put $X_{n}=h_{x_{n} x_{n+1}}$.
For the family of operators $\mathcal{L}_{t}: \mathbb{C}^{d} \rightarrow \mathbb{C}^{d}$,

$$
\begin{equation*}
\left(\mathcal{L}_{t} f\right)_{j}=\sum_{k=1}^{d} e^{i t h_{j k}} p_{j k} f_{k}, j=1, \ldots, d \tag{3.40}
\end{equation*}
$$

$v=1$ and $\ell=\mu_{0}$, the initial distribution, we have (3.1).
Define $b_{r, j, k}=h_{r j}+h_{j k}$ for all $j, r=1, \ldots, d$ and $k=2, \ldots, d$. Put $d(s)=$ $\max \left\{\left(b_{r, j, k}-b_{r, 1, k}\right) s\right\}$ where $\{$.$\} denotes the fractional part. We further assume$ that $\mathbf{h}$ is $\beta$-Diophantine, that is, there exists $K \in \mathbb{R}$ such that for all $|s|>1$,

$$
\begin{equation*}
d(s) \geq \frac{K}{|s|^{\beta}} . \tag{3.41}
\end{equation*}
$$

If $\beta>\frac{1}{d^{2}(d-1)-1}$ then almost all $\mathbf{h}$ are $\beta$-Diophantine.
Because $S_{n}$ can take at most $\mathcal{O}\left(n^{d^{2}-1}\right)$ distinct values, $S_{n}$ has a maximal jump of order at least $n^{-\left(d^{2}-1\right)}$. Therefore, the process $X_{n}^{\mathbf{h}}=h_{x_{n} x_{n-1}}$ does not admit the order $2\left(d^{2}-1\right)$ Edgeworth expansion.

The Perron-Forbenius theorem implies that the operator $\mathcal{L}_{0}$ satisfies (A2). Because (3.40) is a finite sum, it is clear that $t \mapsto \mathcal{L}_{t}$ is analytic on $\mathbb{R}$. So we also have (A1). Also the spectral radius of $\mathcal{L}_{t}$ is at most 1 . Assume $\mathcal{L}_{t}$ has an eigenvalue on the unit circle, say $e^{i \lambda}$, with eigenvector $f$, then,

$$
e^{i \lambda} f_{j}=\left(\mathcal{L}_{t} f\right)_{j}=\sum_{k=1}^{d} e^{i t h_{j k}} p_{j k} f_{k}
$$

Assuming $\max _{j}\left|f_{j}\right|=\left|f_{r}\right|$,

$$
\left|f_{r}\right|=\left|e^{i \lambda} f_{r}\right|=\left|\sum_{k=1}^{d} e^{i t h_{j k}} p_{j k} f_{k}\right| \leq \sum_{k=1}^{d} p_{j k}\left|f_{k}\right| \Longrightarrow \sum_{k=1}^{d} p_{j k}\left(\left|f_{k}\right|-\left|f_{r}\right|\right) \geq 0
$$

Because $\left|f_{k}\right|-\left|f_{r}\right| \leq 0$ for all $k$ and $p_{j k} \geq 0$ for all $j$ and $k$ we have $\left|f_{k}\right|=\left|f_{r}\right|$ for all $k$. Therefore, there exist a $d$-vector $H$ such that $f_{k}=R e^{i H(k)}$ for all $k$. Then,

$$
\begin{aligned}
e^{i \lambda} R e^{i H(j)} & =\sum_{k=1}^{d} e^{i t h_{j k}} p_{j k} R e^{i H(k)} \\
0 & =\sum_{k=1}^{d} p_{j k}\left(e^{i\left(t h_{j k}+H(k)-H(j)-\lambda\right)}-1\right) \\
\Longrightarrow t h_{j k} & =\lambda+H(j)-H(k) \bmod 2 \pi
\end{aligned}
$$

But this is a contradiction. Therefore, (A3) holds. Next we notice that,

$$
\begin{align*}
\left|\left(\mathcal{L}_{t}^{2} f\right)_{r}\right|=\left|\sum_{j=1}^{d} \sum_{k=1}^{d} e^{i t\left(h_{r j}+h_{j k}\right)} p_{r j} p_{j k} f_{k}\right| & =\left|\sum_{k=1}^{d}\left(\sum_{j=1}^{d} e^{i t\left(h_{r j}+h_{j k}\right)} p_{r j} p_{j k}\right) f_{k}\right| \\
& \leq\|f\|\left(\sum_{k=1}^{d}\left|\sum_{j=1}^{d} e^{i t b_{r, j, k}} p_{r j} p_{j k}\right|\right) \tag{3.42}
\end{align*}
$$

Now we estimate $\left|b_{r, k}(t)\right|$ where

$$
b_{r, k}(t)=\sum_{j=1}^{d} e^{i t b_{r, j, k}} p_{r j} p_{j k}=e^{i t b_{r, 1, k}} \sum_{j=1}^{d} e^{i t\left(b_{r, j, k}-b_{r, 1, k}\right)} p_{r j} p_{j k}
$$

Then we have,

$$
\begin{aligned}
\left|b_{r, k}(t)\right|^{2} & =\sum_{j=1}^{d} p_{r j}^{2} p_{j k}^{2}+2 \sum_{j>l}^{d} p_{r j} p_{j k} p_{r l} p_{l k} \cos \left(\left(b_{r, j, k}-b_{r, l, k}\right) t\right) \\
& =\left(\sum_{j=1}^{d} p_{r j} p_{j k}\right)^{2}-2 \sum_{j>l}^{d} p_{r j} p_{j k} p_{r l} p_{l k}\left[1-\cos \left(\left(b_{r, j, k}-b_{r, l, k}\right) t\right)\right] \\
& =\left(\sum_{j=1}^{d} p_{r j} p_{j k}\right)^{2}-2 C d(t)^{2}+\mathcal{O}\left(d(t)^{3}\right), C>0 \\
\left|b_{r, k}(t)\right| & =\sum_{j=1}^{d} p_{r j} p_{j k}-\tilde{C} d(t)^{2}+\mathcal{O}\left(d(t)^{3}\right), \tilde{C}>0
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
\sum_{k=1}^{d}\left|\sum_{j=1}^{d} e^{i t b_{r, j, k}} p_{r j} p_{j k}\right| & =\sum_{k=1}^{d}\left(\sum_{j=1}^{d} p_{r j} p_{j k}\right)-\bar{C} d(t)^{2}+\mathcal{O}\left(d(t)^{3}\right) \\
& =1-\bar{C} d(t)^{2}+\mathcal{O}\left(d(t)^{3}\right), \bar{C}>0
\end{aligned}
$$

From the Diophantine condition (3.41), we can conclude that there exists $\theta>0$ such that for all $|t|>1$,

$$
\left\|\mathcal{L}_{t}^{2}\right\| \leq 1-\theta d(t)^{2} \Longrightarrow\left\|\mathcal{L}_{t}^{N}\right\| \leq\left(1-\theta d(t)^{2}\right)^{\lceil N / 2\rceil} \leq e^{-\theta d(t)^{2} N / 2} \leq e^{-\theta t^{-2 \beta} N / 2}
$$

When $1<|t|<N^{\frac{1-\epsilon}{2 \beta}}$, we have, $\left\|\mathcal{L}_{t}^{N}\right\| \leq e^{-\theta N^{\epsilon} / 2}$ which gives us (A4) with $r_{1}=\frac{1-\epsilon}{2 \beta}$ where $\epsilon>0$ can be made as small as required. Because for small $\epsilon,\left\lceil\frac{r+1}{2(1-\epsilon)}\right\rceil=$ $\left\lceil\frac{r+1}{2}\right\rceil$, choosing $q>\frac{r+1}{2} \beta$, we conclude that for $f \in F_{0}^{q+2}$ weak global and for $f \in F_{r+1}^{q+2}$ weak local Edgeworth expansions of order $r$ for the process $X_{n}^{\mathbf{h}}$ exist. Also, $S_{N}$ admits averaged Edgeworth expansions of order $r$ for $f \in F_{0}^{2}$. In the special case of $\beta>\frac{1}{d^{2}(d-1)-1}$, these hold for a full measure set of $\mathbf{h}$ even though the order $r$ strong expansion does not exist for $r+1 \geq d^{2}$.

### 3.5.3 More general Markov chains.

### 3.5.3.1 Chains with smooth transition density.

First we consider the case where $x_{n}$ is a time homogeneous Markov process on a compact connected manifold $\mathcal{M}$ with smooth transition density $p(x, y)$ which is bounded away from 0 , and $X_{n}=h\left(x_{n-1}, x_{n}\right)$ for a piece-wise smooth function $h: \mathcal{M} \times \mathcal{M} \rightarrow \mathbb{R}$. We assume that $h(x, y)$ can not be written in the form

$$
\begin{equation*}
h(x, y)=H(y)-H(x)+c(x, y) \tag{3.43}
\end{equation*}
$$

where $c(x, y)$ is piece-wise constant.
In particular, there is no constant $c$ and a function $H$ such that $h(x, y)=$ $H(y)-H(x)+c$. Also, the transition probability $P(x, d y)$ of $X_{n}$ has a non-degenrate absolute continuous component. Then, by [25], the CLT holds with $\sigma^{2}>0$.

To check the assumption 3.43 we need the following:
Lemma 3.5.1. (3.43) holds iff there exists $o \in \mathcal{M}$ such that the function $x \mapsto$ $h(o, x)+h(x, y)$ is piece-wise constant for each $y$.

Proof. If (3.43) holds then for each $o \in \mathcal{M}$

$$
h(o, x)+h(x, y)=c(o, x)+c(x, y)+H(y)-H(o)
$$

where $c(o, x)+c(x, y)$ is piece-wise constant in $x$ for each $y$.
Conversely, suppose for some $o \in \mathcal{M}, x \mapsto h(o, x)+h(x, y)$ is piece-wise constant for each $y$. Fix $y$. Let $c=h(o, o)$ and $H(x)=h(o, x)-h(o, o)$. Then, $h(o, o)+h(o, y)$ and $h(o, x)+h(x, y)$ differ by a piece-wise constant function. Then
(3.43) holds because $h(o, x)+h(x, y)-(h(o, o)+h(o, y))=h(x, y)+H(x)-H(y)-c$ is piecewise constant.

Let $\mathbb{B}=L^{\infty}(\mathcal{M})$ and consider the family of integral operators,

$$
\left(\mathcal{L}_{t} u\right)(x)=\int p(x, y) e^{i t h(x, y)} u(y) d y
$$

Let $\mu$ be the initial distribution of the Markov chain and $\left\{\mathcal{F}_{n}\right\}$ be the filtration adapted to the processes. Then, using the Markov property,

$$
\mathbb{E}_{\mu}\left[e^{i t S_{n}}\right]=\mathbb{E}_{\mu}\left[e^{i t S_{n-1}} \mathcal{L}_{t} 1\right]
$$

By induction we can conclude

$$
\mathbb{E}_{\mu}\left(e^{i t S_{n}}\right)=\int \mathcal{L}_{t}^{n} 1 d \mu
$$

Because $h$ is bounded, expanding $e^{i t h(x, y)}$ as a power series in $t$, we see that $t \mapsto \mathcal{L}_{t}$ is analytic for all $t$. This shows that (A1) is statisfied.

From the Weierstrass theorem there exist functions $q_{k}, r_{k}$ on $\mathcal{M}$ such that $p(x, y)$ is a uniform limit of functions of the form $\sum_{k=1}^{n} q_{k}(x) r_{k}(y)$. Therefore, $\mathcal{L}_{t}$ is a uniform limit of finite rank operators and is compact. Compact operators have a point spectrum hence the essential spectral radius of $\mathcal{L}_{t}$ vanishes. It is also immediate that $\left\|\mathcal{L}_{t}\right\| \leq 1$ for all $t$. Hence the spectrum is contained in the closed unit disk.

In addition, $\mathcal{L}_{0}: L^{\infty}(\mathcal{M}) \rightarrow L^{\infty}(\mathcal{M})$ given by

$$
\left(\mathcal{L}_{0} u\right)(x)=\int p(x, y) u(y) d y
$$

is a positive operator. Note that $\left(\mathcal{L}_{0} 1\right)(x)=1$ for all $x$. Thus, 1 is an eigenvalue of $\mathcal{L}_{0}$ with eigenfunction 1 . Also, eigenvalue 1 is simple and all other eigenvalues
$\beta$ are such that $|\beta|<1$. This follows from a direct application of Birkhoff Theory (see [2]). Thus, we have (A2).

Next we show that if $\beta \in \operatorname{sp}\left(\mathcal{L}_{t}\right), t \neq 0$ then $|\beta|<1$. If not, then there exists $\lambda$ and $u \in L^{\infty}(\mathcal{M})$ such that

$$
\int p(x, y) e^{i t h(x, y)} u(y) d y=e^{i \lambda} u(x)
$$

Suppose $\sup _{x}|u(x)|=R$ then for each $\epsilon>0$ there exists $x_{\epsilon}$ such that

$$
R-\epsilon \leq\left|u\left(x_{\epsilon}\right)\right|=\left|e^{i \lambda} u\left(x_{\epsilon}\right)\right|=\left|\int p(x, y) e^{i t h(x, y)} u(y) d y\right| \leq \int p(x, y)|u(y)| d y
$$

Therefore,

$$
\int p(x, y)[|u(y)|-R] d y \geq-\epsilon
$$

But $|u(y)|-R \leq 0$. Hence, $|u(y)|=R$ a.e. Therefore, $u(y)=R e^{i \theta(y)}$ a.e. for some function $\theta$ and we may assume $\theta \in[0,2 \pi)$.

$$
\begin{align*}
& \int p(x, y) e^{i t h(x, y)} R e^{i \theta(y)} d y=R e^{i \lambda} e^{i \theta(x)} \\
\Longrightarrow & \int p(x, y)\left[e^{i(t h(x, y)-\lambda+\theta(y)-\theta(x))}-1\right] d y=0 \\
\Longrightarrow & t h(x, y)-\lambda+\theta(y)-\theta(x) \equiv 0 \quad \bmod 2 \pi \tag{3.44}
\end{align*}
$$

Fix $y$ and $t$. Then, for all $z, x \mapsto h(y, x)+h(x, z)$ does not depend on $x$ modulo $2 \pi$ i.e. it is piece-wise constant for all $t \neq 0$. By Lemma 3.5.1, $h(x, y)$ satisfies (3.43). This contradiction proves (A3).

Recall that if $\mathcal{K}$ is integral operator

$$
(\mathcal{K} u)(x)=\int k(x, y) u(y) d y
$$

then

$$
\|\mathcal{K}\|=\sup _{x} \int|k(x, y)| d y
$$

In our case $\mathcal{L}_{t}^{2}$ has the kernel,

$$
\mathfrak{l}_{t}(x, y)=\int e^{i t[h(x, z)+h(z, y)]} p(x, z) p(z, y) d z
$$

By Lemma 3.5.1 for each $x$ and $y$ the function $z \mapsto(h(x, z)+h(z, y))$ is not piecewise constant. So its derivative (whenever it exists) is not identically 0 . Thus there is an open set $V_{x, y}$ and a vector field $e$ such that $\partial_{e}[h(x, z)+h(z, y)] \neq 0$ on $V_{x, y}$. Integrating by parts in the direction of $e$ we conclude that

$$
\lim _{t \rightarrow \infty} \int_{V_{x, y}} e^{i t[h(x, z)+h(z, y)]} p(x, z) p(z, y) d z=0
$$

By compactness there are constants $r_{0}, \varepsilon_{0}$ such that for $|t| \geq r_{0}$ and all $x$ and $y$ in $\mathcal{M},\left|\mathfrak{l}_{t}(x, y)\right| \leq \mathfrak{l}_{0}(x, y)-\varepsilon_{0}$. It follows that

$$
\begin{equation*}
\left\|\mathcal{L}_{t}^{2}\right\|=\sup _{x} \int_{\mathcal{M}}\left|\mathfrak{l}_{t}(x, y)\right| d y \leq \int_{\mathcal{M}} \mathfrak{l}_{0}(x, y) d y-\varepsilon_{0} \tag{3.45}
\end{equation*}
$$

The first term here equals

$$
\iint_{\mathcal{M} \times \mathcal{M}} p(x, z) p(z, y) d z d y=1
$$

Hence for $|t| \geq r_{0},\left\|\mathcal{L}_{t}^{2}\right\| \leq 1-\varepsilon_{0}$ and so $\left\|\mathcal{L}_{t}^{N}\right\| \leq\left(1-\varepsilon_{0}\right)^{\lceil N / 2\rceil}$. This proves (A4) with no restriction on $r_{1}$. Therefore, $S_{N}$ admits Edgeworth expansions of all orders.

Next we look at the case when (3.43) fails but the constants are not lattice valued. Then, arguments for (A1), (A2) and (A3) hold. In particular, (3.44) cannot
hold since it implies that

$$
\left(h(x, y)+\frac{\theta(y)}{t}-\frac{\theta(x)}{t}\right) \in \frac{\lambda}{t}+\frac{2 \pi}{t} \mathbb{Z}
$$

However, we have to impose a Diophantine condition on the values that $h(x, y)$ can take in order to obtain a sufficient control over $\left\|\mathcal{L}_{t}^{N}\right\|$ and obtain (A4).

For fixed $x, y$ let the range of $z \mapsto h(x, z)+h(z, y)$ be $S=\left\{c_{1}, \ldots, c_{d}\right\}$. Note that these $c_{i}$ 's may depend on $x$ and $y$. However, there can be at most finitely many values that $h(x, z)+h(z, y)$ can take as $x$ and $y$ vary on $\mathcal{M}$ because $h$ is piece-wise smooth. So we might as well assume that $S$ is this complete set of values. Also, take $U_{k}$ to be the open set on which $z \mapsto h(x, z)+h(z, y)$ takes value $c_{k}$. Take $b_{k}=c_{k}-c_{1}$ and define $d(s)=\max \left\{b_{k} s\right\}$. Assume further that there exists $K>0$ such that for all $|s|>1$,

$$
d(s) \geq \frac{K}{|s|^{\beta}}
$$

If $\beta>(d-1)^{-1}$ for almost all $d$-tuples $\mathbf{c}=\left(c_{1}, \ldots, c_{d}\right)$, the above holds.
Note that,

$$
\begin{aligned}
\left|\mathcal{L}_{t}^{2} u(x)\right| & =\int\left|\int e^{i t[h(x, z)+h(z, y)]} p(x, z) p(z, y) d z\right||u(y)| d y \\
& \leq\|u\| \int\left|\sum_{k=1}^{d} e^{i t c_{k}} \int_{U_{k}} p(x, z) p(z, y) d z\right| d y=\|u\| \int\left|\sum_{k=1}^{d} p_{k} e^{i t b_{k}}\right| d y
\end{aligned}
$$

where and $p_{k}=\int_{U_{k}} p(x, z) p(z, y) d z$. Therefore, $p_{1}+\cdots+p_{d}=p(x, y)$.
Now the situation is similar to that of (3.42) and a similar calculation yields,

$$
\left|\sum_{k=1}^{d} p_{k} e^{i t b_{k}}\right|=p(x, y)-C d(t)^{2}+\mathcal{O}\left(d(t)^{3}\right), C>0
$$

Therefore,

$$
\left\|\mathcal{L}_{t}^{2}\right\| \leq \int\left[p(x, y)-C d(t)^{2}+\mathcal{O}\left(d(t)^{3}\right)\right] d y=1-\tilde{C} d(s)^{2}
$$

From this we can repeat the analysis done in the finite state Markov chains example following (3.42). In particular, when $1<|t|<N^{\frac{1-\epsilon}{2 \beta}}$, there exists $\theta>0$ such that $\left\|\mathcal{L}_{t}^{N}\right\| \leq e^{-\theta N^{\epsilon}}$ which gives us (A4).

Finally, when (3.43) fails and $h$ takes integer values with span $1, X_{n}$ is a lattice random variable and we can discuss the existence of the lattice Edgeworth expansion. In this case $S_{N}$ admits the lattice expansion of all orders. To this end, only the condition $\widetilde{(\mathrm{A} 3)}$ needs to be checked. First note that $\mathcal{L}_{0}=\mathcal{L}_{2 \pi k}$ for all $k \in \mathbb{Z}$. Also, assuming $\mathcal{L}_{t}$ has an eigenvalue on the unit circle, we conclude (3.44),

$$
t h(x, y)-\lambda+\theta(y)-\theta(x) \equiv 0 \quad \bmod 2 \pi
$$

This implies $t(h(x, y)+h(y, x)) \in 2 \pi \mathbb{Z}+2 \lambda$. Note that LHS belongs a lattice with span $t$ and RHS is a lattice with span $2 \pi$. Because $t$ is not a multiple of $2 \pi$ this equality cannot happen. Therefore, when $t \notin 2 \pi \mathbb{Z}, \operatorname{sp}\left(\mathcal{L}_{t}\right) \subset\{|z|<1\}$ and we have the claim.

### 3.5.3.2 Chains without densities.

We consider a more general case where transition probabilities may not have a density. We claim we can recover (A1)-(A4) if the transition operator takes the form

$$
\mathcal{L}_{0}=a \mathcal{J}_{0}+(1-a) \mathcal{K}_{0}
$$

where $a \in(0,1)$ and $\mathcal{J}_{0}$ and $\mathcal{K}_{0}$ are Markov operators on $L^{\infty}(\mathcal{M})$ (i.e. $\mathcal{J}_{0} f \geq 0$ if $f \geq 0$ and $\mathcal{J}_{0} 1=1$ and similarly for $\left.\mathcal{K}_{0}\right)$,

$$
\mathcal{J}_{0} f(x)=\int p(x, y) f(y) d \mu(y)
$$

and

$$
\mathcal{K}_{0} f(x)=\int f(y) Q(x, d y)
$$

where $p$ is a smooth transition density and $Q$ is a transition probability measure. Let $h(x, y)$ be piece-wise smooth and put,

$$
\mathcal{J}_{t}(f)=\mathcal{J}_{0}\left(e^{i t h} f\right) \text { and } \mathcal{K}_{t}(f)=\mathcal{K}_{0}\left(e^{i t h} f\right)
$$

Defining $\mathcal{L}_{t}=a \mathcal{J}_{t}+(1-a) \mathcal{K}_{t}$ we can conclude $t \mapsto \mathcal{L}_{t}$ is analytic and that

$$
\mathbb{E}_{\mu}\left(e^{i t S_{n}}\right)=\int \mathcal{L}_{t}^{n} 1 d \mu
$$

Now we show that conditions (A2), (A3) and (A4) are satisfied. Because $\left\|\mathcal{J}_{t}\right\| \leq 1$ and $\left\|\mathcal{K}_{t}\right\| \leq 1$ we have $\left\|\mathcal{L}_{t}\right\| \leq 1$. Thus the spectral radius of $\mathcal{L}_{t}$ is $\leq 1$. Because $a \mathcal{J}_{t}$ is compact, $\mathcal{L}_{t}$ and $(1-a) \mathcal{K}_{t}$ have the same essential spectrum. See [33, Theorem IV.5.35]. However the spectral radius of the latter is at most $(1-a)$. Hence, the essential spectral radius of $\mathcal{L}_{t}$ is at most $(1-a)$.

Because both $\mathcal{J}_{0}$ and $\mathcal{K}_{0}$ are Markov operators we can conclude that 1 is an eigenvalue of $\mathcal{L}_{0}$ with constant function 1 as the corresponding eigenfunction. From the previous paragraph the essential spectral radius of $\mathcal{L}_{0}$ is at most $(1-a)$. Because $\mathcal{L}^{n}$ is norm bounded it cannot have Jordan blocks. So 1 is semisimple.

Suppose, $\mathcal{L}_{t} u=e^{i \theta} u$. Without loss of generality we may assume $\|u\|_{\infty}=1$. Assuming there exists a positive measure set $\Omega$ with $|u(x)|<1-\delta$ we can conclude that, for all $x$,

$$
\begin{aligned}
|u(x)|=\left|L_{t} u(x)\right| & =\left|a \mathcal{J}_{t} u(x)+(1-a) \mathcal{K}_{t} u(x)\right| \\
& \leq a \int_{\Omega}|u(y)| p(x, y) d \mu(y)+a \int_{\Omega^{c}}|u(y)| p(x, y) d \mu(y)+(1-a)
\end{aligned}
$$

$$
\leq 1-a \delta \mu(\Omega)
$$

This is a contradiction. Therefore, $|u(x)|=1$. Put $u(x)=e^{i \gamma(x)}$. Then,

$$
1=a \int e^{i(t h(x, y)+\gamma(y)-\gamma(x)-\theta)} p(x, y) d \mu(y)+(1-a) e^{-i(\theta+\gamma(x))} \mathcal{K}_{t} u
$$

Hence, $\int e^{i(t h(x, y)+\gamma(y)-\gamma(x)-\theta)} p(x, y) d \mu(y)=1 \Longrightarrow \mathcal{J}_{t} u=e^{i \theta} u$. From section 3.5.3.1, this can only be true when $t=0$ and in this case $\theta=0$ and $u \equiv 1$. This concludes that $\mathcal{L}_{t}, t \neq 0$ has no eigenvalues on the unit disk and the only eigenvalue of $\mathcal{L}_{0}$ on the unit disk is 1 and its geometric multiplicity is 1 . As 1 is semisimple, it is simple as required. This concludes proof of (A2) and (A3).

From the previous case, there exists $r>0$ and $\epsilon \in(0,1)$ such that such that for all $|t|>r$ we have $\left\|\mathcal{J}_{t}^{2}\right\| \leq 1-\epsilon$. From this we have,

$$
\left\|\mathcal{L}_{t}^{2}\right\|=\left\|a^{2} \mathcal{J}_{t}^{2}+a(1-a) \mathcal{J}_{t} \mathcal{K}_{t}+(1-a) a \mathcal{K}_{t} \mathcal{J}_{t}+(1-a)^{2} \mathcal{K}_{t}^{2}\right\| \leq 1-a^{2} \epsilon
$$

Hence, for all $|t|>r$, for all $N,\left\|\mathcal{L}_{t}^{N}\right\| \leq\left(1-a^{2} \epsilon\right)^{\lfloor N / 2\rfloor}$ which gives us (A4) with no restrictions on $r_{1}$. Therefore, $S_{N}$ admits Edgeworth expansions of all orders as before.

As in the previous section, an analysis can be carried out when (3.43) fails. The conclusions are exactly the same.

### 3.5.4 One dimensional piecewise expanding maps.

Here we check assumptions (3.1), (A1)-(A4) for piecewise expanding maps of the interval using the results of $[5,37]$.

Let $f:[0,1] \rightarrow[0,1]$ be such that there is a finite partition $\mathcal{A}_{0}$ of $[0,1]$ (except possibly a measure 0 set) into open intervals such that for all $I \in \mathcal{A}_{0},\left.f\right|_{I}$ extends to a $C^{2}$ map on an interval containing $\bar{I}$. In other words $f$ is a piece-wise $C^{2}$ map. Further, assume that $f^{\prime} \geq \lambda>1$ i.e. $f$ is uniformly expanding. Next, let $\mathcal{A}_{n}=\bigvee_{k=0}^{n} T^{-j} \mathcal{A}_{0}$ and suppose for each $n$ there is $N_{n}$ such that for all $I \in \mathcal{A}_{n}$, $f^{N_{n}} I=[0,1]$. Such maps are called covering.

Statistical properties of piece-wise $C^{2}$ covering expanding maps of an interval, are well-understood. For example, see [37]. In particular, such a function $f$ has a unique absolutely continuous invariant measure with a strictly positive density $h \in \operatorname{BV}[0,1]$ and the associated transfer operator

$$
\mathcal{L}_{0} \varphi(x)=\sum_{y \in f^{-1}(x)} \frac{\varphi(y)}{f^{\prime}(y)}
$$

has a spectral gap.
Let $g$ be $C^{2}$ except possibly at finite number of points and admitting a $C^{2}$ extension on each interval of smoothness. Define $X_{n}=g \circ f^{n}$ and consider it as a random variable with $x$ distributed according to some measure $\rho(x) d x, \rho \in$ $\operatorname{BV}[0,1]$.

Define a family of operators $\mathcal{L}_{t}: \mathrm{BV}[0,1] \rightarrow \mathrm{BV}[0,1]$ by

$$
\mathcal{L}_{t} \varphi(x)=\sum_{y \in f^{-1}(x)} \frac{e^{i t g(y)}}{f^{\prime}(y)} \varphi(y)
$$

where $t=0$ corresponds to the transfer operator. Because $g$ is bounded, writing $e^{i t g(y)}$ as a power series we can conclude $t \rightarrow \mathcal{L}_{t}$ is analytic for all $t$. This gives (A1).
(A2) follows from the fact that $\mathcal{L}_{0}$ has a spectral gap. We further assume that
$g$ is not cohomologous to a piece-wise constant function.

In particular, $g$ is not a BV coboundary.
The assumption (3.46) is reasonable. Indeed, suppose that $g$ is piece-wise constant taking values $c_{1}, c_{2} \ldots c_{k}$. Then $S_{n}$ takes less than $n^{k-1}$ distinct values so the maximal jump is of order at least $n^{-(k-1)}$ so $S_{n}$ can not admit Edgeworth expansion of order $(2 k-2)$ in contrast to the case where (3.46) holds as we shall see below.

A direct computation gives,

$$
\mathbb{E}\left(e^{i t S_{n} / \sqrt{n}}\right)=\int_{0}^{1} \mathcal{L}_{t / \sqrt{n}}^{n} \rho(x) d x
$$

Therefore, there exists $A$ such that,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \mathbb{E}\left(e^{i t \frac{S_{n}-n A}{\sqrt{n}}}\right)=e^{-t^{2} \sigma^{2} / 2} \tag{3.47}
\end{equation*}
$$

where $\sigma^{2} \geq 0$. It is well know that $\sigma^{2}>0 \Longleftrightarrow g$ is a BV coboundary (see [24]). From (3.47) it is clear that $S_{n}$ satisfies the CLT.

To show (A3) holds, we first normalize the family of operators,

$$
\overline{\mathcal{L}}_{t} v(x)=\sum_{f(y)=x} \frac{e^{i t g(y)} h(y)}{f^{\prime}(y) h \circ f(y)} v(y)
$$

Then, $\overline{\mathcal{L}}_{t}=H^{-1} \circ \mathcal{L}_{t} \circ H$ where $H$ is multiplication by the function $h$. Therefore, $\mathcal{L}_{t}$ and $\overline{\mathcal{L}}_{t}$ have the same spectrum. However, the eigenfunction corresponding to the eigenvalue 1 of $\overline{\mathcal{L}}_{0}$ changes to the constant function 1.

Assume $e^{i \theta}$ is an eigenvalue of $\overline{\mathcal{L}}_{t}$. Then, there exists $u \in \mathrm{BV}[0,1]$ with $\overline{\mathcal{L}}_{t} u(x)=e^{i \theta} u(x)$. Observe that,

$$
\begin{aligned}
\overline{\mathcal{L}}_{0}|u|(x) & =\sum_{f(y)=x} \frac{|u(y)| h(y)}{f^{\prime}(y) h \circ f(y)} \\
& \geq\left|\sum_{f(y)=x} \frac{e^{i t g(y)} u(y) h(y)}{f^{\prime}(y) h \circ f(y)}\right|=\left|\overline{\mathcal{L}}_{t} u(x)\right|=\left|e^{i \theta} u(x)\right|=|u(x)|
\end{aligned}
$$

Also note that, $\overline{\mathcal{L}}_{0}$ is a positive operator. Hence, $\overline{\mathcal{L}}_{0}^{n}|u|(x) \geq|u(x)|$ for all $n$. However,

$$
\lim _{n \rightarrow \infty}\left(\overline{\mathcal{L}}_{0}^{n}|u|\right)(x)=\int|u(y)| \cdot 1 d y
$$

because 1 is the eigenfunction corresponding to the top eigenvalue. So for all $x$,

$$
\int|u(y)| d y \geq|u(x)|
$$

This implies that $|u(x)|$ is constant. WLOG $|u(x)| \equiv 1$. So we can write $u(x)=$ $e^{i \gamma(x)}$. Then,

$$
\begin{gathered}
\overline{\mathcal{L}}_{t} u(x)=\sum_{f(y)=x} \frac{h(y)}{f^{\prime}(y) h \circ f(y)} e^{i(t g(y)+\gamma(y))}=e^{i(\theta+\gamma(x))} \\
\Longrightarrow \sum_{f(y)=x} \frac{h(y)}{f^{\prime}(y) h \circ f(y)} e^{i(t g(y)+\gamma(y)-\gamma(f(y))-\theta)}=1
\end{gathered}
$$

for all $x$. Since,

$$
\overline{\mathcal{L}}_{0} 1=\sum_{f(y)=x} \frac{h(y)}{f^{\prime}(y) h \circ f(y)}=1
$$

and $e^{i(t g(y)+\gamma(y)-\gamma(x)-\theta)}$ are unit vectors, it follows that

$$
\begin{equation*}
\operatorname{tg}(y)+\gamma(y)-\gamma(f(y))-\theta=0 \quad \bmod 2 \pi \tag{3.48}
\end{equation*}
$$

for all $y$. Because $g$ is not cohomologous to a piecewise constant function we have a contradiction. Therefore, $\overline{\mathcal{L}}_{t}$ and hence $\mathcal{L}_{t}$ does not have an eigenvalue on the unit circle when $t \neq 0$.

To complete the proof of (A3) one has to show that the spectral radius of $\mathcal{L}_{t}$ is at most 1 and that the essential spectral radius of $\mathcal{L}_{t}$ is strictly less than 1 . This is clear from Lasota-Yorke type inequality in [5, Lemma 1]. In fact, there is a uniform $\kappa \in(0,1)$ such that $r_{\text {ess }}\left(\mathcal{L}_{t}\right) \leq \kappa$ for all $t$.

Next, we describe in detail how the estimate in [5, Proposition 1] gives us (A4). To make the notation easier we assume $t>0$ and we replace $|t|$ by $t$. [5, Proposition 1] implies that there exist $c$ and $C$ such that if $K_{1}$ large enough (we fix one such $\left.K_{1}\right)$ then for all $t>K_{1}$,

$$
\begin{equation*}
\left\|\mathcal{L}_{t}^{\lceil c \ln t\rceil} u\right\|_{t} \leq e^{-C\lceil c \ln t\rceil}\|u\|_{t} \tag{3.49}
\end{equation*}
$$

where $\|h\|_{t}=(1+t)^{-1}\|h\|_{\mathrm{BV}}+\|h\|_{\mathrm{L}^{1}}$. Therefore,

$$
\left\|\mathcal{L}_{t}^{k[c \ln t\rceil} u\right\|_{t} \leq e^{-C\lceil c \ln t\rceil}\left\|\mathcal{L}^{(k-1)\lceil c \ln t\rceil} u\right\|_{t} \leq \cdots \leq e^{-C k[c \ln t\rceil}\|u\|_{t}
$$

Also, $\left\|\mathcal{L}_{t}\right\|_{t} \leq 1$. So, if $n=k\lceil c \ln t\rceil+r$ where $0 \leq r<\lceil c \ln t\rceil$ then

$$
\left\|\mathcal{L}_{t}^{n} u\right\|_{t} \leq e^{-C k[c \ln t]}\left\|\mathcal{L}_{t}^{r} u\right\|_{t} \leq e^{-C n \frac{k[c \ln t]}{k[c \ln t]+r}}\|u\|_{t} \leq e^{-C n \frac{k}{k+1}}\|u\|_{t}
$$

However,

$$
(1+t)^{-1}\|h\|_{\mathrm{BV}} \leq\|h\|_{t} \leq\left[1+(1+t)^{-1}\right]\|h\|_{\mathrm{BV}}
$$

Therefore,

$$
(1+t)^{-1}\left\|\mathcal{L}_{t}^{n} u\right\|_{\mathrm{BV}} \leq\left[1+(1+t)^{-1}\right] e^{-C n \frac{k}{k+1}}\|u\|_{\mathrm{BV}}
$$

which gives us

$$
\left\|\mathcal{L}_{t}^{n}\right\|_{\mathrm{BV}} \leq(t+2) e^{-C n \frac{k}{k+1}}
$$

and here $k=k(n, t)=\left\lfloor\frac{n}{\lceil c \ln t\rceil}\right\rfloor$. When $K_{1} \leq|t| \leq n^{r_{1}}, k_{\text {min }}=\left\lfloor\frac{n}{\left\lceil c \ln n^{\left.r_{1}\right\rceil}\right.}\right\rfloor$ and $\frac{k_{\min }}{k_{\min }+1} \rightarrow 1$ as $n \rightarrow \infty$. Also, $1 \geq \frac{k}{k+1} \geq \frac{k_{\min }}{k_{\min }+1}$ and,

$$
\left\|\mathcal{L}_{t}^{n}\right\|_{\mathrm{BV}} \leq(t+2) e^{-C n \frac{k(n, t)}{k(n, t)+1}} \leq 2 n^{r_{1}} e^{-C n \frac{k_{\min }}{k_{\min }+1}}
$$

Choosing $n_{0}$ such that for all $n>n_{0}, \frac{k_{\min }}{k_{\min }+1}>\frac{1}{2}$ (so this choice of $n_{0}$ works for all $t$ ) we can conclude that,

$$
\left\|\mathcal{L}_{t}^{n}\right\|_{\mathrm{BV}} \leq 2 n^{r_{1}} e^{-C n / 2}
$$

This proves (A4) for all choices of $r_{1}$. In particular given $r$, we can choose $r_{1}>\frac{r-1}{2}$ in the above proof. This implies that Edgeworth expansions of all orders exist.

### 3.5.5 Multidimensional expanding maps.

Let $\mathcal{M}$ be a compact Riemannian manifold and $f: \mathcal{M} \rightarrow \mathcal{M}$ be a $C^{2}$ expanding map. Let $g: \mathcal{M} \rightarrow \mathbb{R}$ be a $C^{2}$ function which is non homologous to constant. The proof of Lemma 3.13 in [13] shows that this condition is equivalent to $g$ not being infinitesimally integrable in the following sense. The natural extension of $f$ acts on the space of pairs $\left(\left\{y_{n}\right\}_{n \in \mathbb{N}}, x\right)$ where $f\left(y_{n+1}\right)=y_{n}$ for $n>0$ and $f y_{1}=x$. Given such pair let

$$
\Gamma\left(\left\{y_{n}\right\}, x\right)=\lim _{n \rightarrow \infty} \frac{\partial}{\partial x}\left[\sum_{k=0}^{n-1} g\left(f^{k} y_{n}\right)\right]=\lim _{n \rightarrow \infty} \frac{\partial}{\partial x}\left[\sum_{k=1}^{n} g\left(y_{k}\right)\right]=\sum_{k=1}^{\infty} \frac{\partial}{\partial x} g\left(y_{k}\right) .
$$

$g$ is called infinitesimally integrable if $\Gamma\left(\left\{y_{n}\right\}, x\right)$ actually depends only on $x$ but not on $\left\{y_{n}\right\}$.

Let $X_{n}=g \circ f^{n}$. We want to verify (A1)-(A4) when $x$ is distributed according to a smooth density $\rho$. Note that assumption (3.1) holds with $v=\rho, \ell$ being the

Lebesgue measure and

$$
\left(\mathcal{L}_{t} \phi\right)(x)=\sum_{y \in f^{-1}(x)} \frac{e^{i t g(y)}}{\left|\operatorname{det}\left(\frac{\partial f}{\partial x}\right)\right|} \phi(y) .
$$

We will check (A1)-(A4) for $\mathcal{L}_{t}$ acting on $C^{1}(\mathcal{M})$. The proof of (A1)-(A3) is the same as in section 3.5.4. In particular, for (A3) we need Lasota-Yorke inequality (see (3.52) below) which is proven in [13, equation (19)].

The proof of (A4) is also similar to section 3.5.4, so we just explain the differences. As before we assume that $t>0$. Given a small constant $\kappa$ let

$$
\|\phi\|_{t}=\max \left(\|\phi\|_{C^{0}}, \frac{\kappa\|D \phi\|_{C^{0}}}{1+t}\right)
$$

Then by [13, Proposition 3.16]

$$
\begin{equation*}
\left\|\mathcal{L}_{t}^{n} \phi\right\|_{t} \leq\|\phi\|_{t} \tag{3.50}
\end{equation*}
$$

provided that $n \geq C_{1} \ln t$.
By [13, Lemma 3.18] if $g$ is not infinitesimally integrable then there exists a constant $\eta<1$ such that

$$
\begin{equation*}
\left\|\mathcal{L}_{t}^{n} \phi\right\|_{L^{1}} \leq \eta^{n}\|\phi\|_{t} \tag{3.51}
\end{equation*}
$$

The Lasota-Yorke inequality says that there is a constant $\theta<1$, such that

$$
\begin{equation*}
\left\|D\left(\mathcal{L}_{t}^{n} \phi\right)\right\|_{C^{0}} \leq C_{3}\left(t\|\phi\|_{C^{0}}+\theta^{n}\|D \phi\|_{C^{0}}\right) \tag{3.52}
\end{equation*}
$$

Also,

$$
\begin{equation*}
\left\|\mathcal{L}_{t}^{n} \phi\right\|_{C^{0}} \leq\left\|\mathcal{L}_{0}^{n}(|\phi|)\right\|_{C^{0}} \leq C_{4}\left(\||\phi|\|_{L^{1}}+\theta^{n}\||\phi|\|_{\text {Lip }}\right) \tag{3.53}
\end{equation*}
$$

where the last step relies on $\mathcal{L}_{0}$ having a spectral gap on the space of Lipshitz functions. Combing (3.50) through (3.53), we conclude that $\mathcal{L}_{t}$ satisfies (3.49). The rest of the argument is the same as in section 3.5.4.

## Chapter 4: Large Deviation Principles.

### 4.1 Asymptotics for Cramér's Theorem.

In this section, we focus on sequences of i.i.d. random variables. First, we prove the existence of weak asymptic expansions for Cramér's LDP - Theorem 1.2. Next, we deduce existence of the strong expansion in special cases. As expected, a stronger assumption on the regularity of the law of the random variables is required for the second step.

### 4.1.1 Weak asymptotic expansions.

We recall that a random variable $X$ is called $l$-Diophantine if there exist positive constants $t_{0}$ and $C$ such that $\left|\mathbb{E}\left(e^{i t X}\right)\right|<1-\frac{C}{|t|^{l}}$ for $|t|>t_{0}$. It is known that when $X$ is $l$-Diophantine and $r+2$ moments exist weak Edgeworth expansions exist. For example, see [4] and Section 3.5.1.

Given a random variable $X$ with distribution function $F$, we define $Y_{X, \gamma}$ to be a random variable with distribution function $G^{\gamma}$ given by,

$$
\begin{equation*}
d G^{\gamma}(y)=\frac{e^{y \gamma} d F(y)}{\mu(\gamma)} \tag{4.1}
\end{equation*}
$$

where $\mu(\gamma)=\int e^{y \gamma} d F(y)$. Therefore,

$$
\begin{equation*}
\mathbb{E}\left[Y_{X, \gamma}\right]=\frac{\int y e^{y \gamma} d F(y)}{\int e^{y \gamma} d F(y)} \tag{4.2}
\end{equation*}
$$

In Section 3.1 we defined the function spaces $F_{k}^{m}: f \in F_{k}^{m}$ if $f$ is $m$ times continuously differentiable and $C_{k}^{m}(f)=\left(\max _{0 \leq j \leq m}\left\|f^{(j)}\right\|_{\mathrm{L}^{1}}+\max _{0 \leq j \leq k}\left\|x^{j} f\right\|_{\mathrm{L}^{1}}\right)<\infty$. We call a function $f$, (left) exponential of order $\alpha$, if $\lim _{x \rightarrow-\infty}\left|e^{-\alpha x} f(x)\right|=0$. Denote by $F_{m, \alpha}^{k}$ the collection of all $f \in F_{m}^{k}$ with $f^{(k)}$ is exponential of order $\alpha$.

We note that due to assumption $f \in F_{m}^{k}, f^{(k)}$ being exponential of order $\alpha$ is enough to guarantee that $f^{(l)}$ is exponential of order $\alpha$ for all $0 \leq l \leq k$. To see this suppose $f, f^{\prime} \in L^{1}$. Then, $\lim _{|x| \rightarrow \infty} f(x)=0$. Suppose $f^{\prime}$ is exponential of order $\alpha$. Then, given $\epsilon>0$ there is $M>0$ such that for $x<-M,-\epsilon e^{\alpha x}<f^{\prime}(x)<\epsilon e^{\alpha x}$. So, $-\epsilon \int_{-\infty}^{x} e^{\alpha y} d y \leq \int_{-\infty}^{x} f^{\prime}(y) d y \leq \epsilon \int_{-\infty}^{x} e^{\alpha y} d y \Longrightarrow-\frac{\epsilon}{\alpha} e^{\alpha x} \leq f(x) \leq \frac{\epsilon}{\alpha} e^{\alpha x}$ for $x<-M$. So $f$ is also of exponential order $\alpha$. Since $f^{(l)} \in L^{1}$ for all $0 \leq l \leq k$, we can repeat the same argument starting from $k$ and conclude that all lower order derivatives are of exponential order $\alpha$.

It is clear that $F_{m, \alpha}^{k} \subset F_{m, \beta}^{k}$ if $\alpha>\beta$. Finally, define, $F_{m, \infty}^{k}=\bigcap_{\alpha>0} F_{m, \alpha}^{k}$. This intersection is non-empty. For example, the family of Gaussian functions and $C_{c}^{\infty}(\mathbb{R})$ are in $F_{m, \alpha}^{k}$ for all $\alpha>0$.

Recall from Chapter 1 that for a function $f: \mathbb{R} \rightarrow(-\infty, \infty]$ with $f \neq \infty$, $D_{f}=\{x \in \mathbb{R} \mid f(x)<\infty\}$ and $f^{*}(x)=\sup _{t \in \mathbb{R}}[t x-f(t)]$. If $f$ is convex, lower semicontinuous with $\check{D}_{f}=(a, b)$ and $f \in C^{2}(a, b)$ with $f^{\prime \prime}>0$ on $(a, b)$ then, $\check{D}_{f^{*}}=$ $(A, B)$ where $A=\lim _{t \rightarrow a+} f^{\prime}(t)$ and $B=\lim _{t \rightarrow b-} f^{\prime}(t), f^{*}$ is continuously differentiable on $(A, B)$. For any $f$ satisfying the above properties, for any $x \in \check{D}_{f^{*}}$ the supremum in
the definition of $f^{*}(x)$ is achieved at a unique point. $f$ is called steep if $\lim _{t \rightarrow a+}\left|f^{\prime}(t)\right|=$ $\lim _{t \rightarrow b-}\left|f^{\prime}(t)\right|=\infty$.

Theorem 4.1.1. Let $X$ be a non-constant, real-valued, and centred random variable. Assume that the logarithmic moment generating function $h(\theta)=\log \mathbb{E}\left(e^{\theta X}\right)$ is finite on a neighbourhood of 0 . Further assume that there is $l \in \mathbb{N}$ such that for all $\theta \in \stackrel{\circ}{D}_{h}$, $Y_{X, \theta}$ is l-Diophantine. Let $X_{n}$ be a sequence of i.i.d. copies of $X$. Let $r \in \mathbb{N}$ and $a \in(0, \sup (\operatorname{supp} X))$. Let $\theta_{a}$ be the unique $\theta$ such that

$$
I(a)=\sup _{\theta \in \tilde{D}_{h}}\left(a \theta-\log \int e^{y \theta} d F(y)\right)=a \theta_{a}-\log \int e^{y \theta_{a}} d F(y)
$$

Take $q>\frac{l(r+2)}{2}+1$ and $\alpha>\theta_{a}$. Then, for every $f \in F_{r+1, \alpha}^{q}$ we have,

$$
\begin{equation*}
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+\frac{1}{2}}} \int P_{p}(z) f_{\theta_{a}}(z) d z+C_{r+1}^{q}\left(f_{\theta_{a}}\right) \cdot o_{r, \theta_{a}}\left(\frac{1}{N^{\frac{r+1}{2}}}\right) \tag{4.3}
\end{equation*}
$$

where $f_{\theta}(x)=e^{-\theta x} f(x)$ and $P_{p}(z)$ polynomials depending on a.

Proof. Assuming $F$ to be the distribution function of $X$ we can define $Y_{X, \gamma}$ by (4.1). Let $Y_{i}$ 's be i.i.d. copies of $Y_{X, \gamma}$ and take $\tilde{S}_{N}=Y_{1}+\cdots+Y_{N}$. A simple computation gives us,

$$
d G_{N}^{\gamma}(y)=\frac{e^{y \gamma} d F_{N}(y)}{\mu(\gamma)^{N}}
$$

where $F_{N}$ is the distribution function of $S_{N}$ and $G_{N}^{\gamma}$ is the distribution function of $\tilde{S}_{N}$. Now, we formally compute,

$$
\begin{aligned}
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{a \gamma N} & =\mathbb{E}\left(e^{a \gamma N} f\left(S_{N}-a N\right)\right) \\
& =\mathbb{E}\left(e^{\gamma S_{N}} f_{\gamma}\left(S_{N}-a N\right)\right) \\
& =\int e^{\gamma y} 2 \pi f_{\gamma}(y-a N) d F_{N}(y)
\end{aligned}
$$

$$
\begin{aligned}
& =\mu(\gamma)^{N} \int 2 \pi f_{\gamma}(y-a N) d G_{N}^{\gamma}(y) \\
& =\mu(\gamma)^{N} \mathbb{E}_{\gamma}\left(2 \pi f_{\gamma}\left(\tilde{S}_{N}-a N\right)\right)
\end{aligned}
$$

where $f_{\gamma}(s)=\frac{1}{2 \pi} e^{-s \gamma} f(s)$. Hence,

$$
\begin{equation*}
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{(a \gamma-\log \mu(\gamma)) N}=\mathbb{E}_{\gamma}\left(2 \pi f_{\gamma}\left(\tilde{S}_{N}-a N\right)\right) \tag{4.4}
\end{equation*}
$$

Put $\gamma=\theta_{a}$. Then, $Y_{X, \gamma}$ has mean $a$ (see [17, Chapter 2]).
Since $f \in F_{r+1, \alpha}^{q}$ with $\theta_{a}<\alpha$ we have $f_{\theta_{a}} \in F_{r+1}^{q}$. We prove this when $r=0$ and $q=1$. The argument for general $q$ and $r$ is similar. Suppose, $f(x), f^{\prime}(x), x f(x) \in$ $L^{1}$ and $f^{\prime}(x)$ is continuous. It is immediate that $\left(e^{-\theta_{a} x} f(x)\right)^{\prime}=-\theta_{a} e^{-\theta_{a} x} f(x)+$ $e^{-\theta_{a} x} f^{\prime}(x)$ is continuous. We need to show, $e^{-\theta_{a} x} f(x),\left(e^{-\theta_{a} x} f(x)\right)^{\prime}, x e^{-\theta_{a} x} f(x) \in L^{1}$. Since $f$ and $f^{\prime}$ are of exponential order, it is enough to show, $e^{-\theta_{a} x} g(x), x e^{-\theta_{a} x} g(x) \in$ $L^{1}$ if $g$ is exponential of order $\alpha\left(>\theta_{a}\right)$. This is true because there is $M>0$ such that for $x<-M,\left|e^{-\theta_{a} x} f(x)\right|<e^{\left(\alpha-\theta_{a}\right) x}$ and $\left|x e^{-\theta_{a} x} f(x)\right|<-x e^{\left(\alpha-\theta_{a}\right) x}$.

Therefore, from [4], RHS of (4.4) admits the weak Edgeworth expansion whose coefficients are determined by moments of $Y_{X, \theta_{a}}$. Therefore, we have that for all functions $f \in F_{r+1, \alpha}^{q}$

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+\frac{1}{2}}} \int P_{p, l}(z) f_{\theta_{a}}(z) d z+C_{r+1}^{q}\left(f_{\theta_{a}}\right) \cdot o\left(\frac{1}{N^{\frac{r+1}{2}}}\right) .
$$

## Remark 4.1.1.

1. The assumption of $X$ being centred is just to simplify the notation. One can easily reformulate the results for non-centred $X$ using the corresponding results
for $X-\mathbb{E}(X)$. Therefore, from now on we discuss results for centred random variables only.
2. A similar result holds for $a \in(\inf (\operatorname{supp} X), 0)$. In fact, one can deduce the corresponding results for $a<0$ by considering $-X$ and $(-a)>0$. But, for simplicity we focus only on $a>0$ hereafter.
3. Note that the requirement to expand $\mathbb{E}_{\gamma}\left(f_{\theta_{a}}\left(\tilde{S}_{N}-a N\right)\right)$ is $f_{\theta_{a}} \in F_{r+1}^{q}$ which is indeed the case when $f \in F_{\theta_{a}, \alpha}^{q}$ for some $\alpha>\theta_{a}$. In particular, this result holds for $f \in C_{c}^{q}(\mathbb{R})$.
4. In addition, if $h(\theta)$ is steep then $\sup (\operatorname{supp} X)=\infty$ (see [30, Chapter 1]) and the expansion holds for all $a>0$.

We note that for a large class of random variables $X, Y_{X, \theta}$ is $l$-Diophantine. For example, if $X$ is $0-$ Diophantine then so is $Y_{X, \theta}$ because $X$ is absolutely continuous with respect to $Y_{X, \theta}$ (see [1, Lemma 4]). Also, we claim that if $X$ is compactly supported and $l$-Diophantine for $l>0$ then so is $Y_{X, \theta}$.

We recall from [4], that a random variable $X$ with distribution function $F$ is $l-$ Diophantine if and only if there exists $C_{1}, C_{2}>0$ such that for all $|x|>C_{1}$,

$$
\inf _{y \in \mathbb{R}} \int_{\mathbb{R}}\{a x+y\}^{2} d F(a) \geq \frac{C_{2}}{|x|^{l}}
$$

where $\{z\}=\operatorname{dist}(z, \mathbb{Z})$. If $X$ is compactly supported (say on $[c, d])$ then,

$$
\begin{aligned}
\int_{\mathbb{R}}\{a x+y\}^{2} d G^{\theta}(a) & =\frac{1}{\int_{c}^{d} e^{\theta a} d F(a)} \int_{c}^{d}\{a x+y\}^{2} e^{\theta a} d F(a) \\
& \geq \frac{e^{\theta c}}{\int_{\mathbb{R}} e^{\theta a} d F(a)} \int_{c}^{d}\{a x+y\}^{2} d F(a)
\end{aligned}
$$

Thus, for all $|x|>C_{1}$,

$$
\inf _{y \in \mathbb{R}} \int_{\mathbb{R}}\{a x+y\}^{2} d G^{\theta}(a) \geq \frac{e^{\theta c}}{\int_{c}^{d} e^{\theta a} d F(a)} \frac{C_{2}}{|x|^{l}}
$$

So the random variable $Y_{X, \theta}$ with distribution function $G^{\theta}$ is $l$-Diophantine as claimed earlier. From this we obtain the following corollary.

Corollary 4.1.2. Let $X$ be a non-constant, real-valued, compactly supported and l-Diophantine centred random variable. Let $X_{n}$ be a sequence of i.i.d. copies of $X$. Let $r \in \mathbb{N}$ and $a \in(0, \sup (\operatorname{supp} X))$. Let $\theta_{a}$ be the unique $\theta$ such that

$$
I(a)=\sup _{\theta \in \dot{D}_{h}}\left(a \theta-\log \int e^{y \theta} d F(y)\right)=a \theta_{a}-\log \int e^{y \theta_{a}} d F(y)
$$

Then, for every $f \in F_{r+1, \alpha}^{q}$ with $q>\frac{l(r+2)}{2}+1$ and $\alpha>\theta_{a}$ we have,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+\frac{1}{2}}} \int P_{p}(z) f_{\theta_{a}}(z) d z+C_{r+1}^{q}\left(f_{\theta_{a}}\right) \cdot o_{r, \theta_{a}}\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

for some polynomials $P_{p}(z)$ depending on $a$.

### 4.1.2 Strong asymptotic expansions.

We prove a lemma that gives conditions for the point-wise limit of a sequence of functions uniformly bounded in $F_{r+1}^{q}$ to satisfy the asymptotic expansions.

Lemma 4.1.3. Let $q \geq 0$. Suppose $\left\{f_{k}\right\}$ is a sequence in $F_{r+1}^{q}, S_{N}$ admits the weak local Edgeworth expansion for $f_{k}, C_{r+1}^{q}\left(f_{k}\right) \leq C$ for all $k$, $f_{k}$ are uniformly bounded in $L^{\infty}(\mathbb{R}), f_{k} \rightarrow f$ point-wise and for all $p$,

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int P_{p}(z) f_{k}(z) d z=\int P_{p}(z) f(z) d z \tag{4.5}
\end{equation*}
$$

Then,

$$
\sqrt{N} \mathbb{E}\left(f\left(S_{N}\right)\right)=\frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p}} \int P_{p}(z) f(z) d z+C \cdot o_{r, \beta}\left(N^{-r / 2}\right) .
$$

Proof. For large $N$,

$$
\begin{align*}
\left|\sqrt{N} \mathbb{E}\left(f_{k}\left(S_{N}\right)\right)-\frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p}} \int P_{p}(z) f_{k}(z) d z\right| & \leq C_{r+1}^{q}\left(f_{k}\right) \cdot o_{r, \beta}\left(N^{-r / 2}\right) \\
& \leq C \cdot o_{r, \beta}\left(N^{-r / 2}\right) \tag{4.6}
\end{align*}
$$

LDCT gives us that,

$$
\lim _{k \rightarrow \infty} \mathbb{E}\left(f_{k}\left(S_{N}\right)\right)=\mathbb{E}\left(f\left(S_{N}\right)\right)
$$

This along with assumption (4.5) allows us to take the limit $k \rightarrow \infty$ in the RHS of (4.6) and to conclude,

$$
\left|\sqrt{N} \mathbb{E}\left(f\left(S_{N}\right)\right)-\frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p}} \int P_{p}(z) f(z) d z\right| \leq C \cdot o_{r, \beta}\left(N^{-r / 2}\right)
$$

which implies the result.

Remark 4.1.2. The same would hold if we replace weak local by weak global. However, our focus here is on weak local expansions.

The next theorem specifies when the existence of weak expansions imply the existence of strong expansions.

Theorem 4.1.4. Let $X_{n}$ be a sequence of random variables not necessarily i.i.d. Suppose $S_{N}=X_{1}+\cdots+X_{N}$ admits the weak asymptotic expansion of order $r$ for large deviations in the range $(0, L)$ for $f \in F_{r+1, L_{+}}^{1}$ where $L_{+}>L$ when $L<\infty$ and $L_{+}=\infty$ if $L=\infty$. That is,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+1 / 2}} \int P_{p}(z) f_{\theta_{a}}(z) d z+C_{r+1}^{1}\left(f_{\theta_{a}}\right) \cdot o_{r, \theta_{a}}\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

for all $a \in(0, L)$ where $I(a)$ and $\theta_{a}$ as in (4.11). Then, $S_{N}$ admits the strong asymptotic expansion of order $r$ for large deviation in $(0, L)$.

Proof. If $f \in C_{c}^{\infty}$ then $f_{\theta} \in F_{r+1}^{1}$ for all $\theta$. Therefore, we approximate $1_{[0, \infty)}$ by a sequence $f_{k}$ of $C_{c}^{\infty}$ functions such that $\left(f_{k}\right)_{\theta_{a}}$ are uniformly bounded in $F_{r+1}^{1}$ (see Appendix A. 3 for such a sequence) and invoke Lemma 4.1.3 to establish,

$$
\mathbb{P}\left(S_{N} \geq a N\right) e^{I(a) N}=\frac{1}{2 \pi} \sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+1 / 2}} \int_{0}^{\infty} P_{p}(z) e^{-\theta_{a} z} d z+C \cdot o_{r, \theta_{a}}\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

Remark 4.1.3. Note that the coefficients of the strong expansion are $C_{p}(a)=$ $\frac{1}{2 \pi} \int_{0}^{\infty} P_{p}(z) e^{-\theta_{a} z} d z$ obtained by replacing $f$ with $1_{[0, \infty)}$ in coefficients of the weak expansions. Since $f_{k}$ 's are bounded in $F_{r+1}^{1}$, we can do this without altering the order of the error. However, for any $q>1,1_{[0, \infty)}$ is not a pointwise limit of a sequence of functions $f_{k}$ in $F_{r}^{q}$ with $C_{r+1}^{q}\left(f_{k}\right)$ bounded. To see this, assume that $\left\|f_{k}\right\|_{1},\left\|f_{k}^{\prime}\right\|_{1},\left\|f_{k}^{\prime \prime}\right\|_{1}$ are uniformly bounded and $f_{k} \rightarrow 1_{[0, \infty)}$ point-wise. Then, for all $\phi \in C_{c}^{\infty}(\mathbb{R})$,

$$
\int \delta^{\prime} \phi=-\int \delta \phi^{\prime}=\int 1_{[0, \infty)} \phi^{\prime \prime}=\lim _{k \rightarrow \infty} \int f_{k} \phi^{\prime \prime}=\lim _{k \rightarrow \infty}-\int f_{k}^{\prime} \phi^{\prime}=\lim _{k \rightarrow \infty} \int f_{k}^{\prime \prime} \phi
$$

This implies that $\frac{\left|\phi^{\prime}(0)\right|}{\|\phi\|_{\infty}} \leq \sup _{k}\left\|f_{k}^{\prime \prime}\right\|_{1}$ for all $\phi \in C_{c}^{\infty}(\mathbb{R})$. Clearly, this is a contradiction. Therefore, Theorem 4.1.1 does not automatically give us strong expansions.

Now we are in a position to state and prove the main result of this section, which extends Cramér's LDP for i.i.d. random variables when the random variables have a sufficiently regular density.

Theorem 4.1.5. Let $X$ be a non-constant real valued centred random variable. Assume that the logarithmic moment generating function $h(\theta)=\log \mathbb{E}\left(e^{\theta X}\right)$ is finite on a neighbourhood of 0 . Further assume that, $X$ is 0 -Diophantine. Let $r \in \mathbb{N}$. Then for all $a \in(0, \sup (\operatorname{supp} X))$, there are constants $C_{p}(a)$ such that

$$
\mathbb{P}\left(S_{N} \geq a N\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{C_{p}(a)}{N^{p+\frac{1}{2}}}+o\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

where

$$
C_{p}(a)=\frac{1}{2 \pi} \int_{0}^{\infty} e^{-\theta_{a} z} P_{p}(z) d z
$$

for some polynomials $P_{p}(z)$ depending on $a$,

$$
I(a)=\sup _{\theta \in \mathbb{R}}\left(a \theta-\log \int e^{y \theta} d F(y)\right)
$$

and $\theta_{a}$ is this unique point the supremum is achieved.

Proof. If $X$ is 0 -Diophantine then so is $Y_{X, \theta}$ as $X$ is absolutely continuous with respect to $Y_{X, \theta}$ (see [1, Lemma 4]). Since, $Y_{X, \theta}$ has moments of all orders, $Y_{X, \theta}$ admits the strong Edgeworth expansion of all orders. Therefore, for each $r \in \mathbb{N}, Y_{X, \theta}$ admits the weak local Edgeworth expansion of order $r$ for $f \in F_{r}^{1}$ (see Appendix A.2).

From (4.4) we know that,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\mathbb{E}_{\gamma}\left(2 \pi f_{\theta_{a}}\left(\tilde{S}_{N}-a N\right)\right)
$$

where summands of $\tilde{S}_{N}$ have mean $a$. The assumptions allow us to expand RHS using the weak local Edgeworth expansion and obtain,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+\frac{1}{2}}} \int P_{p}(z) f_{\theta_{a}}(z) d z+C_{r+1}^{1}\left(f_{\theta_{a}}\right) \cdot o_{r, \beta}\left(N^{-r / 2}\right)
$$

for $f \in C_{c}^{\infty}(\mathbb{R})$.
Now, we approximate $1_{[0, \infty)}$, by a sequence $f_{k} \in C_{c}^{\infty}(\mathbb{R})$ such that $\left(f_{k}\right)_{\theta_{a}}$ are bounded in $F_{r+1}^{1}$ (see Appendix A. 3 for such a sequence) and use Theorem 4.1.4 to obtain the required expansion.

Remark 4.1.4. This gives us an alternative proof of [1, Theorem 2] for $X$ satisfying the Cramér's condition (which corresponds to Case 1 there).

There are two ways the coefficients $C_{p}(a)$ depend on $a$. First note that $\theta_{a}$ depends on the choice of $a$. Also, from Section 3.3, we know exactly how the coefficients of $P_{p}$ depend on the first $p+2$ asymptotic moments of $\tilde{S}_{N}$ and thus, on the first $p+2$ moments of $Y_{X, \theta_{a}}$. So the dependence of $C(a)$ on $a$ is explicit and one can compute these coefficients. In addition, $C_{p}(a)$ does not depend on $r$ because $P_{p}(z)$ 's do not.

### 4.2 Higher order asymptotics in the non-i.i.d. case.

Let $X_{n}$ be a sequence of random variables that are not necessarily i.i.d. with asymptotic mean 0 . Suppose that there exist a Banach space $\mathbb{B}$, a family of bounded linear operators $\mathcal{L}_{z}: \mathbb{B} \rightarrow \mathbb{B}$ and vectors $v \in \mathbb{B}, \ell \in \mathbb{B}^{\prime}$ such that

$$
\begin{equation*}
\mathbb{E}\left(e^{z S_{N}}\right)=\ell\left(\mathcal{L}_{z}^{N} v\right), z \in \mathbb{C} \tag{4.7}
\end{equation*}
$$

and satisfying the following,
(B1) There exists $\delta>0$ such that $z \mapsto \mathcal{L}_{z}$ is continuous on the strip $|\operatorname{Re}(z)|<\delta$ and holomorphic on the disc $|z|<\delta$.
(B2) 1 is an isolated and simple eigenvalue of $\mathcal{L}_{0}$, all other eigenvalues of $\mathcal{L}_{0}$ have absolute value less than 1 and its essential spectrum is contained strictly inside the disk of radius 1 (spectral gap).
(B1) and (B2) along with perturbation theory of operators (see [33]) imply that there is $\delta_{0} \in(0, \delta)$ such that

$$
\begin{equation*}
\mathcal{L}_{z}=\mu(z) \Pi_{z}+\Lambda_{z},|z|<\delta_{0} \tag{4.8}
\end{equation*}
$$

where $\mu(z)$ is the top eigenvalue of $\mathcal{L}_{z}, \Pi_{z}$ is the corresponding eigen-projection, $\Pi_{z} \Lambda_{z}=\Lambda_{z} \Pi_{z}=0$ and $z \mapsto \mu(z), z \mapsto \Pi_{z}$ and $z \mapsto \Lambda_{z}$ are holomorphic. In addition, $\left\|\frac{d^{k}}{d z^{k}} \Lambda_{z}^{N}\right\|<\beta_{k}^{N}$ with $0<\beta_{k}<1$. Therefore,

$$
\mathcal{L}_{z}^{N}=\mu(z)^{N} \Pi_{z}+\Lambda_{z}^{N}
$$

Combining this with (4.7) we have,

$$
\begin{equation*}
\mathbb{E}\left(e^{z S_{N}}\right)=\mu(z)^{N} \ell\left(\Pi_{z} v\right)+\ell\left(\Lambda_{z} v\right) \tag{4.9}
\end{equation*}
$$

Then, plugging in $z=0$ and taking $N \rightarrow \infty$, we conclude that $\ell\left(\Pi_{0} v\right)=1$. Also, taking the derivative at $z=0$, dividing by $N$ and taking the limit as $N \rightarrow \infty$, we obtain,

$$
\left.\frac{d}{d z} \mu(z)\right|_{z=0}=\lim _{N \rightarrow \infty} \frac{\mathbb{E}\left(S_{N}\right)}{N}=0
$$

Taking the second derivative at $z=0$, dividing by $N^{2}$ and taking the limit as $N \rightarrow \infty$, we obtain,

$$
\left.\frac{d^{2}}{d z^{2}} \mu(z)\right|_{z=0}=\lim _{N \rightarrow \infty} \frac{\mathbb{E}\left(S_{N}^{2}\right)}{N^{2}}
$$

In addition, it follows from [24][Theorem 2.4] that there exists $\sigma^{2} \geq 0$ such
that $\frac{S_{N}}{\sqrt{N}} \xrightarrow{d} \mathcal{N}\left(0, \sigma^{2}\right)$. Since our interest is in $S_{N}$ that satisfies the CLT we would asumme from now on that $\sigma^{2}>0$.

We also assume the following:
(B3) $\mu(\theta)>0$ for all $\theta \in\left(-\delta_{0}, \delta_{0}\right)$ (Here $\delta_{0}$ as in (4.8)).
Define $\Omega(\theta)=\log \mu(\theta)$ for $\theta \in\left(-\delta_{0}, \delta_{0}\right)$. Then, $\Omega(0)=\log \mu(0)=0$ and $\Omega^{\prime}(0)=\frac{\mu^{\prime}(0)}{\mu(0)}=0$. Also, $\Omega^{\prime \prime}(0)=\frac{\mu^{\prime \prime}(0) \mu(0)-\mu^{\prime}(0)^{2}}{\mu(0)^{2}}=\mu^{\prime \prime}(0)=\sigma^{2}>0$. Since $\Omega^{\prime \prime}$ is continuous, there exists $\delta_{1} \in\left(0, \delta_{0}\right)$ such that $\Omega$ is strictly convex on $\left(-\delta_{1}, \delta_{1}\right)$. Note that due to convexity, $\Omega^{\prime}\left(-\delta_{1}\right)<0<\Omega^{\prime}\left(\delta_{1}\right)$. In addition, when $\theta \neq 0$, $\mu(\theta)>\mu(0)=1$ by convexity.

Next, we consider the Legendre transform of $\Omega, I$ given by,

$$
I(a)=\sup _{\theta \in\left(-\delta_{1}, \delta_{1}\right)}[a \theta-\Omega(\theta)], \text { for } a \in\left[0, \Omega^{\prime}\left(\delta_{1}\right)\right)
$$

which itself is a strictly convex function.
Because $\Omega^{\prime}$ is strictly increasing and continuous on [0, $\left.\Omega^{\prime}\left(\delta_{1}\right)\right]$, $a-\Omega^{\prime}(\theta)=0$ has a unique solution $\theta_{a}$ which depends continuously on $a$. Note that $I(a) \geq 0$ for all $a$ and $I(a)=0 \Longleftrightarrow a=0$. Also, $I(a)$ is continuous because $I$ is convex and $I(0)=0$. In addition, $I\left(\Omega^{\prime}\left(\delta_{1}\right)\right)=a \delta_{1}-\Omega\left(\delta_{1}\right)$.

Now, we are in a position to prove a Large Deviation Principle for $S_{N}$ using Theorem 1.3. The following lemma shows that Theorem 1.3 applies in our case.

Lemma 4.2.1. Suppose (B1), (B2) and (B3) hold. Then, there exists $0<\delta_{2} \leq \delta_{1}$ such that for $\theta \in\left(-\delta_{2}, \delta_{2}\right)$,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{E}\left(e^{\theta S_{N}}\right)=\log \mu(\theta)
$$

Proof. Because $\ell\left(\Pi_{0} v\right)>0$, there exists $\delta_{2}$ and $m>0$ such that for $\theta \in\left[-\delta_{2}, \delta_{2}\right]$
$\ell\left(\Pi_{\theta} v\right)>2 m$. Because $\left\|\Lambda_{\theta}^{N}\right\|<C \mu(\theta)^{N}$ for large $N$, we have that

$$
\lim _{N \rightarrow \infty} \mu(\theta)^{-N} \ell\left(\Lambda_{\theta}^{N} v\right)=0
$$

Hence, there exists $N_{0}$ such that for $N>N_{0}$,

$$
m<\ell\left(\Pi_{\theta} v\right)+\mu(\theta)^{-N} \ell\left(\Lambda_{\theta}^{N} v\right)<3 m
$$

Hence,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \ln \left[\ell\left(\Pi_{\theta} v\right)+\mu(\theta)^{-N} \ell\left(\Lambda_{\theta}^{N} v\right)\right]=0
$$

Now, for $\theta \in\left(-\delta_{2}, \delta_{2}\right)$ we can rewrite (4.9) as

$$
\frac{1}{N} \log \mathbb{E}\left(e^{\theta S_{N}}\right)=\log \mu(\theta)+\frac{1}{N} \log \left[\ell\left(\Pi_{\theta} v\right)+\mu(\theta)^{-N} \ell\left(\Lambda_{\theta}^{N} v\right)\right]
$$

This implies that,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{E}\left(e^{\theta S_{N}}\right)=\log \mu(\theta)
$$

Combining this lemma with Theorem 1.3 and the analysis proceeding it, we have the following LDP.

Theorem 4.2.2. Suppose $(B 1),(B 2)$ and (B3) hold. Then, there exists $\delta_{2} \in\left(0, \delta_{1}\right]$ such that for all $a \in\left(0, \frac{\log \mu\left(\delta_{2}\right)}{\delta_{2}}\right)$,

$$
\begin{equation*}
\lim _{N \rightarrow \infty} \frac{1}{N} \log \mathbb{P}\left(S_{N} \geq a N\right)=-I(a) \tag{4.10}
\end{equation*}
$$

where

$$
\begin{equation*}
I(a)=\sup _{\theta \in\left(-\delta_{2}, \delta_{2}\right)}[a \theta-\log \mu(\theta)]=a \theta_{a}-\log \mu\left(\theta_{a}\right) \tag{4.11}
\end{equation*}
$$

and $\theta_{a}$ is the unique $\theta$ solving $(\log \mu(\theta))^{\prime}=\frac{\mu^{\prime}(\theta)}{\mu(\theta)}=a$.

Remark 4.2.1. The range of a for which the LDP holds, is constrained by the assumptions $(B 1),(B 2)$ and (B3). We require a positive top eigenvalue $\mu(\theta)$ to exist, $\log \mu(\theta)$ to be strictly convex and $\ell\left(\Pi_{\theta} v\right)>0$. Larger the range of $\theta$ for which these hold, larger the range of $a$. In particular, if these hold for all $\theta \in \mathbb{R}$, then by convexity $B=\lim _{\delta \rightarrow \infty} \frac{\log \mu(\delta)}{\delta}$ exists as an extended real number and for all $a \in(0, B)$ the LDP holds.

Next, we compute higher order asymptotics of this LDP. To this end, we make two more assumptions about $\mathcal{L}_{z}$.
(B4) For all $\theta \in\left(-\delta_{2}, \delta_{2}\right)$, for all real numbers $t \neq 0, \operatorname{sp}\left(\mathcal{L}_{\theta+i t}\right) \subset\{|z|<\mu(\theta)\}$.
(B5) There are positive real numbers $r_{1}, r_{2}, C, K$ and $N_{0}$ such that for all $\theta \in$ $\left(-\delta_{2}, \delta_{2}\right)$, for all $N>N_{0}$ and for all $K<|t|<N^{r_{1}},\left\|\mathcal{L}_{\theta+i t}^{N}\right\| \leq C \frac{\mu(\theta)^{N}}{N^{r_{2}}}$.

Remark 4.2.2. As in Remark 3.1.1 it follows that by slightly decreasing $r_{1}$ we can assume $r_{2}$ to be as large as required for large enough $N$.

Pick $a \in\left(0, \frac{\log \mu\left(\delta_{2}\right)}{\delta_{2}}\right)$. Then,

$$
\begin{aligned}
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{a \theta N} & =\mathbb{E}\left(e^{\theta S_{N}} e^{-\left(S_{N}-a N\right) \theta} f\left(S_{N}-a N\right)\right) \\
& =\frac{1}{2 \pi} \int \widehat{f}_{\theta}(t) e^{-i a t N} \ell\left(\mathcal{L}_{\theta+i t}^{N} v\right) d t
\end{aligned}
$$

where $f_{\theta}(x)=\frac{1}{2 \pi} e^{-\theta x} f(x)$. Now define, $\overline{\mathcal{L}}_{\theta+i t}=\frac{e^{-i a t}}{\mu(\theta)} \mathcal{L}_{\theta+i t}$. Then,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{a \theta N}=\mu(\theta)^{N} \int \widehat{f}_{\theta}(t) \ell\left(\overline{\mathcal{L}}_{\theta+i t}^{N} v\right) d t
$$

From this we have,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{[a \theta-\log \mu(\theta)] N}=\int \widehat{f}_{\theta}(t) \ell\left(\overline{\mathcal{L}}_{\theta+i t}^{N} v\right) d t
$$

In particular,

$$
\begin{equation*}
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\int \widehat{f}_{\theta_{a}}(t) \ell\left(\overline{\mathcal{L}}_{\theta+i t_{a}}^{N} v\right) d t \tag{4.12}
\end{equation*}
$$

Note that for $\left|\theta_{a}+i t\right|<\delta_{0}$ the top eigenvalue of $\bar{L}_{\theta_{a}+i t}$ is $\bar{\mu}\left(\theta_{a}+i t\right)=\frac{e^{-i a t}}{\mu\left(\theta_{a}\right)} \mu\left(\theta_{a}+\right.$ $i t)$. As a function of $t, \bar{\mu}\left(\theta_{a}+i t\right)$ is analytic in a neighbourhood of 0 by (4.8). Further,

$$
\bar{\mu}\left(\theta_{a}\right)=1, \bar{\mu}^{\prime}\left(\theta_{a}\right)=\left.\frac{d}{d t} \bar{\mu}(z)\right|_{t=0}=-i a+i \frac{\mu^{\prime}\left(\theta_{a}\right)}{\mu\left(\theta_{a}\right)}=0, \bar{\mu}^{\prime \prime}\left(\theta_{a}\right)=-\frac{\mu^{\prime \prime}\left(\theta_{a}\right)}{\mu\left(\theta_{a}\right)}=-\sigma_{a}^{2}
$$

with $\sigma_{a}>0$. Thus, there exists $\bar{\delta}$ such that

$$
\begin{equation*}
\left|\bar{\mu}\left(\theta_{a}+i t\right)\right|<e^{-\sigma_{a}^{2} t^{2} / 4},|t|<\bar{\delta} \tag{4.13}
\end{equation*}
$$

We also notice that,

$$
\lim _{N \rightarrow \infty} \frac{\ell\left(\Lambda_{\theta}^{N} v\right)}{\mu(\theta)^{N}}=0
$$

because the spectral radius of $\Lambda_{\theta}$ is strictly smaller than $\mu(\theta)$. Combining this with $\mathbb{E}\left(e^{\theta S_{N}}\right)=\mu(\theta)^{N} \ell\left(\Pi_{\theta} v\right)+\ell\left(\Lambda_{\theta}^{N} v\right)$ we conclude that for all $\theta$,

$$
\ell\left(\Pi_{\theta} v\right)=\lim _{N \rightarrow \infty} \frac{\mathbb{E}\left(e^{\theta S_{N}}\right)}{\mu(\theta)^{N}}
$$

The following lemma allows us to obtain asymtotics of (4.12). We note that it is analogous to Theorem 3.1.4 where asymptotics of $\mathbb{E}\left(f\left(S_{N}-a N\right)\right)$ for $f \in F_{r+1}^{q+2}$ are discussed and can be proven using the ideas in the proof of Theorem 3.1.4. One just has to replace $\mathcal{L}_{t}$ by $\mathcal{L}_{\theta_{a}+i t}$ there and introduce the corresponding changes.

Lemma 4.2.3. Suppose (B1) through (B5) hold. Let $r \in \mathbb{N}$. Then, there exist $\delta_{2} \in(0, \delta)$ such that for all $a \in\left(0, \frac{\log \mu\left(\delta_{2}\right)}{\delta_{2}}\right)$ there are polynomials $P_{p}(z)$ such that for $g \in F_{r+1}^{q+1}$ with $q>\frac{r+1}{2 r_{1}}$,

$$
\int \widehat{g}(t) \ell\left(\overline{\mathcal{L}}_{\theta_{a}+i t}^{N} v\right) d t=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+1 / 2}} \int P_{p}(z) g(z) d z+C_{r+1}^{q+2}(g) \cdot o_{r, \theta_{a}}\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

where $\theta_{a}$ is as in (4.11).

Proof. We state how to estimate LHS away from 0 . The rest of the proof, which contains the construction of polynomials $P_{p}$, is identical to that of Theorem 3.1.5 with it replaced by $\theta_{a}+i t$.

Fix $\delta>0$ as in (4.13). By (B4), for $\delta \leq|t| \leq K$, there exists $c_{0} \in(0,1)$ such that $\left\|\overline{\mathcal{L}}_{\theta_{a}+i t}^{n}\right\| \leq c_{0}^{n}$. Thus, $\left|\int_{\delta<|t|<K} \widehat{g}(t) \ell\left(\overline{\mathcal{L}}_{\theta_{a}+i t}^{n} v\right) d t\right| \leq C\|g\|_{1} c_{0}^{n}$.

WLOG assuming $r_{2}>r_{1}+(r+1) / 2$,

$$
\begin{aligned}
\left|\int_{K<|t|<n^{r_{1}}} \widehat{g}(t) \ell\left(\overline{\mathcal{L}}_{\theta_{a}+i t}^{n} v\right) d t\right| \leq C\|g\|_{1} \int_{K<|t|<n^{r_{1}}}\left\|\mathcal{L}_{\theta_{a}+i t}^{n}\right\| d t & \leq \frac{C\|g\|_{1}}{n^{r_{2}-r_{1}}} \\
& =\|g\|_{1} o\left(n^{-(r+1) / 2}\right)
\end{aligned}
$$

Since, $g \in F_{r+1}^{q+2}$, we have that $t^{q} \widehat{g}(t)=(-i)^{q} \widehat{g^{(q)}}(t)$ and $\widehat{g^{(q)}}$ is integrable. Integrability of $\widehat{g^{(q)}}$ along with $q>\frac{r+1}{2 r_{1}}$ implies,

$$
\begin{align*}
\left|\int_{|t|>n^{r_{1}}} \widehat{g}(t) \ell\left(\overline{\mathcal{L}}_{\theta_{a}+i t}^{n} v\right) d t\right| \leq \int_{|t|>n^{r_{1}}}|\widehat{g}(t)| d t & \leq \int_{\left||t|>n^{r_{1}}\right.}\left|\frac{\widehat{g^{(q)}}(t)}{t^{q}}\right| d t  \tag{4.14}\\
& \leq \frac{\left\|g^{(q)}\right\|_{1}}{n_{1}^{r_{1} q}}=\left\|\widehat{g^{(q)}}\right\|_{1} o\left(n^{-(r+1) / 2}\right)
\end{align*}
$$

Therefore,

$$
\begin{equation*}
\left|\int_{|t|>\delta} \widehat{g}(t) \ell\left(\overline{\mathcal{L}}_{\theta_{a}+i t}^{n} v\right) d t\right|=o\left(n^{-(r+1) / 2}\right) \tag{4.15}
\end{equation*}
$$

## Remark 4.2.3.

1. The proof is almost identical to the proof of Theorem 3.1.4 and hence, the coefficients of polynomials $P_{p}$ can be computed as shown in Section 3.3. In particular, they depend on exponential moments of $S_{N}$.
2. Since $\theta_{a}$ depends on $a$, the coeffients of the polynomials $P_{p}$ also depend on $a$.

As a direct consequence of Lemma 4.2.3 and equation (4.12), we have the following theorem.

Theorem 4.2.4. Suppose (B1) through (B5) hold. Let $r \in \mathbb{N}$. Then, for $a \in$ $\left(0, \frac{\log \mu\left(\delta_{2}\right)}{\delta_{2}}\right)$ there exist $\theta_{a} \in\left(0, \delta_{2}\right)$ and polynomials $P_{p}(z)$ such that for $f \in F_{r+1, \alpha}^{q+2}$ with $q>\frac{r+1}{2 r_{1}}$ and $\alpha>\delta_{2}$,

$$
\mathbb{E}\left(f\left(S_{N}-a N\right)\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{1}{N^{p+1 / 2}} \int P_{p}(z) f_{\theta_{a}}(z) d z+C_{r+1}^{q+2}\left(f_{\theta_{a}}\right) \cdot o_{r, \theta_{a}}\left(\frac{1}{N^{\frac{r+1}{2}}}\right)
$$

where $f_{\theta}(x)=\frac{1}{2 \pi} e^{-\theta x} f(x), I$ and $\theta_{a}$ as in (4.11).
Remark 4.2.4. In particular, the theorem holds for all $f \in C_{c}^{\infty}(\mathbb{R})$.
This is the weak asymptotic expansion which gives us the required higher order asymptotics for (4.10), the LDP in Theorem 4.2.2.

Next, we replace (B5) by the following stronger assumption which allows us to conclude existence of strong expansions for the LDP. Compare this assumption with assumption (A5) in Chapter 3.
$\widetilde{(\mathrm{B} 5)}$ There are positive real numbers $r_{1}, r_{2}, r_{3}, C, K$ and $N_{0}$ such that for all $\theta \in$ $\left(-\delta_{2}, \delta_{2}\right)$, for all $N>N_{0}$ and for all $|t|>K,\left\|\mathcal{L}_{\theta+i t}^{N}\right\| \leq C \frac{\mu(\theta)^{N}}{N^{r_{2}}|t|^{r_{3}}}$.
As in the case of (B5), we can assume $r_{2}$ and $r_{3}$ to be large after slightly reducing $r_{1}$. Therefore we have the following theorem.

Theorem 4.2.5. Suppose (B1) through (B4) and $\widetilde{(B 5)}$ hold. Let $r \in \mathbb{N}$. Then, there exists $0<\delta_{2} \leq \delta$ such that $S_{N}$ admits a weak asymptotic expansions for the $L D P$ in the range $\left(0, \frac{\log \mu\left(\delta_{2}\right)}{\delta_{2}}\right)$ for $f \in F_{r+1, \alpha}^{1}$ with $\alpha>\delta_{2}$.

In particular, for all $a \in\left(0, \frac{\log \mu\left(\delta_{2}\right)}{\delta_{2}}\right)$ there exist constants $C_{p}(a)$ such that

$$
\mathbb{P}\left(S_{N} \geq a N\right) e^{I(a) N}=\sum_{p=0}^{\lfloor r / 2\rfloor} \frac{C_{p}(a)}{N^{p+1 / 2}}+C_{r, \theta_{a}} o\left(\frac{1}{N^{\frac{r+1}{2}}}\right) .
$$

where

$$
C_{p}(a)=\frac{1}{2 \pi} \int_{0}^{\infty} e^{-\theta_{a} z} P_{p}(z) d z
$$

for some polynomials $P_{0}(z), \ldots, P_{r}(z)$ depending on a and unique $\theta_{a} \in\left(0, \delta_{2}\right)$ such that

$$
I(a)=\sup _{\theta \in\left(-\delta_{2}, \delta_{2}\right)}[a \theta-\log \mu(\theta)]=a \theta_{a}-\log \mu\left(\theta_{a}\right) .
$$

Proof. The proof of the first part is similar to that of Theorem 4.2.4. The only difference is the estimate (4.14).

Since $f \in F_{r+1, \alpha}^{1}$, we have $g=f_{\theta} \in F_{r+1}^{1}$. So $t \widehat{g}(t)=(-i) \widehat{g^{\prime}}(t)$. WLOG assume $r_{3}>\frac{r+1}{2 r_{1}}$. Then,

$$
\begin{aligned}
\left|\int_{|t|>n^{r_{1}}} \widehat{g}(t) \ell\left(\overline{\mathcal{L}}_{\theta_{a}+i t}^{n} v\right) d t\right| \leq C \int_{|t|>n^{r_{1}}}|\widehat{g}(t)|\left\|\overline{\mathcal{L}}_{\theta_{a}+i t}^{n}\right\| d t & \leq C \int_{|t|>n^{r_{1}}}\left|\frac{\widehat{g^{\prime}}(t)}{t^{1+r_{3}}}\right| d t \\
& \leq \frac{C\left\|g^{\prime}\right\|_{1}}{n^{r_{1} r_{3}}} \\
& =\left\|g^{\prime}\right\|_{1} o\left(n^{-(r+1) / 2}\right)
\end{aligned}
$$

Now, the existence of the strong expansion follows from the first part of the theorem and Theorem 4.1.4.

As in the i.i.d. case, $C_{p}(a)$ does not depend $r$ because $\theta_{a}$ and $P_{p}$ do not. Also, there are two ways the coefficients $C(a)$ depend on $a$. First note $\theta_{a}$ depends on the choice of $a$. Also, from Section 3.3, we know exactly how the coefficients of $P_{p}$ depend on the derivatives of the $\mu(z)$ and $\ell\left(\Pi_{z}(\cdot)\right)$ at $\theta_{a}$ and thus, on the exponential
moments of $S_{N}$. Since this dependence of $C(a)$ on $a$ is explicit, one can compute these coefficients.

### 4.3 An application to Markov Chains.

Take $x_{n}$ to be a time homogeneous Markov process on a compact connected manifold $\mathcal{M}$ with smooth transition density $p(x, y)$ which is bounded away from 0 , and $X_{n}=h\left(x_{n-1}, x_{n}\right)$ for smooth function $h: \mathcal{M} \times \mathcal{M} \rightarrow \mathbb{R}$. We assume that $h(x, y)$ can not be written in the form

$$
\begin{equation*}
h(x, y)=H(y)-H(x)+c(x, y) \tag{4.16}
\end{equation*}
$$

where $c(x, y)$ is piece-wise constant. (An equivalent condition is given in Lemma 3.5.1). This is exactly the setting we worked in Section 3.5.3.1.

We need the following lemma to establish (B1) through (B5).
Lemma 4.3.1. Let $K(x, y)$ be a smooth positive function on $\mathcal{M} \times \mathcal{M}$. Let $P$ be an operator on $L^{\infty}(\mathcal{M})$ given by $P u(x)=\int_{\mathcal{M}} K(x, y) u(y) d y$. Then, $P$ has a simple leading eigenvalue $\lambda>0$ and the corresponding eigenfunction $g$ is positive and smooth.

Proof. From the Weierstrass theorem, $K(x, y)$ is a uniform limit of functions of the form $\sum_{r \leq n} J_{r}(x) L_{r}(y)$. Therefore, $P$ can be approximated by finite rank operators. So $P$ is compact. Since $P$ is an operator which leaves the cone of positive functions invariant, by a direct application of Birkhoff Theory (see [2]), $P$ has a leading eigenvalue $\lambda$ which is positive and simple. The corresponding eigenfunction $g$ is also positive.

Because $P$ is compact, there is $l \in(0, \lambda)$ such that $\operatorname{sp}_{L^{\infty}}(P) \cap\{|z|>r\}=\{\lambda\}$. Next, we consider $P$ acting on $C^{1}(\mathcal{M})$. Observe that,

$$
\frac{d}{d x}(P u)(x)=\int_{\mathcal{M}} \frac{\partial K}{\partial x}(x, y) u(y) d y
$$

So, $\|P u(x)\|_{C^{1}} \leq C\|u\|_{\infty}$ for some $C$. Since $\|\cdot\|_{\infty} \leq\|\cdot\|_{C^{1}}$ unit ball with respect to $\|\cdot\|_{C^{1}}$ is relatively compact with respect to $\|\cdot\|_{\infty}$. Therefore the essential spectral radius is 0 by [24, Lemma 2.2]. This gives us, $\operatorname{sp}_{C^{1}}(P) \cap\{|z|>r\} \subseteq\{\lambda\}$.

To see that equality holds, note that the constant function $1 \in C^{1}(\mathcal{M})$. By positivity of $P$,

$$
1 \geq \frac{g}{\sup g} \Longrightarrow P^{n} 1 \geq \frac{P^{n} g}{\sup g} \Longrightarrow P^{n} 1 \geq \frac{\lambda^{n} g}{\sup g} \Longrightarrow\left|\left\|P^{n}\right\|\right| \geq \lambda^{n}\left\|\frac{g}{\sup g}\right\|_{C^{1}} \geq \lambda^{n}
$$

where $|\|\cdot\||$ is the operator norm of $P$ acting on $C^{1}(\mathcal{M})$. Therefore, the spectral radius of $P$ is $\geq \lambda$. This establishes that $g \in C^{1}$. We can repeat the argument and show $g \in C^{r}$ for $r \in \mathbb{N}$.

Take $\mathbb{B}=L^{\infty}(\mathcal{M})$ and consider the family of integral operators,

$$
\left(\mathcal{L}_{z} u\right)(x)=\int_{\mathcal{M}} p(x, y) e^{z h(x, y)} u(y) d y, z \in \mathbb{C} .
$$

Let $\mu$ be the initial distribution of the Markov chain. Then, using the Markov property, we have $\mathbb{E}_{\mu}\left[e^{z S_{n}}\right]=\mu\left(\mathcal{L}_{z}^{N} 1\right)$. Now, we check conditions (B1) through (B5).

Conditions (B1) and (B2) coincide with the conditions (A1) and (A2) in Chapter 3 and we verify them in Section 3.5.3.1. In particular, (B1) holds with $\delta=\infty$. Note that, for all $\theta, \mathcal{L}_{\theta}$ is of the form $P$ in Lemma 4.3.1. Therefore, (B3) holds for all $\theta$. Take $\lambda(\theta)$ be the top eigenvalue and $g_{\theta}$ to be the corresponding eigenfunction. Then, $g_{\theta}$ is smooth.

To show (B4) and (B5) we define a new operator $Q_{\theta}$ as follows.

$$
\left(Q_{\theta} u\right)(x)=\frac{1}{\lambda(\theta)} \int_{\mathcal{M}} e^{\theta h(x, y)} p(x, y) u(y) \frac{g_{\theta}(y)}{g_{\theta}(x)} d(y) .
$$

It is easy see to that $p_{\theta}(x, y)=\frac{e^{\theta h(x, y)} p(x, y)}{g_{\theta}(x) \lambda(\theta)}$ and $d m_{\theta}(y)=g_{\theta}(y) d(y)$ defines a new Markov chain $x_{n}^{\theta}$ with the associated Markov operator $Q_{\theta}$. That is, $Q_{\theta}$ is a positive operator and $Q_{\theta} 1=\frac{1}{\lambda(\theta)} \int_{\mathcal{M}} e^{\theta h(x, y)} p(x, y) \frac{g_{\theta}(y)}{g_{\theta}(x)} d y=1$ because $g_{\theta}$ is the eigenfunction corresponding to eigenvalue $\lambda(\theta)$ of $\mathcal{L}_{\theta}$.

Now, we can repeat the argument in Section 3.5.3.1 to establish properties of the perturbed operator given by

$$
\left(Q_{\theta+i t}\right) u(x)=\int_{\mathcal{M}} e^{i t h(x, y)} p_{\theta}(x, y) d m_{\theta}(y)
$$

Since (4.16) does not hold we conclude that $\operatorname{sp}\left(\mathcal{L}_{\theta+i t}\right) \subset\{|z|<1\}$.
Take $G_{\theta}$ to be the operator on $L^{\infty}(\mathcal{M})$ that corresponds to multiplication by $g_{\theta}$. Then, $\mathcal{L}_{\theta+i t}=\lambda(\theta) G_{\theta} Q_{\theta+i t} G_{\theta}^{-1}$. Therefore, $\operatorname{sp}\left(\mathcal{L}_{\theta+i t}\right)$ is the $\operatorname{sp}\left(Q_{\theta+i t}\right)$ scaled by $\lambda(\theta)$. This implies $\operatorname{sp}\left(\mathcal{L}_{\theta+i t}\right) \subset\{|z|<\lambda(\theta)\}$ as required.

Since (4.16) does not hold, the asymptotic variance $\sigma_{\theta}^{2}$ of $X_{n}^{\theta}=h\left(x_{n-1}^{\theta}, x_{n}^{\theta}\right)$ is positive. Taking $\gamma(\theta+i t)$ to be the top eignevalue of $Q_{\theta+i t}, \lambda(\theta+i t)=\lambda(\theta) \gamma(\theta+i t)$. Thus, $(\log \lambda(\theta))^{\prime \prime}=-\left.\frac{d^{2}}{d t^{2}} \log \lambda(\theta+i t)\right|_{t=0}=-\left.\frac{d^{2}}{d t^{2}} \log \gamma(\theta+i t)\right|_{t=0}=-\frac{\gamma^{\prime \prime}(\theta)}{\gamma(\theta)}+$ $\left(\frac{\gamma^{\prime}(\theta)}{\gamma(\theta)}\right)^{2}=-\gamma^{\prime \prime}(\theta)+\gamma^{\prime}(\theta)^{2}(\because \gamma(\theta)=1)$. Put $S_{N}^{\theta}=X_{1}^{\theta}+\cdots+X_{N}^{\theta}$. Since, $\mathbb{E}\left(e^{i t S_{N}^{\theta}}\right)=\int Q_{\theta+i t}^{N} 1 d \mu$, from (3.37), we have that $\gamma^{\prime}(\theta)^{2}-\gamma^{\prime \prime}(\theta)=\sigma_{\theta}^{2}$. Thus, $(\log \lambda(\theta))^{\prime \prime}=\sigma_{\theta}^{2}>0$. Therefore, $\log \lambda(\theta)$ is a strictly convex function.

Note that, $\mathcal{L}_{\theta}=\lambda(\theta) \Pi_{\theta}+\Lambda_{\theta}$ where $\Pi_{\theta}$ is the projection onto the top eigenspace. From [27, Chapter III], $\Pi_{\theta}=g_{\theta} \otimes \varphi_{\theta}$ where $\varphi_{\theta}$ is the top eigenfunction of $Q_{\theta}^{*}$, the adjoint of $Q_{\theta}$. Because $Q_{\theta}^{*}$ itself is a positive compact operator acting on $\left(L^{\infty}\right)^{*}$ (the
space of finitely additive finite signed measures), $\varphi_{\theta}$ is a finite positive measure. Hence, $\mu\left(\Pi_{\theta} 1\right)=\varphi_{\theta}(1) \mu\left(g_{\theta}\right)>0$ for all $\theta$.

As a result, Lemma 4.2 .1 holds with $\delta_{2}$ arbitrary large and hence, Theorem 4.2.2 holds with $\delta_{2}$ arbitrary large. So the rate function $I(a)$ in Theorem 4.2.2 is finite for $a \in(0, B)$ where $B=\lim _{\theta \rightarrow \infty} \frac{\log \lambda(\theta)}{\theta}$. We observe that $B<\infty$ because $h$ is bounded i.e. $\frac{S_{N}}{N} \leq\|h\|_{\infty}$ surely. In fact, we claim $B=\lim _{N \rightarrow \infty} \frac{B_{N}}{N}$ where $B_{N}=\sup _{x_{0}, \ldots, x_{N}} \sum_{j=1}^{N} h\left(x_{j-1}, x_{j}\right)$ (the supremum taken over all possible realizations of the Markov chain $x_{n}$ ).

First note that $B_{N}$ is subadditive. So $\lim _{N \rightarrow \infty} \frac{B_{N}}{N}$ exists and is equal to $\inf _{N} \frac{B_{N}}{N}$. Given, $a>B$ there exists $N_{0}$ such that for all $N>N_{0}, \frac{S_{N}}{N} \leq \frac{B_{N}}{N}<a$. Therefore, $\mathbb{P}\left(S_{N} \geq a N\right)=0$ for all $N>N_{0}$ and hence, $I(a)=\infty$. Next, given $a<B$, for all $N, B_{N}>a N$. Fix $N$. Then, there exists a realization $x_{1}, \ldots, x_{N}$ such that $a N<\sum_{j=1}^{N} h\left(x_{j-1}, x_{j}\right) \leq B$. Since $h$ is uniformly continuous on $\mathcal{M} \times \mathcal{M}$, there exists $\delta>0$ such that by choosing $y_{j}$ from a ball of radius $\delta$ centred at $x_{j}$ i.e. $y_{j} \in \mathbb{B}\left(x_{j}, \delta\right)$, we have $a N<\sum_{j=1}^{N} h\left(y_{j-1}, y_{j}\right) \leq B$. We estimate the probability of choosing such a realization $y_{1}, \ldots, y_{N}$ and obtain a lower bound for $\mathbb{P}\left(S_{N} \geq a N\right)$ :

$$
\begin{aligned}
\mathbb{P}\left(S_{N} \geq a N\right) & \geq \int_{\mathbb{B}\left(x_{N}, \delta\right)} \ldots \int_{\mathbb{B}\left(x_{1}, \delta\right)} \int_{\mathbb{B}\left(x_{0}, \delta\right)} p\left(y_{N-1}, y_{N}\right) \ldots p\left(y_{0}, y_{1}\right) d \mu\left(y_{0}\right) d y_{1} \ldots d y_{N} \\
& \geq \mu\left(\mathbb{B}\left(x_{0}, \delta\right)\right)\left(\min _{x, y \in \mathcal{M}} p(x, y)\right)^{N} \operatorname{vol}\left(\mathbb{B}_{\delta}\right)^{N}
\end{aligned}
$$

Therefore, $I(a)<\infty$ as required.
Also, because $g_{\theta}$ is smooth we can repeat the argument in Section 3.5.3.1 to obtain (3.45) for $Q_{\theta+i t}$. That is, there is $\epsilon_{\theta}$ and $r_{\theta}$ such that $\left\|Q_{\theta+i t}^{2}\right\| \leq\left(1-\epsilon_{\theta}\right)$ for
all $|t|>r_{\theta}$. Therefore,

$$
\left\|\mathcal{L}_{\theta+i t}^{N}\right\|=\lambda(\theta)^{N}\left\|G_{\theta} Q_{\theta+i t}^{N} G_{\theta}^{-1}\right\| \leq \lambda(\theta)^{N}\left\|G_{\theta}\right\|\left\|Q_{\theta+i t}^{N}\right\|\left\|G_{\theta}^{-1}\right\| \leq C \lambda(\theta)^{N}\left(1-\epsilon_{\theta}\right)^{\lfloor N / 2\rfloor} .
$$

This establishes (B5).
Since the rate in (B5) is exponential and Theorem 4.2.2 holds for $(0, B)$, we conclude that for all $r \in \mathbb{N}$, these Markov chains admit weak expansions for large deviations of order $r$ in the range $(0, B)$ for $F_{r+1, B+}^{3}$ where $B+=\infty$, if $B=\infty$ and $B+>B$, if $B<\infty$.

We need a stronger assumption on $h$ to establish ( $\widetilde{\mathrm{B} 5})$. Suppose,

For all $x, y$ critical points of $z \mapsto(h(x, z)+h(z, y))$ are non-degenerate.

Since critical points of $z \mapsto(h(x, z)+h(z, y))$ are non-degenerate we can use the stationary phase asymptotics in [48, Chapter VIII.2], to obtain,

$$
\left|\int_{\mathcal{M}} e^{i t(h(x, z)+h(z, y)} p(x, z) p(z, y) e^{\theta(h(x, z)+h(z, y))} d z\right| \leq \frac{M}{|t|^{d / 2}}
$$

where $M$ is a constant and $d$ is the dimension of $\mathcal{M}$. Therefore, $\left\|\mathcal{L}_{\theta+i t}^{2}\right\| \leq \frac{M}{|t|^{d / 2}}$. Choose $K=(2 M)^{2 / d}$. Then for all $|t|>K,\left\|Q_{\theta+i t}^{2}\right\| \leq \frac{1}{2}$ and hence,

$$
\left\|\mathcal{L}_{\theta+i t}^{N}\right\| \leq\left\|\mathcal{L}_{\theta+i t}^{N-2}\right\|\left\|\mathcal{L}_{\theta+i t}^{2}\right\| \leq\left(\frac{1}{2}\right)^{\lfloor(N-2) / 2\rfloor} \frac{M}{t^{d / 2}},|t|>K
$$

By convexity, $\lambda(\theta)>1$. Thus,

$$
\left\|\mathcal{L}_{\theta+i t}^{N}\right\| \leq M\left(\frac{1}{2}\right)^{\lfloor(N-2) / 2\rfloor} \frac{\lambda(\theta)^{N}}{t^{d / 2}},|t|>K
$$

This establishes $\widetilde{(\mathrm{B} 5)}$.

In particular, when $h$ depends only on one variable, i.e. $h(x, y)=H(x)$ for some $H$, we have that $h(x, z)+h(z, y)=H(x)+H(z)$. Then, the condition (4.17) reduces to critical points of $H$ being non-degenerate.

Again, because Theorem 4.2.2 hold for all $(0, B)$ and the rate in $(\widetilde{\mathrm{B} 5})$ is exponential, we conclude that these strongly ergodic Markov chains admit strong expansions for large deviations of all orders in the range $(0, B)$.

## Chapter A: Appendix

## A. 1 Convergence of $\mathcal{X}$.

We need some background information. Given a piecewise smooth function $g: \mathbb{R}^{d} \rightarrow \mathbb{R}$ of compact support its Siegel transform is a function on the space of lattices defined by

$$
\mathcal{S}(g)(\mathcal{L})=\sum_{\mathbf{w} \in \mathcal{L} \backslash\{0\}} g(\mathbf{w}) .
$$

We need an identity of Siegel, see ( [38, Section 3.7] or [46, Lecture XV]) saying that

$$
\begin{equation*}
\mathbf{E}_{\mathcal{L}}(\mathcal{S}(g))=\int_{\mathbb{R}^{d}} g(\mathbf{w}) d \mathbf{w} \tag{A.1}
\end{equation*}
$$

In particular, if $B$ is a set in $\mathbb{R}^{d}$ with piecewise smooth boundary not containing $\mathbf{0}$ then

$$
\begin{equation*}
\mathbf{P}_{\mathcal{L}}(\mathcal{L} \cap B \neq \emptyset) \leq \mathbf{P}\left(\mathcal{S}\left(\mathbb{1}_{B}\right)(\mathcal{L}) \geq 1\right) \leq \mathbf{E}_{\mathcal{L}}\left(\mathcal{S}\left(\mathbb{1}_{B}\right)\right)=\operatorname{Vol}(B) \tag{A.2}
\end{equation*}
$$

Proof of Lemma 2.1.2. Let $\mathcal{L}^{+}=\{\mathbf{w} \in \mathcal{L}: y(\mathbf{w})>0\}$. Since $\frac{\sin (2 \pi \chi(\mathbf{w}))}{y(\mathbf{w})}$ is even it is enough to restrict the attention to $\mathbf{w} \in \mathcal{L}^{+}$.

Throughout the proof we fix two numbers $\varepsilon>0, \tau<1$ such that $\varepsilon \ll(1-\tau) \ll$ 1. It is easy to see using (A.2) and Borel-Cantelli Lemma that for almost every lattice
$\mathcal{L}$, there exists $C$ and $\beta$ such that $y(\mathbf{w})>\frac{C}{\|\mathbf{w}\|^{\beta}}$. It follows that

$$
\sum_{\mathbf{w} \in \mathcal{L}^{+}:\|x(\mathbf{w})\| \geq\|\mathbf{w}\|^{\varepsilon}} \frac{\sin 2 \pi \chi(\mathbf{w})}{y(\mathbf{w})} e^{-\|x(\mathbf{w})\|^{2}} \leq \sum_{\mathbf{w} \in \mathcal{L}^{+}} C\|\mathbf{w}\|^{\beta} e^{-\|\mathbf{w}\|^{2 \varepsilon}}
$$

converges absolutely. Hence it suffices to establish the convergence of

$$
\overline{\mathcal{X}}:=\sum_{\mathbf{w} \in \mathcal{L}^{+}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon}<R^{\varepsilon}} \frac{\sin 2 \pi \chi(\mathbf{w})}{y(\mathbf{w})} e^{-\|x(\mathbf{w})\|^{2}}
$$

Let $R_{j, k}=2^{k}+j 2^{\tau k}, j=0, \ldots 2^{(1-\tau) k}$. To prove the convergence of $\overline{\mathcal{X}}$ we will show that for all $\mathcal{L}$ almost all $\chi$ satisfy two estimates below

$$
\begin{align*}
& \forall \text { sequence }\left\{j_{k}\right\} \overline{\mathcal{X}}_{R_{j_{k}, k}} \text { converges as } k \rightarrow \infty  \tag{A.3}\\
& \max _{j} \sup _{R_{j, k} \leq R \leq R_{j+1, k}}\left|\overline{\mathcal{X}}_{R}-\overline{\mathcal{X}}_{j, k}\right| \rightarrow 0 \text { as } k \rightarrow \infty . \tag{A.4}
\end{align*}
$$

To prove (A.3) let

$$
S_{j, k}=\sum_{\mathbf{w} \in \mathcal{L}^{+}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon}, R_{j, k} \leq\|\mathbf{w}\| \leq R_{j+1, k}} \frac{\sin 2 \pi \chi(\mathbf{w})}{y(\mathbf{w})} e^{-\|x(\mathbf{w})\|^{2}}
$$

Using that $\mathbf{E}_{\chi}(\sin (2 \pi(\chi(\mathbf{w}))))=0$ and for $\mathbf{w}_{1} \neq \pm \mathbf{w}_{2}$ we have

$$
\mathbf{E}_{\chi}\left(\sin \left(2 \pi\left(\chi\left(\mathbf{w}_{1}\right)\right)\right) \sin \left(2 \pi\left(\chi\left(\mathbf{w}_{2}\right)\right)\right)\right)=0
$$

we see that $\mathbf{E}_{\chi}\left(S_{j, k}\right)=0$ and

$$
\begin{aligned}
\operatorname{Var}_{\chi}\left(S_{j, k}\right) & =\sum_{\mathbf{w} \in \mathcal{L}^{+}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon}, R_{j, k} \leq\|\mathbf{w}\| \leq R_{j+1, k}} \frac{e^{-2\|\mathbf{x}(\mathbf{w})\|^{2}}}{2 y^{2}(\mathbf{w})} \\
& \leq \frac{1}{2^{2 k+1}} \operatorname{Card}\left(\mathbf{w}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon}, R_{j, k} \leq\|\mathbf{w}\| \leq R_{j+1, k}\right) \\
& \leq \frac{C(\mathcal{L})}{2^{2 k}} \operatorname{Vol}\left(\mathbf{w}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon}, R_{j, k} \leq\|\mathbf{w}\| \leq R_{j+1, k}\right) \\
& \leq C(\mathcal{L}) 2^{(\tau+\varepsilon(d-1)-2) k}
\end{aligned}
$$

Hence by Chebyshev inequality for each $j$

$$
\mathbf{P}_{\chi}\left(S_{j, k} \geq 2^{-(1-\tau+\varepsilon) k}\right) \leq C(\mathcal{L}) 2^{(\varepsilon d-\tau) k}
$$

and so

$$
\mathbf{P}_{\chi}\left(\exists j: S_{j, k} \geq 2^{-(1-\tau+\varepsilon) k}\right) \leq C(\mathcal{L}) 2^{(1+\varepsilon d-2 \tau) k}
$$

Thus if $\varepsilon$ is sufficiently small and $\tau$ is sufficiently close to 1 then Borel-Cantelli Lemma shows that for almost every $\chi$, if $k$ is large enough, then for all $j S_{j, k} \leq$ $2^{-(1-\tau+\varepsilon) k}$ and thus $\sum_{j} S_{j, k} \leq 2^{-\varepsilon k}$ proving (A.3). Likewise,

$$
\begin{aligned}
\sup _{R_{j, k} \leq R \leq R_{j+1, k}} & \left|\overline{\mathcal{X}}_{R}-\overline{\mathcal{X}}_{j, k}\right| \\
& \leq \sum_{\mathbf{w} \in \mathcal{L}^{+}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon},\|\mathbf{w}\| \in\left[R_{j, k}, R_{j+1, k}\right]} \frac{1}{|y(\mathbf{w})|} e^{-\|x(\mathbf{w})\|^{2}} \\
& \leq C(\mathcal{L}) 2^{-2 k} \operatorname{Vol}\left(\mathbf{w}:\|x(\mathbf{w})\| \leq\|\mathbf{w}\|^{\varepsilon}, R_{j, k} \leq\|\mathbf{w}\| \leq R_{j+1, k}\right) \\
& \leq \bar{C}(\mathcal{L}) 2^{\tau+\varepsilon(d-1)-1}
\end{aligned}
$$

proving (A.4). Lemma 2.1.2 is established.

## A. 2 Hierarchy of Expansions.

In the discussion below, we do not assume the abstract setting introduced in section 3.1. Therefore the hierarchy of asymptotic expansions provided here holds true in general.

We observe that the classical Edgeworth expansion is the strongest form of asymptotic expansion among the expansions for non-lattice random variables. The following proposition and remark A.2.1 establish this fact.

Proposition A.2.1. Suppose $S_{N}$ admits order $r$ Edgeworth expansions, then it also admits order $r$ weak global expansion for $f \in F_{0}^{1}$ and order $r$ averaged expansions for $f \in L^{1}$. Further, if the polynomials $P_{p}$ in the Edgeworth expansion has opposite parity as $p$ then $S_{N}$ admits order $r-1$ weak local expansion for $f \in F_{r}^{1}$.

Remark A.2.1. Section 3.5.2 contains examples for which the weak and averaged forms of expansions exist but the strong expansion does not. Therefore none of the above implications are reversible.

Proof of Proposition A.2.1. Suppose $f \in F_{0}^{1}$. Let $F_{n}=\mathbb{P}\left(\frac{S_{n}-n A}{\sqrt{n}} \leq x\right)$ and put

$$
\mathcal{E}_{r, n}(x)=\mathfrak{N}(x)+\sum_{p=1}^{r} \frac{P_{p}(x)}{n^{p / 2}} \mathfrak{n}(x) .
$$

Observe that $F_{n}(x)-\mathcal{E}_{n}(x)=o\left(n^{-r / 2}\right)$ uniformly in $x$ and,

$$
d \mathcal{E}_{r, n}(x)=\mathfrak{n}(x) d x+\sum_{p=1}^{r} \frac{1}{n^{p / 2}}\left[P_{p}^{\prime}(x) \mathfrak{n}(x)+P_{p}(x) \mathfrak{n}^{\prime}(x)\right] d x=\sum_{p=0}^{r} \frac{R_{p}(x)}{n^{p / 2}} \mathfrak{n}(x) d x
$$

where $R_{p}$ are polynomials given by $R_{p}=P_{p}^{\prime}+P_{p} Q$ and $Q$ is such that $\mathfrak{n}^{\prime}(x)=$
$Q(x) \mathfrak{n}(x)$. Next, we observe that,

$$
\begin{aligned}
\mathbb{E}\left(f\left(S_{n}-n A\right)\right) & =\mathbb{E}\left(f\left(\frac{S_{n}-n A}{\sqrt{n}} \sqrt{n}\right)\right)=\int f(x \sqrt{n}) d F_{n}(x) \\
& =\int f(x \sqrt{n}) d \mathcal{E}_{r, n}(x)+\int f(x \sqrt{n}) d\left(F_{n}-\mathcal{E}_{r, n}\right)(x)
\end{aligned}
$$

Now we integrate by parts and use $\mathcal{E}_{r, n}(\infty)=F_{n}(\infty)=1$ and $\mathcal{E}_{r, n}(-\infty)=F_{n}(-\infty)=$ 0 to obtain,

$$
\begin{aligned}
\mathbb{E}\left(f\left(S_{n}-n A\right)\right) & =\int f(x \sqrt{n}) d \mathcal{E}_{r, n}(x)+\left.\left(F_{n}-\mathcal{E}_{r, n}\right)(x) f(x \sqrt{n})\right|_{-\infty} ^{\infty} \\
& \quad-\int\left(F_{n}-\mathcal{E}_{r, n}\right)(x) \sqrt{n} f^{\prime}(x \sqrt{n}) d x \\
& =\int \sum_{p=0}^{r} \frac{1}{n^{p / 2}} R_{p}(x) \mathfrak{n}(x) f(x \sqrt{n}) d x+o\left(n^{-r / 2}\right) \int \sqrt{n} f^{\prime}(x \sqrt{n}) d x \\
& =\sum_{p=0}^{r} \frac{1}{n^{p / 2}} \int R_{p}(x) \mathfrak{n}(x) f(x \sqrt{n}) d x+o\left(n^{-r / 2}\right) .
\end{aligned}
$$

This is the order $r$ weak global Edgeworth expansion. The existence of the order $r-1$ weak local expansion follows from this. This is our next theorem. So we postpone its proof.

For $f \in L^{1}$ substituting $x$ by $x+\frac{y}{\sqrt{n}}$ in the Edgeworth expansion for $S_{n}$ we have

$$
\begin{aligned}
\mathbb{P}\left(\frac{S_{n}-n A}{\sqrt{n}} \leq x+\frac{y}{\sqrt{n}}\right)- & \mathfrak{N} \\
= & \left(x+\frac{y}{\sqrt{n}}\right) \\
& =\sum_{p=1}^{r} \frac{1}{n^{p / 2}} P_{p}\left(x+\frac{y}{\sqrt{n}}\right) \mathfrak{n}\left(x+\frac{y}{\sqrt{n}}\right)+o\left(n^{-r / 2}\right) .
\end{aligned}
$$

For fixed $x$, the error is uniform in $y$. Therefore, multiplying the equation by $f(y)$ and then integrating we can conclude that the order $r$ averaged expansion exists.

Remark A.2.2. We have seen from the derivation of the Edgeworth expansion in section 3.2 that $P_{p}(x)$ and $p$ have opposite parity in the weakly dependent case. This
implies that $P_{p, g}$ has the same parity as $p$. This is true in the i.i.d. case as well. Even though this assumption may look artificial in the general case, it is reasonable. When using characteristic functions to derive the expansions, one is likely to end up with Hermite polynomials which is the reason behind the parity relation.

Next, we compare the the relationships among the weak and averaged forms of Edgeworth expansions.

Proposition A.2.2. Suppose $S_{N}$ admits order $r$ weak global Edgeworth expansion for $f \in F_{r}^{q+1}$ for some $q \geq 0$. If the polynomials $P_{p, g}$ in the global Edgeworth expansion has the same parity as $p$ then $S_{N}$ admits order $r-1$ weak local expansion for $f$.

Proof. Assume, $f \in F_{r}^{1}$. Then, from the Plancherel formula,

$$
\int_{\mathbb{R}} \sqrt{n} f(x \sqrt{n}) P_{p, g}(x) \mathfrak{n}(x) d x=\frac{1}{2 \pi} \int_{\mathbb{R}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) A_{p}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t
$$

where $A_{p}(t)$ are polynomials constructed using the following relation,

$$
P_{p, g}(t) e^{-\frac{t^{2}}{2 \sigma^{2}}}=\frac{1}{\sqrt{2 \pi \sigma^{2}}} A_{p}\left(-i \frac{d}{d t}\right)\left[e^{-\frac{t^{2}}{2 \sigma^{2}}}\right] .
$$

By construction $P_{p, g}$ and $A_{p}$ has the same parity. This means $A_{p}$ has the same parity as $p$.

First replace

$$
\int P_{p, g}(x) \mathfrak{n}(x) f(x \sqrt{n}) d x
$$

by

$$
\frac{1}{2 \pi \sqrt{n}} \int_{\mathbb{R}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) A_{p}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t
$$

in the weak global expansion to obtain,

$$
\sqrt{n} \mathbb{E}\left(f\left(S_{n}-n A\right)\right)=\frac{1}{2 \pi} \sum_{p=0}^{r} \frac{1}{n^{p / 2}} \int_{\mathbb{R}} \widehat{f}\left(\frac{t}{\sqrt{n}}\right) A_{p}(t) e^{-\frac{\sigma^{2} t^{2}}{2}} d t+o\left(n^{-(r-1) / 2}\right)
$$

Then substituting for $\widehat{f}$ with its order $r-1$ Taylor expansion,

$$
\sqrt{n} \mathbb{E}\left(f\left(S_{n}-n A\right)\right)=\frac{1}{2 \pi} \sum_{p=0}^{r} \sum_{j=0}^{r-1} \frac{\widehat{f}^{(j)}(0)}{j!n^{(j+p) / 2}} \int_{\mathbb{R}} t^{j} e^{-\sigma^{2} t^{2} / 2} A_{p}(t) d t+o\left(n^{-(r-1) / 2}\right)
$$

Put

$$
a_{p j}=\int_{\mathbb{R}} t^{j} e^{-\sigma^{2} t^{2} / 2} A_{p}(t) d t=0 \quad \text { and } \quad f^{(j)}(0)=\int(-i t)^{j} f(t) d t
$$

to get,

$$
\sqrt{n} \mathbb{E}\left(f\left(S_{n}-n A\right)\right)=\frac{1}{2 \pi} \sum_{p=0}^{r} \sum_{j=0}^{r-1} \frac{a_{p j}}{j!n^{(j+p) / 2}} \int_{\mathbb{R}}(-i t)^{j} f(t) d t+o\left(n^{-(r-1) / 2}\right)
$$

Since $p$ and $A_{p}$ are of the same parity, when $j+p$ is odd. $a_{p j}=0$. So we collect terms such that $p+j=2 k$ where $k=0, \ldots, r-1$ and write,

$$
P_{k, w}=\sum_{p+j=2 k} \frac{a_{p j}}{j!}(-i t)^{j}
$$

Then, rearranging, simplifying and absorbing higher order terms to the error, we obtain,

$$
\sqrt{n} \mathbb{E}\left(f\left(S_{n}-n A\right)\right)=\frac{1}{2 \pi} \sum_{k=0}^{\lfloor(r-1) / 2\rfloor} \frac{1}{n^{k}} \int_{\mathbb{R}} P_{k, w}(t) f(t) d t+o\left(n^{-(r-1) / 2}\right)
$$

which is the order $r-1$ weak local Edgeworth expansion.

## A. 3 Construction of $\left\{f_{k}\right\}$.

For each $k$, let $f_{k}(x)=\frac{1}{\pi} \tan ^{-1}(k x)+\frac{1}{2}$ for $x \in[-1, k]$. Extend $f_{k}$ to $[-2, k+1]$ in such a way that $f_{k}(-2)=f_{k}(k+1)=0, f_{k}$ is continuously differentiable and satisfying the following conditions.

1. $f_{k}$ is increasing on $[-2, k]$ with derivative on $[-2,-1]$ is bounded above by 1 .
2. $f_{k}$ is decreasing on $[k+1 / 2, k+1]$ with derivative bounded below by -5 .
3. $\left|f_{k}^{\prime}\right| \leq 5$ on $[k, k+1]$.
4. $0 \leq f_{k} \leq 1$ on $[-2, k+1]$ and $f_{k}=0$ elsewhere.

Then, $f_{k}$ is supported on $[-2, k+1]$. Here our choice of bounds 1 and -5 in some sense arbitrary. As long as they are large enough and independent of $k$, we obtain an appropriate sequence of functions.

As an example, when $k=5$, the graph of $f_{5}$ looks like:


For all $\gamma>0$,

$$
\int\left|\left(f_{k}\right)_{\gamma}(x)\right| d x=\int\left|e^{-\gamma x} f_{k}(x)\right| d x \leq \int_{-2}^{\infty} e^{-\gamma x} d x=C_{\gamma, 1}<\infty
$$

because $0 \leq f_{k} \leq 1$.

Since $\left|f_{k}^{\prime}\right| \leq 5$ on $[k, k+1], 0 \leq f_{k} \leq 1$ and $f_{k}$ is increasing on $[-2, k]$,

$$
\begin{aligned}
\int\left|\left(\left(f_{k}\right)_{\gamma}\right)^{\prime}(x)\right| d x & =\int_{-2}^{k+1}\left|\gamma e^{-\gamma x} f_{k}(x)+e^{-\gamma x} f_{k}^{\prime}(x)\right| d x \\
& \leq \int_{-2}^{k+1}\left(\gamma e^{-\gamma x} f_{k}(x)+e^{-\gamma x}\left|f_{k}^{\prime}(x)\right|\right) d x \\
& \leq \int_{-2}^{k} \gamma e^{-\gamma x} d x+\int_{-1}^{k} f_{k}^{\prime}(x) d x+\int_{k}^{k+1}\left(\gamma e^{-\gamma x}+5 e^{-\gamma x}\right) d x \\
& \leq 1+\int_{-2}^{k+1}(5+\gamma) e^{-\gamma x} d x=C_{\gamma, 2}<\infty
\end{aligned}
$$

Also, note that $\left|x^{l} f_{k}(x)\right| \leq x^{l} e^{-\gamma x}$ for all $x \in[-2, k+1]$. Hence,

$$
\int\left|x^{l} f_{k}(x)\right| d x \leq \int_{-2}^{\infty} x^{l} e^{-\gamma x} d x=J_{\gamma, l}<\infty
$$

Put $J_{r}(\gamma)=\max _{1 \leq l \leq r} J_{\gamma, l}$ and $C_{\gamma}(r)=\max \left\{J_{r}(\gamma), C_{\gamma, 1}, C_{\gamma, 2}\right\}$. Then, $C_{\gamma}(r)$ is finite and depends only on $\gamma$ and $r$.

Now, we have the following,

1. $C_{r+1}^{1}\left(\left(f_{k}\right)_{\gamma}\right) \leq C_{\gamma}(r)$ for all $k$.
2. Since, $\frac{1}{\pi} \tan ^{-1}(k x)+\frac{1}{2}$ converges pointwise to $1_{[0, \infty)}(x)$, it is easy to see that $f_{k} \rightarrow 1_{[0, \infty)}$ pointwise.
3. Since for each $p, e^{-\gamma z} P_{p}(z) f_{k}(z)$ converges pointwise to $e^{-\gamma z} P_{p}(z) 1_{[0, \infty)}(z)$, $e^{-\gamma z}\left|P_{p}(z)\right| 1_{[-2, \infty)}$ is integrable and $\left|e^{-\gamma z} P_{p}(z) f_{k}(z)\right| \leq e^{-\gamma z}\left|P_{p}(z)\right| 1_{[-2, \infty)}$, we can apply the LDCT to conclude,

$$
\int P_{p}(z) g_{k}(z) d z=\int_{-2}^{\infty} e^{-\gamma z} P_{p}(z) f_{k}(z) d z \rightarrow \int_{0}^{\infty} e^{-\gamma z} P_{p}(z) d z
$$

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