Limit theorems for random walks on a strip in subdiffusive regimes

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Abstract

We study the asymptotic behaviour of occupation times of a transient random walk (RW) in a quenched random environment (RE) on a strip in a subdiffusive regime. The asymptotic behaviour of hitting times, which is a more traditional object of study, is exactly the same. As a particular case, we solve a long standing problem of describing the asymptotic behaviour of a RW with bounded jumps on a one-dimensional lattice. Our technique results from the development of ideas from our previous work (Dolgopyat and Goldsheid 2012 Commun. Math. Phys. 315 241–77) on the simple RWs in RE and those used in Bolthausen and Goldsheid (2000 Commun. Math. Phys. 214 429–47; 2008 Commun. Math. Phys. 278 253–88) and Goldsheid (2008 Probab. Theory Relat. Fields 141 471–511) for the study of random walks on a strip.

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1. Introduction

The main goal of this work is to describe the asymptotic behaviour of a random walk (RW) in a quenched random environment (RE) on a strip in a subdiffusive regime. As a corollary we obtain a solution to a long standing problem about the asymptotic behaviour of a RW with bounded jumps in RE on a one-dimensional lattice. These two models are natural generalizations of the one-dimensional RWRE with jumps to the nearest neighbors—the so called simple RWRE (SRWRE). The techniques and ideas used in this paper resulted from the development and combination of those used in [6], where we studied the limiting behaviour of the SRWRE, and in [1, 2, 12], which studied RWRE on a strip. Our main model is the RWRE on a strip and the main quantitative characteristic of the walk that is the occupation time $T_N$ of a large box (see (1.10) for exact definition). In [6] we also studied $T_N$, but on a strip the
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approach we use is very different from the one used for SRWRE. The important difference between SRWRE and other models can be roughly explained by the fact that a transient simple walk has to visit every point on its way to $\infty$, while on a strip it can miss any point with a positive probability. Due to this fact, the expectations of the occupation times of the sites form a Markov process in the ‘simple’ case but this is not true for a walk on a strip. In order to resolve these difficulties, we have to use methods inspired by the theory of dynamical systems such as products of random transformations, Lyapunov exponents, transfer operators combined with more probabilistic techniques such as coupling, large deviations, Poisson processes etc. We believe that the new point of view presented in this paper makes the proofs more transparent even in the classical SRWRE setting.

We now recall the exact definitions of all three models.

**Model 1.** In the simplest 1D case, a random environment is a sequence of independent identically distributed (i.i.d.) random variables $\omega = \{p_n\}_{n \in \mathbb{Z}}$, where $p_n$ are viewed as probabilities of jumps from $n$ to $n+1$. Given $\omega$ and $X_0 = z$, one defines a Markov chain $X_t$, $t = 0, 1, \ldots,$ on $\mathbb{Z}$ with a transition kernel given by

$$P_\omega(X_{t+1} = k+1|X_t = k) = p_k, \quad P_\omega(X_{t+1} = k-1|X_t = k) = 1 - p_k. \quad (1.1)$$

**Model 2.** The RWRE on a strip $S \coloneqq \mathbb{Z} \times \{1, \ldots, m\}$ was introduced in [1] and will be the main object of our study. We say that the set $L_n \coloneqq \{(n, j) : 1 \leq j \leq m\} \subset S$ is the layer $n$ of the strip (or just layer $n$). The walker is allowed to jump from a site in $L_n$ only to a site in $L_{n-1}$, $L_n$, or $L_{n+1}$. Let $X_t = (Z_t, Y_t)$ denote the coordinate of the walk at time $t$, where $t = 0, 1, 2, \ldots, Z_t \in \mathbb{Z}, 1 \leq Y_t \leq m$. An environment $\omega$ on a strip is a sequence of triples of $m \times m$ matrices $\omega = \{P_n, Q_n, R_n\}_{n \in \mathbb{Z}}$ with non-negative matrix elements and such that $P_n + Q_n + R_n$ is a stochastic matrix:

$$(P_n + Q_n + R_n)\mathbf{1} = 1, \quad (1.2)$$

where $\mathbf{1}$ is a vector whose all components are equal to 1. The transition kernel of the walk is given by

$$P_\omega(X_{t+1} = z'|X_t = z) = \begin{cases} P_n(i, j) & \text{if } z = (n, i), z' = (n + 1, j), \\ Q_n(i, j) & \text{if } z = (n, i), z' = (n - 1, j), \\ R_n(i, j) & \text{if } z = (n, i), z' = (n, j). \end{cases} \quad (1.3)$$

The corresponding Markov chain is completely defined if we set $X(0) = z$.

Throughout the paper we suppose that the following conditions are satisfied:

$$(P_n, Q_n, R_n)_{n \in \mathbb{Z}} \text{ is an i.i.d. sequence} \quad (1.4)$$

There is an $\varepsilon > 0$ such that $P$-almost surely for all $i, j \in [1, m]$

$$\|R_n\| < 1 - \varepsilon, \quad ((I - R_n)^{-1} P_n)(i, j) > \varepsilon, \quad ((I - R_n)^{-1} Q_n)(i, j) > \varepsilon. \quad (1.5)$$

**Remarks.**

1. The matrices $P_n$, $Q_n$, and $R_n$ are comprised of probabilities of jumps from sites in $L_n$ to sites in $L_{n+1}$, $L_{n-1}$, and $L_n$ respectively. Condition (1.2) is equivalent to ‘the nearest layer jumps only’ property of the walk.

2. Note that $((I - R_n)^{-1} P_n)(i, j)$ and $((I - R_n)^{-1} Q_n)(i, j)$ are the probabilities for a RW starting from $(n, i)$ to reach $(n + 1, j)$ and, respectively, $(n - 1, j)$ at its first exit from layer $n$. 

We chose to work under conditions (1.5) in order to simplify the proofs. In fact all main results can be proved under the following much milder conditions.

There is $\epsilon > 0$ and integer $l \geq 1$ such that $P$-almost surely $\forall i \in [1, m]$

$$\|R^l_n\| < 1 - \epsilon, \quad ((I - R_n)^{-1}P_n)(i, 1) > \epsilon, \quad ((I - R_n)^{-1}Q_n)(i, 1) > \epsilon. \quad (1.6)$$

Let us describe explicitly the probability spaces hidden behind the above definitions. By $(\Omega, \mathcal{F}, P)$ we denote the probability space describing REs, where $\Omega$ is the set of all environments, $\mathcal{F}$ is the natural sigma-algebra of subsets of $\Omega$ and $P$ is a probability measure on $(\Omega, \mathcal{F})$. The RWRE is specified by the choice of $\Omega$ and $P$. Next, let $X_z = \{X(\cdot) : X(0) = z\}$ be the space of all trajectories of the walk starting from $z \in L_0$. A quenched (fixed) environment $\omega$ thus provides us with a conditional probability measure $P_{\omega,z}$ on $X_z$ with a naturally defined probability space $(X_z, \mathcal{F}_{X_z}, P_{\omega,z})$. In turn, these two measures generate a semi-direct product measure $P_z := P \ltimes P_{\omega,z}$ which is the annealed probability measure on $(\Omega \times X_z, \mathcal{F} \times \mathcal{F}_{X_z})$.

The expectations with respect to $P_{\omega,z}$, $P$, and $P_z$ will be denoted by $E_{\omega,z}$, $E$, and $E_z$, respectively.

Remark. The notations $X_z$, $P_{\omega,z}$, $E_z$, etc emphasize the dependence of these objects on the starting point $z$ of the walk. However, we often use the simplified version of these notations such as $P_{\omega}$, $E_{\omega}$, etc because the asymptotic behaviour of the walk does not depend on $z$ and it is usually clear from the context what the starting point of the walk is.

Model 3. The RW on $\mathbb{Z}$ with uniformly bounded jumps is another natural extension of the nearest neighbour model. The RE $\omega$ defined as $\{p(x) = (p(x, k))_{-m \leq k \leq m} \}_{x \in \mathbb{Z}}$, where $p(x)$ is a stationary in $x$ sequence of vectors with $\sum_{k=-m}^{m} p(x, k) = 1$ and $p(x, k) \geq 0$. For a given environment $\omega$ the transition kernel of the walk is defined by

$$P_\omega(X(t+1) = x + k \mid X(t) = x) = p(x, k), \quad x \in \mathbb{Z}. \quad (1.7)$$

The following geometric construction transforms this walk into a walk on a strip. Let us view $\mathbb{Z}$ as a subset of the $X$-axis in a two-dimensional plane. Cut the $X$-axis into equal intervals of length $m$ so that each of them contains exactly $m$ consecutive integer points. Turn each such interval around its left most integer point anti-clockwise by $\pi/2$. The image of $\mathbb{Z}$ obtained in this way is a part of a strip with distances between layers equal to $m$. Re-scaling the $X$-axis of the plane by $m^{-1}$ makes the distance between these layers equal to one and the RW on $\mathbb{Z}$ transforms into a RW on a strip with jumps to the nearest layers only. The relevant formulae for matrices $P_n$, $Q_n$, $R_n$ can be found in [1], where this construction was described in a more formal way.

It is obvious that if $p(x)$, $x \in \mathbb{Z}$, is an i.i.d sequence then the just defined triples of matrices $(P_n, Q_n, R_n)$ are i.i.d. It is also easy to see that (1.5) is satisfied if for some $\epsilon > 0$

$$P[p(x, 1) > \epsilon, \ p(x, -1) > \epsilon, \ p(x, m) > \epsilon, \ p(x, -m) > \epsilon] = 1. \quad (1.8)$$

A much wider class of one-dimensional RW with bounded jumps is obtained if instead of (1.8) we suppose only that

$$P[p(x, 1) > \epsilon, \ p(x, -1) > \epsilon] = 1. \quad (1.9)$$

In this case (1.5) may not be satisfied but (1.6) is satisfied.
**Brief comments on the history of the subject.** Two pioneering papers which initiated the development of the theory of RWRE were published in 1975 by Solomon [31] and Kesten *et al* [17]. In [31] the asymptotic properties of the SRWRE were discussed at the level of the Law of Large Numbers and the surprising fact that for a wide class of parameters the SRWRE would be escaping to $\infty$ at a zero speed was discovered. In [17] the limiting distributions of the hitting times and of the position of $X$ were found in the annealed setting. The extensions of the main results from these papers to the RWRE on a strip are explained below in theorems 3, 4, and 7.

In 1982, Sinai [29] described the asymptotic behaviour of a recurrent SRWRE. He discovered a phenomena which is now called the Sinai diffusion.

The methods used in [17, 29, 31] rely heavily on the *jumps to the nearest neighbours only* property of the walk and the limiting distributions described in [17] were obtained for annealed RWRE. Therefore the following questions arose and were known essentially since 1975:

1. Can one describe the limiting behaviour of the quenched RW at least in the case of the SRWRE (model 1)?
2. What are the analogues of (a) P-almost sure results from [31], (b) the annealed limiting statements from [17] for more general models, e.g. model 3?
3. What can be said about more general classes of environments, say stationary environments with appropriate mixing properties?

In the 1982 paper Sinai explicitly stated the questions about the possibility to extend his results to more general models, such as model 3.

The attempts to find answers to question 1 are relatively recent. We shall not discuss them here in any detail. The references concerned with SRWRE along with relevant discussion can be found in [6, 11].

Partial answers to question 2 were obtained in [3–5, 18, 21, 22]. The discussion of these results can be found in [1, 12].

Question 3 was addressed in several publications, see e.g. [1, 11, 12, 23, 33]. And even though in [6] and in this work we consider the so called i.i.d environments (as defined above) we believe that the methods we use are useful for the analysis of RW in stationary RE satisfying appropriate mixing conditions.

Finally, let us mention several results on the RWRE on a strip which are directly related to this work. The criterion for recurrence and transience has been found in [1]. A detailed description of the limiting behaviour in the recurrent regime was given in [2]. A criterion for linear growth and the quenched (and hence annealed) central limit theorem (CLT) was obtained in [12] for wide classes of environments; in particular the CLT for hitting times was established for stationary environments.

**Quantities characterizing the asymptotic behaviour of a RWRE.** Remember that $X_t = (Z_t, Y_t)$ is the coordinate of the walk at time $t$ with $Z_t$ being its $\mathbb{Z}$ component. Denote by $\tilde{T}_N$ the hitting time of layer $L_N$—the time at which the walk starting from a site in $L_0$ reaches $L_N$ for the first time. It is both natural and in the tradition of the field to consider the understanding of the main asymptotic properties of the walk as achieved if the asymptotic behaviour of $Z_t$ as $t \to \infty$ and $\tilde{T}_N$ as $N \to \infty$ is known.

There is of course a strong connection between the asymptotic behaviour of $Z_t$ and $\tilde{T}_N$. Obviously $\tilde{T}_N$ is strictly monotone in $N$ and $Z_{\tilde{T}_N} = N$. This and some other, less trivial relations between these random variables were used in a very efficient way in the study of transient RWs already in [17, 31]. In particular in [17] the asymptotic distribution of $Z_t$ was deduced from that of $\tilde{T}_N$. 
In our recent work [6] on SRWRE we studied a different quantity as the main way of describing the asymptotic behaviour of the RW. Namely, we considered the occupation time $T_N$ of the interval $[0, N-1]$. The asymptotic behaviour of $\tilde{T}_N$ is exactly the same as that of $T_N$ since $|T_N - \tilde{T}_N|$ is a stochastically bounded random variable (see lemma 2.1 from [6]). In this paper, we study a similar quantity—the occupation time of a box $[L_0, L_{N-1}] \overset{\text{def}}{=} \{(n, i) : 0 \leq n \leq N-1\}$.

**Definition.** The occupation time $T_N$ of the box $[L_0, L_{N-1}]$ is the total time the walk $X_t$ starting from a site in $L_0$ spends on this box during its lifetime. In other words

$$T_N = \#\{t: 0 \leq t < \infty, \ X_t \in [L_0, L_{N-1}]\}, \quad (1.10)$$

**Remark.** Note that $T_N \equiv T_{N,z}$ depends on the starting point $z$ of the walk. Also, we use the convention that starting from a site $z$ counts as one visit to $z$.

The paper is organized as follows. We start (section 2) by reviewing the results from [1, 12] which are used in this paper. In section 3 we derive formulae for the expected value of occupation times and state their asymptotic properties; the latter play a major role in the analysis of the asymptotic behaviour of the RW on a strip. In section 4 we define traps and state the main results of the paper (theorems 5 and 6) which are followed by theorem 7 extending to the case of the strip the classical results from [17]. Section 5 is devoted to the proof of the properties of traps followed by the derivation of theorem 5. The proof of theorem 6 is given in section 6. Since this proof is similar to that of the main result in [6], we focus our attention on the differences which are due to the fact that this time we deal with a strip. Section 8 contains the extensions of our results which are not needed in the analysis of the hitting time but are important for understanding of other properties of RWRE (see [8, 10, 16, 20] for related work in the context of SRWRE) and will be used in the future work. The paper has four appendices containing results which are not specific to RWRE. Most of these results are not completely new, but we present them in the form convenient for our purposes. Namely, appendix A contains the estimates of occupation times for general transient Markov chains. Standard facts about the Poisson processes and their relation to stable laws are collected in appendix B. In appendix C we prove a renewal theorem for a system of random contractions. The fact that the assumptions of appendix C are applicable in our setting is verified in section 7. Appendix D contains results about the mixing properties of RWs on the strip satisfying ellipticity conditions.

Some conventions and notations. Letters $C$, $\tilde{C}$, $c$, $\tilde{c}$ denote positive constants, $\varepsilon$ is a strictly positive and small enough number, and $\theta$ is a constant from the interval $(0, 1)$. The values of all these constants may be different in different sections of the paper.

$[L_a, L_b] \overset{\text{def}}{=} \{(n, i) : a \leq n \leq b\}$ is the part of the strip (a box) contained between layers $L_a$ and $L_b$, where $a < b$. We use the notation $[a, b]$ and the term interval $[a, b]$ for the box $[L_a, L_b]$ when the meaning of this notation is clear from the context.

$\mathcal{F}_{a,b}$ is the $\sigma$-algebra of events depending only on the environment in $[L_a, L_b]$.

e_v is a vector whose $v$th coordinate is 1 and all others are zeros.

I is a column vector with all components equal to 1.

If $x = (x(j))$ is a vector and $A = (a(i, j))$ a matrix we put

$$\|x\| \overset{\text{def}}{=} \max_j |x(j)| \quad \text{which implies} \quad \|A\| = \max_i \sum_j |a(i, j)|.$$
We say that $A$ is strictly positive (and write $A > 0$) if all its components satisfy $a(i, j) > 0$. $A$ is called non-negative (and we write $A \geq 0$) if all $a(i, j)$ are non-negative. A similar convention applies to vectors. We shall make use of the following easy fact:

$$\|A\| = \|A_{1}\|.$$ 

$\mathbb{X}$ denotes the set of non-negative unit vectors, $\mathbb{X} = \{x : x \in \mathbb{R}^m, x \geq 0, \|x\| = 1\}$.

$E_\mu(f)$, $\nu(g)$ denote the expectations of functions $f$ and $g$ over measure $\mu$ and $\nu$ respectively defined on the relevant probability spaces.

We often deal with $N^+$, $\ln N$, $\ln \ln N$, etc which are viewed as integer numbers. Strictly speaking, we should write $\lfloor N^+ \rfloor$, $\lfloor \ln \ln N \rfloor$, etc. However, our priority lies with the simpler notation and the exact meaning is always obvious from the context.

2. Review of related results from previous work

The purpose of this review is to list those results from [1, 12] which will be used in this work as well as to put the results of the present work into the right context. We note that many of the statements listed below were proved in [1, 12] under assumptions which are much milder than (1.5).

2.1. Auxiliary sequences of matrices

Let us fix $a \in \mathbb{Z}$ and define for $n \geq a$ two sequences of matrices: $\varphi_n$ and $\psi_n$. To this end put $\varphi_a \overset{\text{def}}{=} 0$ and let $\psi_a$ be a stochastic matrix. For $n > a$ matrices $\varphi_n$ and $\psi_n$ are defined recursively:

$$\varphi_n \overset{\text{def}}{=} (I - R_n - Q_n \varphi_{n-1})^{-1} P_n, \quad \psi_n \overset{\text{def}}{=} (I - R_n - Q_n \psi_{n-1})^{-1} P_n \quad (2.1)$$

Note that the existence of $(I - R_n - Q_n \psi_{n-1})^{-1}$ follows from (1.5).

Properties of matrices $\varphi_n$. We start with the probabilistic definition of $\varphi_n \equiv \varphi_n(a) = (\varphi_n(i, j))$ (which implies equation (2.1) for $\varphi_n$):

$$\varphi_n(i, j) = \mathbb{P}_\omega (\text{RW starting from } (n, i) \text{ hits } L_{n+1} \text{ at } (n+1, j) \text{ before visiting } L_a).$$

Obviously these probabilities are monotone functions of $a$ and hence the limits $\eta_n \overset{\text{def}}{=} \lim_n \varphi_n(a)$ exist for all (!) environments $\omega$. Lemma 4 in [1] implies that if (1.5) is satisfied then $\eta_n > 0$ for $\mathcal{P}$-almost every $\omega$ and for $n > a$

$$\varphi_n(i, j) > \varepsilon, \quad \psi_n(i, j) > \varepsilon \quad \text{for } \mathcal{P}\text{-almost every } \omega. \quad (2.2)$$

Definition of matrices $\zeta_n$. It is easy to see that since $\psi_n$ is stochastic, so are all the $\psi_n$, $n > a$ (lemma 2 in [1]). The following statement from [1] describes the $a \to -\infty$ limits of $\psi_n \equiv \psi_n(\psi_a)$ and defines a stationary sequence of stochastic matrices $\zeta_n$.

**Theorem 1.** Suppose that condition (1.5) is satisfied. Then

(a) For $\mathcal{P}$-a.e. sequence $\omega$ there exists $\zeta_n = \lim_{a \to -\infty} \psi_a(\psi_a)$, where the convergence is uniform in $\psi_a$ and the limit $\zeta_n$ does not depend on the choice of the sequence $\psi_a$.

(b) The sequence $\zeta_n = \zeta_n(\omega)$, $-\infty < n < \infty$, of $m \times m$ matrices is the unique sequence of stochastic matrices which satisfies the following system of equations:

$$\zeta_n = (I - Q_n \zeta_{n-1} - R_n)^{-1} P_n, \quad n \in \mathbb{Z}. \quad (2.3)$$

(c) The enlarged sequence $(P_n, Q_n, R_n, \zeta_n)$, $-\infty < n < \infty$, is stationary and ergodic.
Remark. Statements (a) and (b) imply that \( \zeta_n \equiv \zeta_n(\omega) \) depends only on the ‘past’ of the environment, namely on \( \omega_{\leq n} \equiv ((P_k, Q_k, R_k))_{k \leq n} \).

We need the following corollary of theorem 1 (remark 4 in [1]).

**Corollary 2.1.** Suppose that \( (P, Q, R) \) satisfies condition (1.5) (this can be any triple of matrices from the support of the distribution of \((P_0, Q_0, R_0)\)). Then there is a unique stochastic matrix \( \zeta \) such that

\[
\zeta = (I - Q \zeta - R)^{-1} P. \tag{2.4}
\]

**Proof.** Consider the environment with transition probabilities which do not change from layer to layer and are given by matrices \((P, Q, R)\), that is \( \omega = \{(P, Q, R)\} \). Then for this single environment all conditions of theorem 1 are satisfied. Now statement (a) implies that \( \zeta_n = \zeta_{n-1} \) and setting \( \zeta := \zeta_n = \zeta_{n-1} \) turns equation (2.3) into (2.4).

\( \square \)

**The non-arithmeticity condition.** We are now in a position to introduce the so called non-arithmeticity condition which will be often used in the sequel. Let \((P, Q, R)\) and \( \zeta \) be as in corollary 2.1. Set

\[
A_{(P, Q, R)} = (1 - Q \zeta - R)^{-1} Q \tag{2.5}
\]

and let \( \lambda(P, Q, R) \) be the leading eigenvalue of \( A_{(P, Q, R)} \). We say that the environment satisfies the non-arithmeticity condition if

the distribution of \( \lambda(P, Q, R) \) is non-arithmetic. \( \tag{2.6} \)

**Vectors \( \pi_n \).** Our sequence of stochastic matrices \( \zeta_n \) is such that \( \zeta_n(i, j) \geq \varepsilon > 0 \) for some \( \varepsilon > 0 \). Due to that we can always construct a sequence \( \pi_n \) of probability vectors such that \( \pi_n = \pi_{n-1} \zeta_{n-1} \). Namely, set \( \pi_{n,a} = \tilde{\pi}_a \zeta_a \cdots \zeta_{n-1} \), where \( \tilde{\pi}_a \) is a probability vector.

**Lemma 2.2.** If \( \zeta_n(i, j) \geq \varepsilon > 0 \) for some \( \varepsilon > 0 \) then the following limit exists and does not depend on the choice of the sequence of probability vectors \( \tilde{\pi}_a \):

\[
\pi_n \equiv \lim_{a \to -\infty} \tilde{\pi}_a \zeta_a \cdots \zeta_{n-1}. \tag{2.7}
\]

Moreover, for \( \theta = 1 - m \varepsilon \)

\[
\|\pi_n - \pi_{n,a}\| \leq \theta^{n-1-a} \quad \text{and} \quad \pi_n(i) > \varepsilon \quad \text{for any} \quad i \in [1, m]. \tag{2.8}
\]

**Remarks.**
1. In our case vectors \( \pi_n \equiv \pi(\omega_{\leq n}) \) form a stationary sequence.
2. Lemma 2.2 is a well known fact which follows from the usual contracting properties of products of stochastic matrices. We state it here for future references.

**Matrices \( A_n \) and Lyapunov exponents.** We can finally define the following sequence of matrices:

\[
A_n \equiv (I - Q_n \zeta_{n-1} - R_n)^{-1} Q_n. \tag{2.9}
\]

Obviously, \( A_n \) is a stationary sequence and the top Lyapunov exponent of the product of matrices \( A_n \) is defined as usual by

\[
\lambda \equiv \lim_{n \to \infty} \frac{1}{n} \log \|A_n A_{n-1} \cdots A_1\|. \tag{2.10}
\]

It is well known (see [9]) that with \( P \)-probability 1 the limit in (2.10) exists and does not depend on \( \omega \).
2.2. Recurrence and transience of RWRE

The recurrence criteria was proved in [1] for a RWRE on a strip in very general ergodic setting. We need the following particular case of this result.

**Theorem 2.** Suppose that conditions (1.4) and (1.5) are satisfied. Then the following statements hold for $\mathbb{P}_\omega$-almost surely:

(a) $\lambda < 0$ iff the RW is transient to the right: $\lim_{t \to \infty} X_t = \infty$,
(b) $\lambda > 0$ iff the RW is transient to the left: $\lim_{t \to \infty} X_t = -\infty$,
(c) $\lambda = 0$ iff the RW is recurrent: $\limsup_{t \to \infty} X_t = +\infty$ and $\liminf_{t \to \infty} X_t = -\infty$.

**Remark.** The proof of theorem 2 given in [1] contains the following useful statement: the RW is recurrent or transient to the right if and only if

$$\lim_{a \to -\infty} \phi_k = \zeta_k.$$  \hspace{1cm} (2.11)

2.3. Moment Lyapunov exponents $r(\alpha)$

From now on we consider RWRE which are transient to the right, that is $\lambda < 0$. Let us define a function whose properties are responsible for the speed of growth of our RW.

Let $A_n$ be a sequence of matrices defined by (2.9). For $\alpha \geq 0$ put

$$r(\alpha) \overset{\text{def}}{=} \lim_{n \to \infty} \sup (\mathbb{E} ||A_n \cdots A_1||^\alpha)\frac{1}{n}. \hspace{1cm} (2.12)$$

Note that if $m = 1$ then $\zeta_n = 1$, $A_n = q_n/p_n$, and $r(\alpha) = \mathbb{E}(q_0/p_0)^\alpha$. In this form $r(\alpha)$ was first introduced in [17].

**Lemma 2.3.** Suppose that (1.5) is satisfied. Then:

(a) the following limit exists and is finite for every $\alpha \geq 0$:

$$r(\alpha) = \lim_{n \to \infty} \mathbb{E} ||A_n \cdots A_1||^\alpha \frac{1}{n^\alpha}. \hspace{1cm} (2.13)$$

(b) the convergence in (2.13) is uniform in $\alpha \in [0, \alpha_0]$ for any $\alpha_0 > 0$

(c) $r'(0) = \lambda$.

2.4. Linear and sub-linear growth of the RW

Let as in the Introduction $X_t = (Z_t, Y_t)$ be a RW starting from a site $z \in L_0$, $\tilde{T}_n$ be the hitting time of layer $L_n$ by this walk.

**Theorem 3.** Suppose that (1.5) is satisfied and that $\lambda < 0$. Then:

(i) $r(1) < 1$ implies that for $\mathbb{P}$-a.e. environment $\omega$ with $\mathbb{P}_{\omega, z}$-probability 1

$$\lim_{n \to \infty} n^{-1}(\tilde{T}_n - \mathbb{E}_{\omega, z}\tilde{T}_n) = 0$$

and

$$\lim_{n \to \infty} n^{-1}\mathbb{E}_{\omega, z}\tilde{T}_n = c > 0. \hspace{1cm} (2.14)$$

(ii) $r(1) \geq 1$ implies that for $\mathbb{P}$-a.e. environment $\omega$ with $\mathbb{P}_{\omega, z}$-probability 1

$$\lim_{n \to \infty} n^{-1}\tilde{T}_n = \infty. \hspace{1cm} (2.15)$$

The fact that $\lim_{t \to \infty} t^{-1}Z_t = c^{-1}$ when $r(1) < 1$ and that $\lim_{t \to \infty} t^{-1}Z_t = 0$ when $r(1) \geq 1$ follows from (2.14) and (2.15), respectively.

These results extend the relevant statements from [31] to the case of the strip. Further details can be found in [12].
2.5. The diffusive regime (central limit theorem) for the RW

**Theorem 4.** Assume that (1.5) is satisfied, \( \lambda < 0 \), and \( r(2) < 1 \). Then there exists \( D_1 \) such that for \( P \)-almost every environment \( (\bar{T}_n - E_{\omega, z}\bar{T}_n)/\sqrt{n} \) converges weakly as \( n \to \infty \) to a normal distribution with zero mean and variance \( D_1 \).

**Remark.** It is easy to show that if in addition to the conditions of theorem 4 also (1.4) is satisfied, then there are constants \( c \) and \( D_2 \) such that \( (\bar{T}_n - c n)/\sqrt{n} \) converges weakly to a normal distribution with zero mean and variance \( D_2 \). Consequently \( (\bar{T}_n - c n)/\sqrt{n} \) converges weakly as \( n \to \infty \) in the annealed setting (that is with respect to \( P_z := P \times P_{\omega, z} \)) to a normal distribution with zero mean and variance \( D = D_1 + D_2 \).

3. Occupation times

As stated in the introduction, in this work the study the asymptotic behaviour of the RWRE is conducted in terms of that of the asymptotic behaviour of occupation times. In this section we derive formulae for the expectations of occupation times and discuss some of their properties. Denote the time spent by the walk at site \( x = (n, y) \) by \( \xi_x \). Obviously, the distribution of \( \xi_x \) depends on the starting point of the walk, say \( (k, i) \). Since \( \xi_x \) conditioned on the walk starting from \( x \) has a geometric distribution, it is easy to find the parameters of \( \xi_x \) also for the walk starting from arbitrary \( (k, i) \). Namely, we shall find \( F_{k,i} \) for all \( (k, i) \) in the strip including \( F_x \).

\[
F_{k,i} = \mathbb{E}_{\omega, (k, i)} (\xi_x) \quad \text{for all } (k, i)
\]

Then
\[
P_{\omega, (k, i)} \{ X \text{ reaches } x \} = P_{\omega, (k, i)} \{ \xi_x \geq 1 \} = F_{k,i} F_x^{-1}
\]
and hence
\[
P_{\omega, (k, i)} \{ \xi_x = 0 \} = 1 - F_{k,i} F_x^{-1}, \quad P_{\omega, (k, i)} \{ \xi_x = j \} = F_{k,i} F_x^{-2} (1 - F_x^{-1})^{j-1} \text{ if } j \geq 1.
\]

The expressions for \( F_{k,i} \) will be given in terms of the matrices defined in section 2. Denote by \( F_k \) the \( m \)-dimensional vector with components \( F_{k,i} \), \( 1 \leq i \leq m \) and let \( e_y \in \mathbb{R}^m \) be a vector whose \( y^{th} \) coordinate is 1 and all others are zeros.

**Lemma 3.1.** Suppose that (1.5) holds, \( x = (n, y) \). Then for \( P \)-almost all \( \omega \)

\[
F_k \equiv F_k(n, y) = \sum_{j=n}^{\infty} \zeta_k \ldots \zeta_{j-1} A_j \ldots A_{n+1} u_{n,y}, \quad \text{if } k < n
\]

\[
F_k \equiv F_k(n, y) = \sum_{j=k}^{\infty} \zeta_k \ldots \zeta_{j-1} A_j \ldots A_{n+1} u_{n,y}, \quad \text{if } k \geq n
\]

where
\[
u_{n,y} = (I - Q_n \zeta_{n-1} - R_n)^{-1} e_y.
\]

**Remark.** In the above formulae, we use the conventions that \( A_j \ldots A_{j-1} = I \) if \( j \leq n \) and \( \zeta_k \ldots \zeta_{j-1} = I \) if \( k \geq j \). Thus the first term in (3.1) is \( \zeta_k \ldots \zeta_{n-1} u_{n,y} \) and the first term in (3.2) is \( A_k \ldots A_{n+1} u_{n,y} \).

**Proof.** Consider a box \([L_a, L_b]\) with \( a < k, n < b \) and a walk \( \bar{X} \) on this box starting from \( (k, i) \) with absorption at layers \( L_a \) and \( L_b \). Denote by \( F_{k, i} \) the vector of conditional
expectations, \( \tilde{F}_{k,i} \overset{\text{def}}{=} \mathbb{E}_{\omega,(k,i)}(\xi^a,b) \), where \( \xi^a,b \) is the occupation time of \( x \) by the walk starting from \( (k, i) \). It is easy to see (first step analysis) that \( \tilde{F}_k \) satisfy the following system of equations:

\[
\tilde{F}_k = \Phi_k + \Phi_k \tilde{F}_{k+1} + R_k \tilde{F}_k + Q_k \tilde{F}_{k-1}
\]

if \( a < k < b \) and \( \tilde{F}_n = \tilde{F}_b = 0 \). \( \tag{3.4} \)

where \( \Phi_k = e_\alpha \) if \( k = n \) and \( \Phi_k = 0 \) otherwise. Systems of equations of this form were studied in [12]. The idea is to look for solutions to (3.4) of the form \( \tilde{F}_k = \varphi_k \tilde{F}_{k+1} + d_k \) with \( \varphi_k = 0 \) and \( d_k = 0 \). A simple calculation (see lemma 8 in [12]) shows that \( \varphi_k \) satisfy (2.1) and that

\[
\tilde{F}_k = \varphi_k \ldots \varphi_{n-1} d_n + \varphi_k \ldots \varphi_{d_{n+1}} + \ldots + \varphi_k \ldots \varphi_{d_{b-1}}.
\] \( \tag{3.5} \)

where

\[
d_r = \tilde{u}_r + \tilde{A}_r \tilde{u}_{r-1} + \tilde{A}_r \tilde{A}_{r-2} \tilde{u}_{r-2} + \ldots + \tilde{A}_r \ldots \tilde{A}_j \tilde{u}_{j-1} + \ldots
\]

and

\[
\tilde{u}_r = (I - Q_r \varphi_{r-1} - R_r)^{-1} \Phi_r, \quad \tilde{A}_r \overset{\text{def}}{=} (I - Q_r \varphi_{r-1} - R_r)^{-1} Q_r.
\]

In our case \( d_i = \tilde{A}_i \ldots \tilde{A}_{n+1} \tilde{u}_n \) if \( i > n \), \( d_n = \tilde{u}_n \) and \( d_0 = 0 \) otherwise which turns (3.5) into a version of (3.1), (3.2) with \( \zeta \)'s replaced by \( \varphi \)'s and the sums being finite. Note that:

(a) \( \lim_{a \to -\infty} \varphi_j = \zeta_j \) since \( \lambda < 0 \) (see (2.11) or [1]);
(b) moreover \( \varphi_j \not\overset{\text{pr}}{=} \zeta_j \) and therefore also \( \tilde{A}_j \not\overset{\text{pr}}{=} \tilde{A}_j \) as \( a \to -\infty \);
(c) thus \( \tilde{F}_k \) is monotonically increasing as \( a \) decreases or \( b \) increases with the terms of \( \tilde{F}_k \) converging to the corresponding terms of \( F_k \).

Series (3.1), (3.2) converge \( P \)-almost surely (once again due to \( \lambda < 0 \) and since they have positive entries, it follows that \( \tilde{F}_k \) converges \( P \)-almost surely to \( F_k \) as \( a \to -\infty \) and \( b \to +\infty \) and this proves the lemma.

One immediate corollary from (3.2) is the estimate of the probability of return from \( L_{n+t} \) to \( L_n \). For a \( \theta_0 < 1 \), set

\[
\Omega_{n,l,\theta_0} = \{ \omega : \mathbb{P}_\omega[X \text{ visits } L_n \text{ after } L_{n+t}] \geq \theta_0^t \}.
\] \( \tag{3.6} \)

**Lemma 3.2.** There are \( C > 0, \theta_0, \theta_1 \in (0, 1) \) such that

\[
P( \Omega_{n,l,\theta_0} ) < C \theta_1^t.
\] \( \tag{3.7} \)

**Proof.** It follows from (3.2) that \( ||F_{n+t}(n, y)|| \leq \text{Const} \sum_{j=n+t}^\infty ||A_j \ldots A_{n+1}|| \). Let \( \xi^{n,l} \) be the number of visits to \( L_n \) by the walk starting from \( L_{n+t} \). Then

\[
\mathbb{P}_\omega[X \text{ visits } L_n \text{ after } L_{n+t}] \leq \max_{z \in K_{n+l}} \mathbb{P}_{\omega,z}(\xi^{n,l} \geq 1) \leq \max_{z \in K_{n+l}} \mathbb{E}_{\omega,z}(\xi^{n,l})
\]

\[
\leq \max_y ||F_{n+t}(n, y)|| \leq \text{Const} \sum_{j=n+t}^\infty ||A_j \ldots A_{n+1}||.
\]

Fix any \( \alpha \) such that \( r(\alpha) < 1 \). Then for any \( \theta_0 > 0 \)

\[
P( \sum_{j=n+t}^\infty ||A_j \ldots A_{n+1}|| \geq \theta_0^t ) \leq \theta_0^{-\alpha} \sum_{j=n+t}^\infty E(||A_j \ldots A_{n+1}||)^\alpha \leq \text{Const} \theta_0^{-\alpha} r(\alpha)^t.
\]

and it remains to choose \( \theta_0 \) so that \( \theta_1 := r(\alpha)\theta_0^{-\alpha} < 1 \).

Let us now discuss a corollary which follows from (3.1). This formula will quite often be used when \( n - k > eN^2 \) for some \( e > 0 \) and \( N \to \infty \) in which case the expression for \( F_k(n, y) \) can be simplified. Namely, by lemma 2.2

\[
\zeta_k \ldots \zeta_{j-1} = \pi_j + \mathcal{O}(\theta^{j-k}),
\] \( \tag{3.8} \)
where $\pi_j$ is a rank one matrix all rows of which are equal to $\pi_j$. We can now rewrite (3.1) as

$$F_k(n, y) = \sum_{j=n}^{\infty} \pi_j A_j \ldots A_{n+1} u_n = (1 + O(\theta^{n-k})) \sum_{j=n}^{\infty} \pi_j A_j \ldots A_{n+1} u_{n,y}. \quad (3.9)$$

Fix $\epsilon > 0$ and $\theta \in (0, 1)$. Then it is easy to see that for sufficiently large $N$

$$P \left\{ \theta^N \max_{n \in [0, N]} \sum_{j=n}^{\infty} ||A_j \ldots A_{n+1}|| > N^{-50(1+\epsilon)} \right\} < N^{-100}.$$

Since errors of such order as well as events of such small probability will not play any role for our results and in our proofs (see the statements of the main theorems below), the dependence of $F_k(n, y)$ on the starting point of the walk can be neglected and this vector can be replaced by a single number

$$\rho_{n,y} := \sum_{j=n}^{\infty} \pi_j A_j \ldots A_{n+1} u_n. \quad (3.10)$$

We also set

$$\rho_n := \sum_y \rho_{n,y} = \sum_{j=n}^{\infty} \pi_j A_j \ldots A_{n+1} u_n \quad (3.11)$$

where $u_n = (I - Q_n \zeta_n - R_n)^{-1}$. Obviously, $\rho_n$ is the expectation of the time spent by the walk in layer $n$.

Since $A_j$ is a stationary sequence of positive matrices there exists a stationary sequence of vectors $v_j > 0$, $v_j \in X$ and numbers $\lambda_j > 0$ such that

$$A_j v_j = \lambda_j v_j \quad (3.12)$$

and there is a sequence of functionals $l_j$ and a $\theta < 1$ such that for any vector $u$

$$A_j \ldots A_{n+1} u = l_n(u) \lambda_j \ldots \lambda_{n+1} (1 + O(\theta^{n-j})) v_j.$$

Moreover $l_n$ is well approximated by local functions, that is, there exists $F_{n,n,r}$ measurable functions $l_{n,r}$ such that $||l_n - l_{n,r}|| \leq \theta^r$ whenever $r > r_0$ ($r_0$ depends only on $\theta$ and $\epsilon$ from condition (1.5)).

The foregoing discussion allows us to write

$$\rho_n = l_n(u_n) w_n + R_n \quad (3.13)$$

where

$$w_n = \sum_{j \geq n} \lambda_j \ldots \lambda_{n+1} (\pi_j, v_j) \quad (3.14)$$

and $R_n$ denotes the contribution of subleading terms. Denote

$$w_{n1,n2} = \sum_{j=n1}^{n2} \lambda_j \ldots \lambda_{n+1} (\pi_j, v_j).$$

The following lemma characterizes the tail behaviour of the distribution of $w_n$ and thus also of $\rho_n(y)$ and $\rho_n$. It adjusts to our needs some well known results from [15]; the latter played a major role in many previous studies of the asymptotic behaviour of the RWRE, in particular in [6, 17]. The derivation of this lemma will be given in section 7 using the results of appendix C. The lemma relies on the assumption that

there is $s > 0$ such that $r(s) = 1.$ \quad (3.15)
Note that $\ln(r(\cdot))$ is a strictly convex function (see e.g. section C.3) and therefore the existence of the solution $s$ to (3.15) implies its uniqueness. On the other hand if (3.15) has no positive solutions then $r(\alpha) < 1$ for all $\alpha > 0$. In particular, the walk is diffusive in that case in view of theorem 4.

For the rest of the paper we suppose that (3.15) is satisfied. Further analysis will heavily rely on the asymptotic properties of the of tails of distributions of $\rho_n$ which will follow from those of $w_n$. The latter are described by the following lemma.

**Lemma 3.3.** Suppose that (1.5) is satisfied. Then there are constants $\bar{C}$ and $\bar{s} > s$ such that

(a) If $n_2 - n_1 > \bar{C} \ln t$ then

$$P(w_{n_1} - w_{n_2} \geq 1 | \xi_{n_1} = \xi, \pi_{n_1} = \pi, v_{n_1} = v) \leq Ct^{-\bar{s}}.$$  

(b) If $n_2 - n_1 > \bar{C} \ln t$ and the non-arithmeticity condition (2.6) holds (in addition to (1.5)) then there is a function $f(\xi, \pi, v) > 0$ such that

$$P(w_{n_1, n_2} \geq l | \xi_{n_1} = \xi, \pi_{n_1} = \pi, v_{n_1} = v) \sim f(\xi, \pi, v)t^{-s}.$$  

(Here and below ‘$\sim$’ means the bound which is uniform in the parameter involved. That is, given $\varepsilon$ there exists $t_0$ such that if $t > t_0$ and $n_2 - n_1 > \bar{C} \ln t$ then

$$|t^s P(w_{n_1, n_2} \geq l | \xi_{n_1} = \xi, \pi_{n_1} = \pi, v_{n_1} = v) - f(\xi, \pi, v)| < \varepsilon.$$  

In particular

$$P(w_n \geq l | \xi_n = \xi, \pi_n = \pi, v_n = v) \sim f(\xi, \pi, v)t^{-s}.$$  

(c) There exists $\bar{s} > s$ and $\bar{C} > 0$ such that

$$P(R_n > l | \xi_n = \xi, \pi_n = \pi, v_n = v) \leq Ct^{-\bar{s}}.$$  

The proof of lemma 3.3 is given in section 7.

4. Main results

The description of the asymptotic behaviour of $T_N$ (defined by (1.10)) will be derived from the asymptotic properties of traps. Our first main result describes these properties. Let us introduce the exact definition of a trap on $[0,N]$.

**Definition.** Let $M = M_N := \ln \ln N$ and $\delta > 0$ be a given (small) number, $w_n$ is defined by (3.14). We say that $n$ is a massive site if $w_n \geq \delta N^{1/s}$. A site $n \in [0, N - 1]$ is marked if it is massive and $w_{n+j} < \delta N^{1/s}$ for $1 \leq j \leq M$. For $n$ marked the interval $[n - M, n]$ is called the trap (or $\delta N^{1/s}$-trap) associated with $n$. We call the number $m_n = \sum_{j=n-M}^{n} \rho_j$ the mass of the trap.

Note that this definition implies that distinct traps are disjoint.

The asymptotic distribution of traps is described by the following

**Theorem 5.** Assume that the non-arithmeticity condition (2.6) holds. Then there exists a constant $c$ such that the point process

$$\left\{ \left( \frac{n_j}{N}, \frac{m_{n_j}}{N^{1/s}} \right) : n_j \text{ is } \delta N^{1/s} \text{-marked and } 0 \leq n_j \leq N \right\}$$  

converges as $N \to \infty$ to a Poisson process on $[0, 1] \times [\delta, \infty)$ with the measure $c \delta t^s \mu_\delta(d\delta)$, where $\mu_\delta$ converges to a measure with density $\delta^{s-1}$ as $\delta \to 0$. 

Remarks.

1. Each component of the point process (4.1) is itself a point process converging to a Poisson process.

2. Theorem 5 extends certain statements from [6] (see lemma 4.4 there) to the case of the walk on a strip. It may be worth mentioning that in [6] we used the term *cluster* for what we have now decided to call a *trap*; the latter term seems to better reflect the main properties of this object.

3. The measure $\mu_\delta$ will be described in more explicit terms later. However, its explicit description depends on the choice of the definition of a trap and is only important because it helps to find the limit of $\mu_\delta$ as $\delta \to 0$.

The above theorem plays a major role in the description of the asymptotic behaviour of the walk in the subdiffusive regime $s \in (0, 2)$. To state our second main result, we define $t_N$ which is a normalized version of $T_N$, namely we set

$$t_N = \begin{cases} 
\frac{T_N}{N^{1/s}} & \text{if } 0 < s < 1, \\
\frac{T_N - \mathbb{E}_w(T_N)}{N^{1/s}} & \text{if } 1 \leq s < 2.
\end{cases}$$

The definition of $t_N$ implies that complete understanding of its asymptotic properties should include the study of those of $\mathbb{E}_w(T_N)$. The corresponding normalized quantity is defined as follows:

$$u_N = \begin{cases} 
\frac{\mathbb{E}_w(T_N)}{N^{1/s}} & \text{if } 0 < s < 1, \\
\frac{\mathbb{E}_w(T_N) - u_N}{N} & \text{if } s = 1, \\
\frac{\mathbb{E}_w(T_N) - \mathbb{E}(T_N)}{N^{1/s}} & \text{if } 1 < s < 2,
\end{cases}$$

where $u_N = N \mathbb{E}(\rho_n I_{\rho_n < x_N})$ with $x_N$ defined by $P\{\rho_n > x_N\} = N^{-1}$.

To state our next theorem, we set

$$\Theta^{(N, \delta)} = \left\{ \Theta^{(N, \delta)}_j : 0 \leq n_j \leq N \text{ and } n_j \text{ is } \delta N^{1/s}-marked \right\}$$

**Theorem 6.** Assume that the non-arithmeticity condition (2.6) holds. Then for $0 < s < 2$ and $\delta > 0$ there is a sequence $\Omega_{N, \delta} \subset \Omega$ such that $\lim_{N \to \infty} P(\Omega_{N, \delta}) = 1$ and a sequence of random point processes

$$(\Theta^{(N, \delta)}, \Gamma^{(N, \delta)}) = \left( \left\{ \Theta^{(N, \delta)}_j, \Gamma^{(N, \delta)}_j \right\} \right)$$

such that

(i) The component $\Gamma^{(N, \delta)} = \{\Gamma^{(N, \delta)}_j\}$ is a collection of asymptotically i.i.d. random variables with mean 1 exponential distribution.

(ii) The $t_N$ and $u_N$ can be presented in the following form:

(a) If $0 < s < 1$ then for $\omega \in \Omega_{N, \delta}$

$$t_N = \sum_j \Theta^{(N, \delta)}_j \Gamma^{(N, \delta)}_j + R_N,$$

where $R_N \geq 0$ and $\mathbb{E}(1_{\Omega_{N, \delta}} R_N) = \mathcal{O}(\delta^{1-s})$,

$$u_N = \sum_j \Theta^{(N, \delta)}_j + \hat{R}_N,$$

where $\hat{R}_N \geq 0$, $\mathbb{E}(\hat{R}_N) = \mathcal{O}(\delta^{1-s})$. 

(4.5)
(b) If $s = 1$ then for $\omega \in \Omega_{N,\delta}$ and a given $1/2 < \kappa < 1$
\[ t_N = \sum_j \Theta_j^{(N,\delta)}(\Gamma_j^{(N,\delta)} - 1) + R_N, \quad \text{where } E\left[1_{\Omega_{N,\delta}}(R_N^2)\right] = O(\delta^{2\kappa-1}), \]
\[ u_N = \sum_j \Theta_j^{(N,\delta)} - \bar{c} \ln \delta + \tilde{R}_N, \quad \text{where } E(\tilde{R}_N^2) = O(\delta). \]

(c) If $1 < s < 2$ then for $\omega \in \Omega_{N,\delta}$
\[ t_N = \sum_j \Theta_j^{(N,\delta)}(\Gamma_j^{(N,\delta)} - 1) + R_N, \quad \text{where } E\left[1_{\Omega_{N,\delta}}(R_N^2)\right] = O(\delta^{2-\kappa}), \]
\[ u_N = \sum_j \Theta_j^{(N,\delta)} - \frac{\bar{c}}{(s-1)\delta^{s-1}} + \tilde{R}_N, \quad \text{where } E(\tilde{R}_N^2) = O(\delta^{2-\kappa}). \]

Remarks.

1. Theorem 6 was proven in [6] for SRWRE. The next two remarks are similar to those following theorem 2 in [6]; we nevertheless believe that they are worth of being repeated.
2. The estimates of the remainders in the statements of theorem 6 hold for all $\delta > 0$ but are not uniform in $N$. More precisely, e.g. the relation $E(|\tilde{R}_N^2|) = O(\delta)$ in (b) means that for any $\delta > 0$ there is $N_\delta$ and a constant $C$ (which does not depend on $\delta$) such that
   \[ E(|\tilde{R}_N^2|) \leq C\delta \text{ if } N > N_\delta. \]
3. The dependence of $\Theta^{(N,\delta)}$ on $\omega$ persists as $N \to \infty$ whereas $\Gamma^{(N,\delta)}$ becomes ‘almost’ independent of $\omega$. More precisely, for $K \gg 1$ and sufficiently large $N$ the events $B_k^N := \{|\Theta^{(N,\delta)}| = k\}, 0 \leq k \leq K$, form, up to a set of a small probability, a partition of $\Omega$. If $\omega \in B_k^N$ then $\Gamma^{(N,\delta)}(\omega, X)$ is a collection of $k$ random variables which converge weakly as $N \to \infty$ to a collection of $k$ i.i.d. standard exponential random variables. Thus the only dependence of $\Gamma^{(N,\delta)}(\omega, X)$ on $\omega$ and $\delta$ which persists as $N \to \infty$ is reflected by the fact that $|\Theta^{(N,\delta)}| = |\Gamma^{(N,\delta)}|$. (Remember that $X$ is the trajectory of the walk and the purpose of our notation is to emphasize the dependence of $\Gamma^{(N,\delta)}$ on both $\omega$ and $X$.)

The following statement is a corollary of theorem 6. In the case of the SRWRE, theorem 7 is one of the main results of [17].

**Theorem 7.** The annealed walk has the following properties:

(a) If $s < 1$ then the distribution of $T_N/N^{1/s}$ converges to a stable law with index $s$.
(b) If $1 < s < 2$ then there is a constant $\kappa$ such that the distribution of $(T_N - Nu)/N^{1/s}$
   converges to a stable law with index $s$.
(c) If $s = 1$ then there is a sequence $u_N \sim cN \ln N$ (defined as in (4.3)) such that the distribution of $(T_N - u_N)/N$ converges to a stable law with index 1.

The proof of this theorem will not be given because its derivation from theorems 5 and 6 is easy (see lemma B.2 in appendix B) and also was carried out in [6].

5. The asymptotic properties of traps and roof of theorem 5

5.1. Auxiliary lemmas

The proof of theorem 5 requires understanding of the asymptotic behaviour of traps. The following five lemmas describe the properties of traps that shall be used in the sequel.

We start with $w_n$ defined by (3.14). Observe that we have
\[ w_n = \lambda_n w_{n+1} + (\pi_n, v_n). \]
Lemma 5.1. There exist $\epsilon_1 > 0$, $\epsilon_2 > 0$, $0 < \beta < 1$ and $C > 0$ such that for any $(\tilde{\zeta}, \tilde{\pi}, \tilde{v})$ we have

(a) If $k \leq \epsilon_1 \ln t$ then

$$P(w_n \geq t, w_{n+k} \geq t | \zeta_n = \tilde{\zeta}, \pi_n = \tilde{\pi}, v_n = \tilde{v}) \leq \frac{C^k}{t^s};$$

(b) If $k \geq \epsilon_1 \ln t$ then

$$P(w_n \geq t, w_{n+k} \geq t | \zeta_n = \tilde{\zeta}, \pi_n = \tilde{\pi}, v_n = \tilde{v}) \leq Ct^{-s(\epsilon_2 + 1)}.$$

Proof. For brevity we shall denote $\tilde{P} = P(\cdot | \zeta_n = \tilde{\zeta}, \pi_n = \tilde{\pi}, v_n = \tilde{v})$.

(a) From (5.1) we have

$$w_n = \lambda_n \ldots \lambda_{n+k-1} + O(K^k).$$

So if $K^k \ll t$ then

$$\tilde{P}(w_n \geq t, w_{n+k} \geq t) \leq \tilde{E}(1_{\lambda_n \ldots \lambda_{n+k-1} \geq 1/2} P_{\xi_{n+k}, \pi_{n+k}, v_{n+k}} (w_{n+k} \geq t))$$

$$\leq \tilde{E}(1_{\lambda_n \ldots \lambda_{n+k-1} \geq 1/2} f(\zeta_n, \pi_n, v_n))$$

$$\leq C \tilde{P}(\lambda_n \ldots \lambda_{n+k-1} \geq 1/2) \frac{C^k}{t^s} \leq C \frac{C^k}{t^s}.$$

(b) Consider two cases

(I) $k > \tilde{c} \ln t$ where $\tilde{c}$ is the constant from lemma 3.3. Then

$$\tilde{P}(w_n \geq t, w_{n+k} \geq t) \leq \tilde{P}(w_{n,n-k-1} \geq t - 1, w_{n+k} \geq t) + \tilde{P}(w_n - w_{n,n-k-1} > 1).$$

The second term is $O(t^{-s})$ while the first term equals to

$$\tilde{E}(1_{\lambda_n \ldots \lambda_{n+k-1} \geq 1/2} P_{\xi_{n+k}, \pi_{n+k}, v_{n+k}} (w_{n+k} \geq t))$$

$$\leq C \tilde{P}(w_{n,n-k-1} \geq t - 1) \leq C \frac{C}{t^s}.$$ 

(II) $\epsilon_1 \ln t < k \leq \tilde{c} \ln t$. Fix $\tilde{c} \ll 1$. Then

$$\tilde{P}(w_n \geq t, w_{n+k} \geq t) \leq \tilde{P}(w_{n+k} \geq t^{1/s}) + \tilde{P}(w_n \geq t, t \leq w_{n+k} \leq t^{1/s}).$$

The first term is $O(t^{-s(1+1/s)})$ while the second term is less than

$$\tilde{P}(\lambda_n \ldots \lambda_{n+k-1} \geq t^{-s/2}, w_{n+k} \geq t) + \tilde{c} \ln t \max_{1 \leq i < j < k} \tilde{P}(\lambda_n \ldots \lambda_{n+j} \geq t^{-1/2}, w_{n+k} \geq t).$$

Both terms are estimated in the same way so we only discuss the first one

$$\tilde{P}(\lambda_n \ldots \lambda_{n+k-1} \geq t^{-1/2}, w_{n+k} \geq t) \leq \tilde{E}(1_{\lambda_n \ldots \lambda_{n+k-1} \geq t^{-1/2}} P_{\xi_{n+k}, \pi_{n+k}, v_{n+k}} (w_{n+k} \geq t))$$

$$\leq C \frac{C^k}{t^s} \tilde{P}(\lambda_n \ldots \lambda_{n+k-1} \geq t^{s/2}) \leq Ct^{s(1+1/s)}$$

as claimed.

It may happen that not all massive sites belong to one of the clusters. This situation is controlled by the following

Lemma 5.2. There is $\beta < 1$ such that for $n \in [0, N-1]$

$$P(\rho_n \geq \delta N^{1/s} \text{ and } n \text{ is not in a trap}) \leq \text{Const} \frac{\beta^M}{N}.$$ (5.2)
Proof. Suppose that \( n \) is a massive point which is not in a trap. Then consider all massive points \( n_i \) such that \( n < n_1 < \ldots < n_k < n + M \). Note that such points exist because otherwise \( n \) would have been a marked point. Let now \( n^* > n_k \) be the nearest to \( n_k \) massive point. Then by construction \( n^* \geq n + M \). Also \( n^* \leq n + 2M \) because otherwise \( n_k \) would have been a marked point and \( n \) would belong to the \( n_k \)-trap. Hence the event \[ \{ n \text{ is massive and not in a trap} \} \subseteq \bigcup_{n' \in \left[ n+M, n+2M \right]} \{ w_n \geq \delta N^{1/s}, w_{n'} \geq \delta N^{1/s} \} . \]

By lemma 5.1(b) we obtain \[ P(n \text{ is massive and not in a trap}) \leq \sum_{n' = n+M}^{n+2M} P(w_n \geq \delta N^{1/s}, w_{n'} \geq \delta N^{1/s}) \leq \text{const} \frac{\beta M}{N} \] which proves our statement.

Our next goal is to show that \( P(n \text{ is massive}) \) and \( P(n \text{ is marked}) \) are of the same order.

Lemma 5.3.
(a) For each \( l, R \geq 1 \) the following limit exists \[ f_l(\bar{\zeta}, \bar{\pi}, \bar{v}, R) = \lim_{t \to \infty} tP(w_n \geq tR, w_{n+j} < t \text{ for } j = 1, \ldots, l| \zeta_n = \bar{\zeta}, \pi_n = \bar{\pi}, v_n = \bar{v}). \]

(b) Let \( \bar{f} = \lim_{l \to \infty} f_l(\bar{\zeta}, \bar{\pi}, \bar{v}, R) \). Then \( | \bar{f}(\bar{\zeta}, \bar{\pi}, \bar{v}, R) - f_l(\bar{\zeta}, \bar{\pi}, \bar{v}, R) | \) is of the same order.

Proof. By (5.1) \[ \bar{P}(w_n \geq tR, w_{n+j} < t \text{ for } j = 1, \ldots, l) \sim \mathcal{E} \left( P_{\zeta_n, \pi_n, v_n} \left( \frac{tR}{\lambda_n \ldots \lambda_{n+l-1}}, \frac{t}{\max_j(\lambda_n \ldots \lambda_{n+l-1})} \right) \right) \] so the result follows from lemma 3.3.

(b) Since \[ \bar{P}(w_n \geq tR, w_{n+j} < t \text{ for } j = 1, \ldots, l) - \bar{P}(w_n \geq tR, w_{n+j} < t \text{ for } j = 1, \ldots, l+1) \leq \bar{P}(w_n \geq t, w_{n+l+1} \geq t) \] the result follows by lemma 5.1(a). \( \square \)

Lemma 5.4. \( P(\bar{f}(\zeta, \pi, v, 1) > 0) > 0 \).

Proof. Assume to the contrary that \( P(\bar{f} > 0) = 0 \). Then for each \( \epsilon \) there is \( n_0 \) such that for \( N \geq N_0 \) we have \( P(n \text{ is marked}) \leq \epsilon/N \). Combining lemma 3.3 and lemma 5.2 we obtain that there is a constant \( c \) such that \[ P(\mathbb{T}_n) \geq \frac{c}{N}, \quad \text{where } \mathbb{T}_n = \{ n \text{ is marked and belongs to a trap} \}. \]

If \( n \) is in a trap let \( D_n \) be the distance to the nearest marked point to the left of \( n \). Given \( D \) we write \[ P(\mathbb{T}_n) = P(\mathbb{T}_n \text{ and } D_n < D) + P(\mathbb{T}_n \text{ and } D_n \geq D) \]. The first term equals to \[ \sum_{j=0}^{D-1} P(n \text{ is massive and } n+j \text{ is marked}) \leq \sum_{j=0}^{D-1} P(n+j \text{ is marked}) \leq \frac{\epsilon D}{N} . \]
On the other hand,
\[ P(T_n \text{ and } D_n \geq D) \leq \sum_{j=D}^{M} P(n \text{ and } n + j \text{ are massive}) \leq \frac{C\beta^D}{N}. \]
Thus,
\[ P(T_n) \leq \varepsilon D + \frac{C\beta^D}{N}. \]
Taking \( D \) so large that \( C\beta^D \leq c/3 \) and then choosing \( \varepsilon \leq c/3D \) we obtain a contradiction with (5.3).

We now turn our attention to the mass of the trap. Observe that for \( j \leq M \) we have
\[ \rho_{n-j} \sim w_n \lambda_{n-j} \ldots \lambda_n \gamma_n(v_n-j). \]
Accordingly introduce
\[ a_n = \sum_{j=0}^{\infty} \lambda_{n-j} \ldots \lambda_n \gamma_n(v_n-j) \quad \text{and} \quad a_{n_1,n_2} = \sum_{j=0}^{n_1-n_2} \lambda_{n_2-j} \ldots \lambda_n \gamma_n(v_n-j). \]

In the next result proven in section 7 we use the same notation as in Lemma 3.3.

**Lemma 5.5.**

(a) If \( k > \bar{C} \ln t \) then
\[ P(a_n - a_{n-k,n} \geq 1) \leq Ct^{-\delta} \]
(b) \( P(a_{n-k,n} \geq t) \sim \hat{c} t^{-\delta}. \)
(c) Moreover there exists a measure \( \nu \) such that if \( k > \bar{C} \ln t \) then
\[ t^s E \left( 1_{a_n \geq t} f \left( \zeta_n, \pi_n, v_n, \frac{a_n}{t} \right) \right) \sim \hat{c} \int \int f(\zeta, \pi, v, z) \frac{dz \nu(\zeta, \pi, v)}{z^s}. \]

**Corollary 5.6.**

(a) The following limit exists
\[ h(t, \delta) = \lim_{N \to \infty} N^{1/s} P(n \text{ is marked and } m_n > t N^{1/s}). \]
(b) There is \( c > 0 \) such that \( \lim_{\delta \to 0} h(t, \delta) = ct^{-\delta}. \)

**Proof.** Take \( W \) such that if \( n \) is marked then \( w_n \in [\delta N^{1/s}, \delta W N^{1/s}] \). We have
\[ P(n \text{ is marked and } m_n > t N^{1/s}) = E \left( 1_{a_n \geq t} \frac{1}{\delta N} P_{\delta \pi_n, v_n} \left( w_n > \min \left( 1, \frac{t}{\delta \pi_n} \right) \delta N^{1/s} \right) \right) \sim \frac{1}{\delta N} E \left( 1_{a_n \geq t} \frac{1}{\delta N} f \left( \zeta_n, \pi_n, v_n, \min \left( 1, \frac{\delta a_n}{t} \right) \right) \right). \]
This proves (a). To prove (b) we use Lemma 5.5(c) to get
\[ \frac{1}{\delta N} E \left( 1_{a_n \geq t} \frac{1}{\delta N} f \left( \zeta_n, \pi_n, v_n, \min \left( 1, \frac{\delta a_n}{t} \right) \right) \right) \sim \hat{c} \int \int \frac{dz \nu(\zeta, \pi, v)}{z^s} \frac{dz}{z^{s+t}}. \]
5.2. Proof of theorem 5

We are now in a position to prove theorem 5. The following lemma essentially repeats the statement of theorem 5 with the difference that we can now state it in terms of $h(\cdot, \delta)$ studied above (corollary 5.6).

**Lemma 5.7.** Suppose that all conditions of theorem 5 are satisfied. Then

(a) For a fixed $\delta > 0$ the point process

$$ \{(nN^{-1}, m_nN^{-\frac{1}{2}}) : n \text{ is } \delta N^{-\frac{1}{2}} \text{-marked}\} $$

converges as $N \to \infty$ to a Poisson process on $[0, 1] \times [\delta, \infty)$ with measure $dt' \mu_{\delta}$ such that $\mu_{\delta}([t_1, t_2]) = h(t_2, \delta) - h(t_1, \delta)$.

(b) As $\delta \to 0$ $\mu_{\delta}$ converges to a measure with density $c_{\delta} t_1^{\delta - 1}$.

**Proof.** To prove (a) we use Bernstein’s big block-small block method. Namely, we divide $[0, N]$ into big blocks $I_j$ of length $N^{\frac{1}{3}}$ separated by small blocks of length $N^{\frac{1}{3}}/2$. We take $\varepsilon_2 < \varepsilon_2$, where $\varepsilon_2$ is the constant from lemma 5.1. By lemma 3.3 the probability that there is a trap inside the union of the short intervals tends to 0 so it suffices to consider the union of long intervals. We claim that for each $j$ the probability that $I_j$ contains two or more traps is $o(N^{\frac{1}{3}})$. Indeed due to lemma 5.1, the above probability is bounded by

$$ \sum_{n \in I_j} \sum_{k < N^{\frac{1}{3}}} P(\text{both } n \text{ and } n+k \text{ are marked}) $$

$$ \leq \sum_{n \in I_j} \sum_{k < N^{\frac{1}{3}}} P(\text{both } n \text{ and } n+k \text{ are massive}) $$

$$ \leq C \sum_{n \in I_j} \beta^{\ln N} \sum_{k = \varepsilon_2}^{N^{\frac{1}{3}}} N^{-\varepsilon_2} $$

proving the claim.

Next, let $\{J_i\}$ be a collection of $d$ non-intersecting open intervals in $[\delta, \infty)$ and $K = (t'_1, t'_2)$ be an open interval from $[0, 1]$.

Let

$$ N_i = \text{Card}(n \leq N : n \text{ is marked and } (n, m_n) \in NK \times N^{1/3} J_i). $$

We need to show that, as $N \to \infty$, $\{N_i\}$ converge to independent random variables having Poisson distribution with parameter $(t'_2 - t'_1)\mu_{\delta}(J_i)$ (the proof for all other finite collections of open quadrangles in $[0, 1] \times [\delta, \infty)$ easily follows from this case). We shall now replace $m_n$ by $\tilde{m}_n$ which are defined similarly to $m_n$ but have the property of being i.i.d. random variables and, at the same time, $|m_n - \tilde{m}_n| \leq N^{-\frac{2}{3}}$. Namely, we do the following:

1. Define Markov process $(\psi_n, \hat{\pi}_n)$ starting with initial conditions $\hat{\pi}_n$ and $\psi_n$ which have all entries equal to $1/m$, where $n_j$ is the middle of the short interval preceding $I_j$. This process is defined for $n > n_j$ with the $\psi$ component given by (2.1) and $\hat{\pi}_n = \hat{\pi}_{n-1} \psi_{n-1}$.

2. Similarly to (3.11), set for $n \in I_j$

$$ \hat{\rho}_n := \sum_{j=n}^{n_{j+1}} \hat{\pi}_j \hat{A}_j \ldots \hat{A}_{n+1} \tilde{u}_n, $$

where $\hat{A}_j$, $\tilde{u}_n$ are defined as the corresponding $A$s and $\tilde{u}s$ with $\xi$s replaced by $\psi$s.

3. Finally $\tilde{m}_n$ is defined similarly to $m_n$ but with $\hat{\rho}_n$ replacing $\rho_n$. 

Subdiffusive random walks on a strip

The independence mentioned above is obvious from the construction and the approximation property now follows from proposition D.1:

\[ ||\xi_n - \psi_n|| \leq \text{Const} \theta^{n-s\gamma}. \]

Let now \( \{\Gamma_j\} \) be a sequence of random \( d \)-dimensional vectors such that

\[ \Gamma_j = \begin{cases} e_l & \text{if } I_j \subset NK, \text{ there is exactly one trap } \tilde{n} \in I_j, m_{\tilde{n}} \in N^{1/2}J_l, \\ 0 & \text{in all other cases.} \end{cases} \]

Then \( \Gamma_j \) are iid random vectors and

\[ P(\Gamma_j = e_l) \sim (t'_{2} - t'_{1})\mu_3(J_l)N^{\gamma-1} + o(N^{\gamma-1}). \]

Therefore part (a) follows from the Poisson limit theorem for independent random vectors. Part (b) follows from part (a) and corollary 5.6(b).

6. Proof of theorem 6

The proof of theorem 6 is to an extent similar to that of theorem 2 in [6] and we shall give only an outline of it putting emphases at those parts which are new. As in [6], we start with a lemma which allows us to show the smallness of the contribution to \( T_N \) which comes from the sites where \( \rho_n < \delta N^{1/2} \). We then compute the main contribution to \( T_N \) which comes from the traps in \([0, N]\) described in theorem 5.

Within this section we shall use the following notation: \( \xi_n = \sum_{i=1}^{m} \xi_{n,i} \). Obviously, \( \rho_n = E_\omega(\xi_n) \). In these notations, lemma 6.1 becomes an exact copy of lemma 4.1 from [6].

**Lemma 6.1.** Let \( \delta > 0 \). Then there is \( N_\delta \) such that for \( N > N_\delta \) the following holds:

(a) If \( 0 < s < 1 \) then

\[ E \left( \sum_{w_n < \delta N^{1/2}} \xi_n \right) \leq \text{Const} N^{1/s} \delta^{1-s}. \]

(b) If \( 1 < s < 2 \) then there is a set \( \tilde{\Omega}_{N,\delta} \) such that \( P(\tilde{\Omega}_{N,\delta}) \leq N^{-100} \) and

\[ E \left( \left( \sum_{w_n < \delta N^{1/2}} (\xi_n - \rho_n) \right)^2 \right) \leq \text{Const} N^{2/s} \delta^{2-s}. \]

(c) If \( 0 < s < 1 \) then

\[ E \left( \sum_{w_n < \delta N^{1/2}} \rho_n \right) \leq \text{Const} N^{1/s} \delta^{1-s}. \]

(d) If \( 1 < s < 2 \) then

\[ \text{Var} \left( \sum_{w_n < \delta N^{1/2}} \rho_n \right) \leq \text{Const} N^{2/s} \delta^{2-s}. \]
(e) If $s = 1$ then given $\frac{1}{2} < \kappa < 1$ there is a set $\Omega_{N, \delta}$ such that $P(\Omega_{N, \delta}^c) \leq N^{-100}$ and

$$E \left( I_{\Omega_{N, \delta}} \left( \operatorname{Var}_\omega \left( \sum_{\rho_n < \delta N} (\xi_n - \rho_n) \right) \right)^\kappa \right) \leq \operatorname{Const} N^{2\kappa} g^{2^{\kappa-1}}, \quad (6.1)$$

$$E \left( \sum_{\rho_n < \delta N} \left( \rho_n - E(\rho I_{\rho_n < \delta N}) \right)^2 \right) \leq \operatorname{Const} N^{2\delta}. \quad (6.2)$$

Proof. Parts (a) and (c) follow from lemma 3.3 and Markov inequality (see the proof of lemma 4.1 in [6]).

The proofs of (b), (d), and (e) in lemma 4.1 in [6] do not go through directly in the case of the strip. We shall give a complete proof of (b). The required adaptations in the cases (d) and (e) are the same.

Proof of (b). Set $\chi_n = I_{\rho_n < \delta N}^c$; this notation will be used only within the proof of lemma 6.1. Denote $\tilde{Y}_n = \sum_{\rho_n < \delta N} (\xi_n - \rho_n)$. Then $\mathbb{E}_\omega(\tilde{Y}_n) = 0$ and so it suffices to show that $\operatorname{Var}_\omega(\tilde{Y}_n) = O(\delta^{2-s} N^{2/s})$ except for $\omega$ from a set of small probability. It follows from lemma 3.2 that if $K$ is sufficiently large, $n_2 - n_1 > K \ln N$, and $\omega \notin \Omega_{n_1, K \ln N, \theta}$ (see (3.6) for the definition of this set) then

$$\operatorname{Cov}_\omega(\xi_{n_1}, \xi_{n_2}) \leq C N^{-100}.$$ The last step uses

$$\operatorname{Var}_\omega(\tilde{Y}_n) = o(1) + \sum_{n_2 - K \ln N < n_1 < n_2} 2 \chi_{n_1} \chi_{n_2} \operatorname{Cov}_\omega(\tilde{\xi}_{n_1}, \tilde{\xi}_{n_2}) + \sum_n \chi_n \operatorname{Var}_\omega(\xi_n) \leq 1 + \operatorname{Const} \sum_{n_2 - K \ln N < n_1 < n_2} \rho_{n_1, \rho_{n_2}, n_1, n_2} \quad (6.3)$$

Next, we have to estimate the expectation of the last sum in (6.3). To this end introduce

$$\tau_n = 1 + \|A_{n+1}|| + \|A_{n+2} A_{n+1}|| + \cdots \leq \sum_{j=0}^\infty \|A_{n+j} \cdots A_{n+1}||. \leq \sum_{j=0}^\infty \|A_{n+j} \cdots A_{n+1}||.$$

It is clear from (3.11) and the strong ellipticity condition (1.5) that there are constants $c_1, c_2$ such that $c_1 \tau_n < \nu_n < c_2 \tau_n$. Hence there is a $C$ and $t_0$ such that for $t > t_0$ uniformly in $\xi$

$$P(\tau_n > t|\xi_n = \xi) \leq C \tau^{-s} \quad \text{and} \quad E \left( \phi(\tau_n^2|\xi_n = \xi) \I_{\tau_n^2|\xi_n = \xi} \leq C \right) \leq \operatorname{Const} \tau^{-s}. \quad (6.4)$$

We also have that for $k > 0$

$$\rho_{n-k} \leq C \sum_{j=0}^\infty \|A_{n-k+j} \cdots A_{n-k+1}|| \leq C \sum_{j=0}^{k-1} \|A_{n-k+j} \cdots A_{n-k+1}|| + C \|A_n \cdots A_{n-k+1}|| \tau_n.$$
and
\[ \rho_{n-k}\rho_n \leq C \sum_{j=0}^{k-1} ||A_{n-k+j} \cdots A_{n-k+1}|| ||r_n|| + C ||A_{n-k+1}|| r_n^2. \]

To estimate \( E(t_{n-k}t_n) \) we condition on \( \zeta_n \) and use the fact that the conditioned random variables \( ||A_{n-k+j} \cdots A_{n-k+1}|| \zeta_n \) and \( t_n \zeta_n \) are independent. Therefore
\[ E (||A_{n-k+j} \cdots A_{n-k+1}|| \zeta_n) = E [E (||A_{n-k+j} \cdots A_{n-k+1}|| \zeta_n) E (t_n | \zeta_n)] \leq C \beta \delta^{2-s} N^{\frac{1}{2s}}, \]
where \( \beta := \gamma_1 < 1 \) and \( E(t_n | \zeta_n) \leq \text{Const since } s > 1. \)

Similarly, but this time using also (6.4) we obtain
\[ E (||A_{n+1} \cdots A_{n-k+1}|| \zeta_n) \leq E [E (||A_{n+1} \cdots A_{n-k+1}|| \zeta_n) E (r_{n+1} | t_n, \zeta_n) \leq C \delta^{2-s} N^{\frac{1}{2s}}. \]

Hence
\[ E (\rho_{n-k}\rho_n) \leq C \delta^{2-s} N^{\frac{1}{2s}}. \]
and therefore for \( N > N_\delta \)
\[ E \left( \sum_{k=0}^{K \ln N} \rho_{n-k}\rho_n \right) \leq C \delta^{2-s} \left( \ln N + N^{\frac{1}{2s}} \right). \]
Finally,
\[ E \left( \sum_{n=0}^{N-1} \rho_{n} \rho_{n+1} \rho_{n+2} \right) \leq C N \delta^{2-s} N^{\frac{1}{2s}}. \]
Part (b) is proven. The proofs of parts (d) and (e) follow the proof of the corresponding statement of lemma 4.1 in [6] with the modifications similar to ones presented here.

We are now prepared to explain the main steps of the proof of theorem 6.
We consider the case \( s \in (0, 1) \); other cases are treated similarly. Present the time spent by the walk in \([L_0, L_{N-1}]\) as
\[ T_N = \sum_{n=0}^{N-1} \xi_n = \sum_{n=0}^{N-1} \sum_{i=1}^{m} \xi_{n(i,i)} = S_1 + S_2 + S_3, \]
(6.5)
where
\[ S_1 = \sum_{n: w_n < \delta^{N/2}, n \neq \text{any trap}} \xi_n \]
\[ S_2 = \sum_{n: w_n \geq \delta^{N/2}, n \text{ is not in a trap}} \xi_n \]
\[ S_3 = \sum_{n: n \text{ is in a trap}} \xi_n. \]
By lemma 6.1, (a) we have that \( E(S_1) \leq \text{Const} N^{1/2} \delta^{1-s}. \)
Next, denote
\[ \Omega_{N,A}^{(1)} := \{ \omega : \text{ there is } n \in [0, N-1] \text{ s.t. } w_n \geq \delta N^{1/2}, n \text{ is not in a trap} \}. \]
It follows from (5.2) that
\[ P\left(\tilde{\xi}_{N,\beta}^{(1)}\right) \leq N P\left(\emptyset_n \geq \delta N^{1/\beta} \text{ and } n \text{ is not in a trap}\right) \leq \text{Const}\beta^M. \]

But then
\[ P(S_2 \neq 0) \leq P\left(\tilde{\xi}_{N,\beta}^{(1)}\right) \leq \text{Const}\beta^M \to 0 \quad \text{as } N \to \infty. \]

We thus have that for \( \omega \notin \tilde{\xi}_{N,\beta}^{(1)} \)
\[ t_N = N^{-\frac{1}{\beta}} S_3 + N^{-\frac{1}{\beta}} S_1 = N^{-\frac{1}{\beta}} S_3 + R_N, \]
where \( R_N := N^{-\frac{1}{\beta}} S_1 \) and satisfies the requirements of (a), theorem 6.

It remains to analyse \( S_3 \) which is the main contributor to \( T_N \) coming from the sum over the traps in \([L_0, L_{N+\varepsilon}]\). Let us present it as follows:
\[ N^{-\frac{1}{\beta}} S_3 = \sum_{\text{n marked}} N^{-\frac{1}{\beta}} \sum_{j=0}^{M} \xi_{n-j}. \]

Since \( n \) is marked, we can choose a \( k \) such that \( \rho_{n,k} \geq m^{-1} \delta N^{\frac{1}{\beta}} \). Now present
\[ \sum_{j=0}^{M} \xi_{n-j} = \sum_{j=0}^{m} \left( \xi_{n-j,i} \rho_{n-j,i} - \xi_{n,k} \rho_{n,k} \right) \rho_{n-k} + \sum_{j=0}^{M} \rho_{n-j}. \]

Next, we shall use corollary A.2 to estimate \( \|\xi_{n-j,i}\|/\rho_{n-j,i} \geq \xi_{n,k}/\rho_{n,k} \), where \( \|f\| := \sqrt{\mathbb{E}_n(|f|^2)} \) for a function \( f \) on the space of trajectories of the walk. We have
\[ \left\| \frac{\xi_{n-j,i} - \xi_{n,k}}{\rho_{n-j,i} \rho_{n,k}} \right\| < \sum_{r=n-j}^{n-1} \left\| \frac{\xi_{r,i} - \xi_{r+1,i}}{\rho_{r,i} \rho_{r+1,i}} \right\| C \sum_{r=n-j}^{n-1} \frac{1}{\sqrt{\rho_{r,i}}} + C \frac{1}{\sqrt{\rho_{n,k}}}. \]

Condition (1.5) together with (3.1) imply that there is \( \varepsilon_0 > 0 \) such that \( \rho_{n-j,i} \geq \varepsilon_0 |A_{n-j} \ldots A_n|/|\rho_{n,k}|. \) Hence for \( n - j \) belonging to a trap, that is \( (n - j) \in [n - M, n], \) we have that \( \rho_{n-j,i} \geq c \varepsilon_0^{-1} \rho_{n,k} \geq c N^{-\frac{1}{\beta}} \rho_{n}. \) (Remember that \( M = \ln \ln N \) and therefore these inequalities hold for any \( \varepsilon > 0 \) and \( N > N_0 \).) Thus,
\[ \left\| \sum_{j=0}^{M} \sum_{i=1}^{m} \left( \xi_{n-j,i} / \rho_{n-j,i} \right) \rho_{n-j,i} \right\| \leq \text{Const} \frac{N^{1/2}}{\sqrt{\rho_n}} \sum_{j=0}^{M} \rho_{n-j} \leq \text{Const} \frac{N^{1/2}}{\sqrt{\rho_n}} \sum_{j=0}^{M} \rho_{n-j}. \]

If for \( n \) marked we set
\[ a_n = m_n^{-1} \sum_{j=0}^{M} \sum_{i=1}^{m} \left( \xi_{n-j,i} / \rho_{n-j,i} \right) \rho_{n-j,i} \]
then \( \|a_n\| \leq \text{Const} N^{1/2} / \sqrt{\rho_n} \to 0 \) as \( N \to \infty \) and we have
\[ \sum_{j=0}^{M} \sum_{i=1}^{m} \xi_{n-j-i} / N^{1/\beta} = \left( \xi_{n,k} / \rho_{n,k} \right) + a_n \frac{m_n}{N^{1/\beta}}. \]

Set \( I^{(N,\delta)} = (\xi_{n,k}/\rho_{n,k}) + a_n \), where \( \{n_i\} \) is the collection of marked points. To finish the proof of statement (a) from theorem 6, it remains to check that
\[ \{\xi_{n,k}/\rho_{n,k}\} |_{n_i} \text{ is marked} \]
form a collection of asymptotically independent random variables which also are asymptotically exponential with mean 1.
The convergence to the exponential distribution is an immediate corollary of two facts: the conditional random variable $\xi_n,k| (\xi_n,k \geq 1)$ is geometric and $P_\omega(\xi_n,k = 0) \to 0$ as $N \to \infty$ (to prove this last assertion apply (A.19) with $a = (n, k)$ and $b$ the first point visited by the walker inside layer $L_n$).

To establish the asymptotic independence remember the construction used in the proof of theorem 5. We have established there that the marked points belong to the blocks of length $N^{\epsilon_3}$ which are separated from each other by the blocks of length $N^{-\epsilon_3}$ and, moreover, there is at most one marked point in a large block. By lemma 3.2, the $P_\omega$ probability that the walk would ever return to block $I_{j-1}$ after having reached $I_j$ is of order $O(N^{-\epsilon_3}0)$, where $\theta_0 < 1$, it follows that any random variables which are functions of the part of trajectory of the walk starting at the left end of $I_j$ and restricted to the $N^{\epsilon_3}/2$ neighbourhood of $I_j$ are independent.

Part (a) of theorem 6 is proved. Parts (b) and (c) are dealt with in a similar way.

7. Tail asymptotics

Proof of lemma 3.3. Parts (a) and (b) of the lemma follow from the $z^+$ part of theorem 9 from appendix C applied to the following Markov process

$$\Phi_n = (\pi_n, v_n, \xi_n), \quad g_n = (P_{n+1}, Q_{n+1}, R_{n+1}).$$

(7.1)

Note that due to (1.5) there exists $\bar{\epsilon}$ such that both $A_n$ and $\xi_n$ map $\mathbb{X}$ into $\mathbb{X}_{\bar{\epsilon}} = \{v \in \mathbb{X} : v_j \geq \bar{\epsilon}\}$. In (7.1) $\pi_n$ and $v_n$ are regarded as elements of $\mathbb{X}_{\bar{\epsilon}}$ and $\xi_n$ is an element of the set of stochastic matrices. Recall that given $\Phi_n$ and $g_n$ we can reconstruct $\Phi_{n+1}$ using (2.3), (2.7), (2.9) and (3.12).

In order to apply theorem 9 we need to check three conditions. The first one is eventual contraction (equation (C.1)). Since both $A_n$ and $\xi_n$ map $\mathbb{X}$ into $\mathbb{X}_{\bar{\epsilon}}$ we can apply Birkhoff theorem (see e.g. [22]) which tells us that there is a constant $\bar{\theta} = \bar{\theta}(\bar{\epsilon}) < 1$ such that $A_n$ and $\xi_n$ contract the Hilbert metric on $\mathbb{X}_{\bar{\epsilon}}$ at least by factor $\bar{\theta}$ (the contraction of $\pi_n$ part also follows from lemma 2.2). The contraction of $\xi_n$ part is proven in appendix D.

Second, we need to check (C.2). In our setting we have to show that for each $t$ there is $n$ such that $P(||A_n \ldots A_1|| > t) > 0$. If this were false then there would exist $t_0 > 1$ such that $||A_n \ldots A_1||^\alpha < t_0^\alpha$ for all $n$ with probability 1. This would imply $r(\alpha) \leq 1$ for all $\alpha > 0$. Since $r(0) = 1$ and $\ln r(\alpha)$ is strictly convex we would actually have $r(\alpha) < 1$ for all positive $\alpha$ contradicting (3.15).

Lastly we need to show that (C.3) has no solutions. In our setting (C.3) takes form

$$e^{iu}h(\pi_{n-1}, v_{n-1}, \xi_{n-1}) = e^{iu\ln ||A_n v_{n-1}||}h(\pi_n, v_n, \xi_n).$$

(7.2)

Take $(P, Q, R)$ in the support of the environment measure. Let $\xi_{(P, Q, R)}$ and $A_{(P, Q, R)}$ be defined by (2.4) and (2.5) respectively and denote by $\pi_{(P, Q, R)}$ and $v_{(P, Q, R)}$ the positive eigenvectors of these matrices. Then

$$\Phi_n \equiv (\pi_{(P, Q, R)}, v_{(P, Q, R)}, \xi_{(P, Q, R)}), \quad g_n \equiv (P, R, R)$$

is an admissible trajectory. Evaluating (7.2) along this trajectory we get

$$e^{iu}h(\pi_{(P, Q, R)}, v_{(P, Q, R)}, \xi_{(P, Q, R)}) = e^{iu\ln \lambda_{(P, Q, R)}h(\pi_{(P, Q, R)}, v_{(P, Q, R)}, \xi_{(P, Q, R)})}.$$

From this we conclude that

$$\lambda_{(P, Q, R)} = \frac{\tilde{u}}{u} \in \frac{2\pi}{u} \mathbb{Z}$$

contradicting the non-arithmeticity condition (2.6).
Hence theorem 9 is applicable giving parts (a) and (b) of the lemma. To prove part (c) note that
\[ R_n = \sum_{j \geq n} \pi_j A_j \ldots A_{n+1}[u_n - l_n(u_n)v_n]. \]
Pick a small \( \tilde{\varepsilon} \) and split \( R_n = R'_n + R''_n \) where the first term contain the terms with \( j < n + \tilde{\varepsilon} \ln t \) and the second term contain the terms with \( j \geq n + \tilde{\varepsilon} \ln t \). Choosing \( \tilde{\varepsilon} \) small enough we can ensure that \( R'_n \leq t^2 \).

On the other hand for terms in \( R''_n \) we have \( \theta_n - j \leq t - \tilde{\varepsilon} |\ln \theta_n| \) and hence \( R''_n < C t - \tilde{\varepsilon} |\ln \theta_n| w_n \).

Thus \( P(R_n > t) \leq P(w_n > C t^{1 + \tilde{\varepsilon} |\ln \theta_n|}) \) and so part (c) follows from part (b).

Proof of lemma 5.5. The result follows from \( z^- \) part of theorem 9 applied to the same Markov process (7.1) as in the proof of lemma 3.3.

8. Extensions

Here we discuss some extensions of our results which are not used in the proof of theorem 6 but are helpful in studying other properties of the walk. Applications of these results will be presented in a separate paper.

8.1. Environment inside the trap

Fix \( R \geq 1 \). Let \( T_j = [n_j - M_N, n_j] \) be the \( j \)th trap. We call \( \tilde{n}_j \in T_j \) \( R \)-center of \( T_j \) if \( \tilde{n}_j \) is the rightmost point in \( T_j \) such that \( w_{\tilde{n}_j} > w_n / R \) for all \( n \in T_j \). We choose \( R \) so that for each \( k \) we have
\[ P(\lambda_n \ldots \lambda_{n+k} \in \{ R, R^{-1} \}) = 0. \]
In particular, if for each \( k \) we have \( P(\lambda_n \ldots \lambda_{n+k} = 1) = 0 \) then we can take \( R = 1 \) so that \( \tilde{n}_j \) will be the point with the maximal value of \( w_n \). Denote \( \omega^{(j)} = \tau^{\tilde{n}_j} \omega \), where \( \tau \) is the standard shift on the space of environments. Theorem 5 can be strengthened in the following way.

Theorem 8. Assume that the non-arithmeticity condition (2.6) holds. Then there exists a probability measure \( \tilde{\nu}_3 \) on \( \Omega \) and a constant \( c \) such that the point process
\[ \left\{ \left( \tilde{n}_j / N^{1/3}, \omega^{(j)} \right) \right\} \]
converges as \( N \to \infty \) to a Poisson process on \([0, 1] \times [\delta, \infty) \times \Omega \) with the measure \( cd\nu \mu_{\delta}(m) d\nu_3(\omega) \). As \( \delta \to 0 \) \( \mu_{\delta} \) converges to a measure with density \( c \pi e / m^{s+1} \) and \( \tilde{\nu}_3 \) converges to some measure \( \tilde{\nu} \).

In other words if the walker is trapped, then he sees the environment distributed according to a measure \( \tilde{\nu} \). This statement extends the results obtained in [13, 16].

To prove theorem 8 we first show that as \( N \to \infty \) and then \( \delta \to 0 \)
\[ \left\{ \left( \tilde{n}_j / N^{1/3}, w_{\tilde{n}_j} / N^{1/3}, \omega^{(j)} \right) \right\} \]
converges to a Poisson process with measure \( \tilde{c} dw / w^{1+s} \). (8.2)

The proof of this result is similar to the proof of theorem 5. Namely, call \( n \) \( (R, l) \)-maximal if \( w_n > w_{n+k}/R \) for \( 0 < k \leq l \) and for each \( 0 < k' < l \) there exists \( k'' \leq l \) such that \( w_{n+k''} \leq w_{n+k'}/R \).
Using the same argument as in lemma 5.3 one can show that
\[ P\left(\frac{u_n}{t} \in [\tilde{a}, \tilde{a}] \right) \text{ and } n = (R, l) \text{ maximal } |\tilde{\zeta}_n| = \tilde{\pi}_n = \tilde{\pi}, \nu_n = \tilde{\nu} \sim \tilde{f}_t(\tilde{\zeta}, \tilde{\pi}, \tilde{\nu}) \left(\tilde{a}^{-s} - \tilde{a}^{-s}\right) t^{-s} \]
and \( \tilde{f}_t \to \tilde{f} \) as \( l \to \infty \). Moreover similarly to lemma 5.4 one can show that \( P(\tilde{f} > 0) > 0 \) (otherwise we would get a contradiction with the fact that each trap has a center). In addition we have that for each \( \tilde{g}_{-l} \cdots \tilde{g}_l \)
\[ P\left(\frac{u_n}{t} \in [\tilde{a}, \tilde{a}] \right) \text{ and } n = (R, l) \text{ maximal } |\tilde{\zeta}_n| = \tilde{\pi}_n = \tilde{\pi}, \nu_n = \tilde{\nu}, g_{n-l} = \tilde{g}_{-l} \cdots g_{n+l} = \tilde{g}_l \)
\[ \sim \tilde{f}_t(\tilde{\zeta}, \tilde{\pi}, \tilde{\nu}, \tilde{g}_{-l}, \tilde{g}_l) P_{\tilde{\tau}_{n\tilde{\nu}}, \tilde{\pi}_n, \tilde{\nu}_n} \left(\frac{u_n}{t} \in \left[\tilde{a}l, \tilde{a}l\right]\right) \]
This implies that if \( B \) is \( \mathcal{F}_{-l,l} \) measurable then
\[ P\left(\frac{u_n}{t} \in [\tilde{a}, \tilde{a}], n = (R, l) \text{ maximal } \tau^n \omega \in B \right) \sim \tilde{v}_t(B) \tilde{c}(\tilde{a}^{-s} - \tilde{a}^{-s}) t^{-s} \]
and \( \tilde{v}_t \Rightarrow \tilde{v} \) as \( l \to \infty \). Now the proof of (8.2) proceeds similarly to the proof of theorem 5.
To pass from (8.2) to (8.1) we note that \( m_j/w_{nj} \) is well approximated by \( \sum_{|k|<l} t_{n+k}(u_{n+k})\Lambda_{n,k} \) provided that \( l \) is sufficiently large. Here
\[ \Lambda_{n,k} = \begin{cases} \lambda_{n+1} \cdots \lambda_{n+k} & \text{if } k > 0, \\ 1 & \text{if } k = 0, \\ \lambda_{n+k+1} \cdots \lambda_n & \text{if } k < 0. \end{cases} \]
Accordingly, in the limit \( l \to \infty \), we have \( m_j = w_{nj} \mathcal{H}(\omega^{(j)}) \) for some measurable function \( \mathcal{H} \). Now lemma B.1 shows that (8.2) implies (8.1) with \( d\tilde{v} = \mathcal{H}^t d\tilde{v}^s \).

8.2. Arithmetic case

We note that condition (2.6) was used in section 7 to show that (7.2) does not have solutions. On the other hand if (7.2) has a non-trivial solution then the analysis of appendix C has to be modified. Namely, the non-arithmetic local limit theorem (lemma C.3) has to be replaced by its arithmetic version. This will cause replacing \( t^{-s} \) in the estimates of theorem 9 by \( t^{-s} g(|\ln(t/\Delta)|) \) where \( \{ \ldots \} \) denotes the fractional part, \( \Delta \) is the step of the progression containing the distribution of \( \ln \lambda \) and \( g \) is some continuous function. As a result the estimates of section 5 have to be replaced by
\[ P(\rho_n > t) \sim t^{-s} g_1(|\ln(t/\Delta)|), \quad P(m_j > t) \sim t^{-s} g_2(|\ln(t/\Delta)|). \]
Thus there would exists a measure \( \mu \) on \( \mathbb{R}^+ \) such that \( \mu([t, \infty)) = \tilde{g}(|\ln(t/\Delta)|) t^{-s} \) and the limit points of the distribution of the normalized hitting times will be of the form
\[ \sum_j \Theta_j(\Gamma_j = 1) \]
where \( \Gamma_j \) are iid mean 1 exponential random variables and \( \Theta_j \) is a Poisson process with measure \( \mu_\Delta \) for some \( 0 \leq \Delta < \Delta \) where \( \mu_\Delta(A) = \mu(e^\Delta A) \).

In particular, we would like to note that regardless of condition (2.6) we always have
\[ P(\rho_n > t) \leq Ct^{-s}. \] (8.3)
Appendix A. Occupation times for Markov chains

We recall two facts about general Markov chains with discrete state space. Within this section the notation $\mathbb{P}$ will be used for probabilities concerned with Markov chains and $\mathbb{E}$ for the corresponding expectations. First, the number of visits to a given state conditioned on the event that this state will be visited has a geometric distribution. Second, consider a Markov chain with transition probabilities $p_{ij}$. Let $\tilde{p}_{jk}$ be the probability that the chain starting at $j$ ever visits $k$. Condition the chain on having at least one visit to $k$. Then before coming to $k$ the chain evolves as a Markov chain with transition probabilities

$$ p_{ij}^* = \frac{p_{ij}\tilde{p}_{jk}}{\sum_r p_{ir}\tilde{p}_{rk}}. \quad (A.1) $$

Similarly if we condition the chain on never visiting $k$ then the chain evolves as a Markov chain with transition probabilities

$$ p_{ij}^{**} = \frac{p_{ij}(1 - \tilde{p}_{jk})}{\sum_r p_{ir}(1 - \tilde{p}_{rk})}. \quad (A.2) $$

We shall now use these facts to analyse the joint distribution for the number of visits to different sites.

Namely let $a$ and $b$ be two states of a transient chain $Z$ such that

$$ p_{ab}^{n_0} > \epsilon, \quad p_{ba}^{n_0} > \epsilon, \quad (A.3) $$

where $p_{ab}^{n_0}$ denotes the transition probability after $n_0$ steps.

Let $q_a$ ($q_b$) denote the probability that $a$ (respectively $b$) is visited at least once and $p_a$ ($p_b$) denote the probability that the chain started from $a$ (respectively $b$) does not return to that state again. Let $\xi_a$ ($\xi_b$) be the number of visits to $a$ (respectively $b$).

It is useful to consider our Markov chain only at the times when it visits either $a$ or $b$. Denote the resulting chain by $\tilde{Z}$. This chain is governed by the transition matrix

$$
\begin{pmatrix}
r_{aa} & r_{ab} & r_{ac} \\
r_{ba} & r_{bb} & r_{bc} \\
0 & 0 & 1
\end{pmatrix}, \quad (A.4)
$$

where $c$ denotes the absorbing state which is reached by the particle after the last visit of the set $\{a, b\}$. Denote $\bar{\epsilon} = \epsilon/n_0$. Since $\tilde{Z}$ is obtained from $Z$ by skipping some states, (A.3) implies that

$$\mathbb{P}(\tilde{Z} \text{ visits } b \text{ before } n_0|\tilde{Z}_0 = a) \geq \epsilon. \quad (A.5)$$

(Note that here the time which is $\leq n_0$ is that of the Markov chain $\tilde{Z}$, not $Z$.) On the other hand this probability is equal to

$$\sum_{j=0}^{n_0-1} r_{ab}^j r_{ab} \leq n_0 r_{ab}. $$

This implies that

$$ r_{ab} \geq \bar{\epsilon} \quad \text{and by symmetry} \quad r_{ba} \geq \bar{\epsilon}. \quad (A.5) $$

**Lemma A.1.** Suppose that (A.3) is satisfied. Then there is a constant $C$ depending on $n_0, \epsilon > 0$ such that for all initial conditions of the chain satisfying

$$ q_a > \epsilon, \quad q_b > \epsilon \quad (A.6) $$

the following holds:

$$ \text{Corr}(\xi_a, \xi_b) > 1 - \frac{C}{\mathbb{E}(\xi_a)}. \quad (A.7) $$
Remark. It is implicit in the statement of the lemma that the constant \( C \) above is uniform over all Markov chains satisfying (A.3) and (A.6).

Remark. Remembering that
\[
\mathbb{E}(\xi_a) = \frac{q_a}{p_a},
\]
we see that (A.7) is equivalent to
\[
\text{Corr}(\xi_a, \xi_b) > 1 - \tilde{C} p_a,
\]
where \( \tilde{C} \) depends only on \( \varepsilon \) and \( n_0 \). Therefore this estimate is only interesting if \( p_a \) is very small. In this paper we apply lemma A.1 to traps where the walker spends a lot of time so that the probability of not returning to a site in a trap is small.

Of course by symmetry we also have
\[
\text{Corr}(\xi_a, \xi_b) > 1 - \tilde{C} p_b.
\]

However it is not difficult to see directly that \( p_a \) and \( p_b \) are of the same order. Namely we have
\[
p_a = r_{ac} + r_{ab} \sum_{j=0}^{\infty} r_{bb}^j r_{bc} = r_{ac} + r_{ab} (1 - r_{bb})^{-1} r_{bc}
\]
and similarly \( p_b = r_{bc} + r_{ba} (1 - r_{aa})^{-1} r_{ac} \). It is now obvious that \( p_a \leq r_{ac} + (1 - r_{bb})^{-1} r_{bc} \leq r_{bc} + \frac{r_{bc}}{\tilde{e}} \) and \( p_b \geq r_{bc} + r_{ba} r_{ac} \geq r_{bc} + \tilde{e} r_{ac} \).

This implies that
\[
p_a / p_b \leq \tilde{e}^{-1} \text{ and by symmetry } p_b / p_a \leq \tilde{e}^{-1}.
\]

Proof. We have \( \xi_b = U + V + W \), where \( U \) is the number of visits to \( b \) before the first visit to \( a \), \( W \) is the number of visits to \( b \) after the last visit to \( a \) and \( V = \sum_{j=1}^{n} V_j \), where \( V_j \) is the number of visits to \( b \) between \( j \)-th and \( j + 1 \)-st visit to \( a \). Then \( V_j \) are iid. Let \( v_{ab} = \mathbb{E}(V_j) \).

We claim that the following uniform bounds hold
\[
v_{ab} = O(1), \quad \mathbb{E}(U) = O(1), \quad \mathbb{E}(W) = O(1)
\]
where the implied constants depend only on \( \varepsilon \) and \( n_0 \). Indeed, by (A.1)
\[
P(V_j = 1) = r_{ab}^* (r_{bb}^*)^j r_{ba}^*
\]
so the estimate of \( v_{ab} \) follows from the fact that \( r_{bb}^* = 1 - r_{ba}^* < 1 - r_{ba} < 1 - \tilde{e} \). The estimate of \( \mathbb{E}(U) \) is the same and the estimate of \( \mathbb{E}(V) \) is similar except that we have to use (A.2) instead of (A.1).

Therefore
\[
\mathbb{E}(\xi_b|\xi_a = k + 1) = kv_{ab} + O(1).
\]

Hence
\[
\mathbb{E}(\xi_b) = v_{ab} \mathbb{E}(\xi_a) + O(1) \text{ and by symmetry } \mathbb{E}(\xi_a) = v_{ba} \mathbb{E}(\xi_b) + O(1),
\]
\[
\mathbb{E}(\xi_a \xi_b) = v_{ab} \mathbb{E}(\xi_a^2) + O(\mathbb{E}(\xi_a)).
\]
Combining the last two equalities we obtain
\[
\text{Cov}(\xi_a, \xi_b) = v_{ab} \text{Var}(\xi_a) + O(\mathbb{E}(\xi_a)) = v_{ab} \text{Var}(\xi_a) \left( 1 + O \left( \frac{1}{\mathbb{E}(\xi_a)} \right) \right)
\]
where the last step uses (A.8) and
\[ \text{Var}(\xi_a) = \frac{q_a(2 - q_a - p_a)}{p_a^2}. \] (A.15)

It follows from (A.8) and (A.11) that
\[ n_0^{-1} \varepsilon^2 \leq \mathbb{E}(\xi_a)\mathbb{E}(\xi_b)^{-1} \leq n_0 \varepsilon^{-2}. \] (A.16)

Therefore interchanging roles of \( a \) and \( b \) in (A.14) we get
\[ \text{Cov}(\xi_a, \xi_b) = v_{ba} \text{Var}(\xi_b) \left(1 + \mathcal{O} \left( \frac{1}{\mathbb{E}(\xi_a)} \right) \right). \] (A.17)

Multiplying the two expressions in (A.13) we obtain
\[ v_{ab}v_{ba} = 1 + \mathcal{O} \left( \frac{1}{\mathbb{E}(\xi_a)} \right). \] (A.18)

Finally, multiplying (A.14) and (A.17) and using (A.18) we get
\[ \frac{\text{Cov}^2(\xi_a, \xi_b)}{\text{Var}(\xi_a)\text{Var}(\xi_b)} = 1 + \mathcal{O} \left( \frac{1}{\mathbb{E}(\xi_a)} \right). \]

Let \( \| \ldots \| \) denote the \( L^2 \) norm.

**Corollary A.2.** Under the conditions of Lemma A.1 there exists a constant \( \bar{C} \) (depending only on \( \varepsilon \) and \( n_0 \)) such that
\[ \left\| \frac{\xi_a}{\mathbb{E}(\xi_a)} - \frac{\xi_b}{\mathbb{E}(\xi_b)} \right\| \leq \bar{C} \left[ \frac{1}{\sqrt{\text{Var}(\xi_a)}} + (1 - q_a) + (1 - q_b) \right]. \]

**Proof.** We have
\[ \left\| \frac{\xi_a}{\mathbb{E}(\xi_a)} - \frac{\xi_b}{\mathbb{E}(\xi_b)} \right\| = \left\| \frac{\xi_a - \mathbb{E}(\xi_a)}{\mathbb{E}(\xi_a)} - \frac{\xi_b - \mathbb{E}(\xi_b)}{\mathbb{E}(\xi_b)} \right\| = \left\| \frac{\xi_a - \mathbb{E}(\xi_a)}{\sqrt{\text{Var}(\xi_a)}} - \frac{\xi_b - \mathbb{E}(\xi_b)}{\sqrt{\text{Var}(\xi_b)}} \right\| \leq \frac{\xi_a - \mathbb{E}(\xi_a)}{\sqrt{\text{Var}(\xi_a)}} + \frac{\xi_b - \mathbb{E}(\xi_b)}{\sqrt{\text{Var}(\xi_b)}}. \]

Notice that the first term equals to \( \sqrt{2(1 - \text{Corr}(\xi_a, \xi_b))} \) and so it can be estimated by lemma A.1 while the last two terms can be estimated by (A.8) and (A.15).

To use the above corollary we need to estimate \( 1 - q_a \). To this end we observe that if (A.3) holds then
\[ \mathbb{P}_b(\xi_a = 0) = \sum_{l=0}^{\infty} t_{ba} r_{bc} = \frac{r_{bc}}{1 - r_{bb}} \leq \frac{p_b}{r_{ba}} \leq \frac{p_b}{\varepsilon}. \] (A.19)
Appendix B. Poisson process and stable distributions.

Let \((X, \mu)\) be a measure space. Recall that a Poisson process is a point process with values in \(X\) such that if \(N(A)\) is the number of points in \(A \subset X\) then \(N(A_1), N(A_2) \ldots N(A_k)\) are mutually independent if \(A_1, A_2 \ldots A_k\) are disjoint and \(N(A)\) has the Poisson distribution with parameter \(\mu(A)\). If \(X \subset \mathbb{R}^d\) and \(\mu\) has density \(f\) with respect to the Lebesgue measure we say that \(f\) is the intensity of the Poisson process.

**Lemma B.1** (see [19], sections 2.3 and 5.2).

(a) If \(\{\Theta_j\}\) is a Poisson process on \(X\) and \(\psi: X \rightarrow \tilde{X}\) is a measurable map then \(\tilde{\Theta}_j = \psi(\Theta_j)\) is a Poisson process. If \(X = \tilde{X} = \mathbb{R}\) and \(\psi\) is invertible differentiable map then the intensity of \(\tilde{\Theta}_j\) is \(\tilde{f}(\theta) = f(\psi^{-1}(\theta)) \frac{|d\psi|}{d\theta}^{-1}\).

(b) Let \((\Theta_j, \Gamma_j)\) be a point process on \(X \times \mathbb{Z}\) such that \(\{\Theta_j\}\) is a Poisson process on \(X\) and \(\{\Gamma_j\}\) are \(\mathbb{Z}\)-valued random variables which are i.i.d. and independent of \(\{\Theta_k\}\) then \((\Theta_j, \Gamma_j)\) is a Poisson process on \(X \times \mathbb{Z}\).

(c) If in (b) \(X = \mathbb{Z} = \mathbb{R}\) then \(\tilde{\Theta} = \{\Gamma_j \Theta_j\}\) is a Poisson process. Its intensity is \(\tilde{f}(\theta) = E \left( f \left( \frac{\theta}{\Gamma} \right) \frac{1}{\Gamma} \right)\).

**Lemma B.2** (see [28], theorem 1.4.2).

(a) If \(0 < s < 1\) and \(\Theta_j\) is a Poisson process with intensity \(\theta^{-(1+s)}\) then \(\sum_j \Theta_j\) has stable distribution of index \(s\).

(b) If \(1 < s < 2\) and \(\Theta_j\) is a Poisson process with intensity \(\theta^{-(1+s)}\) then

\[
\lim_{\delta \to 0} \left[ \sum_{j \in \Theta_j} \Theta_j - \frac{1}{(s-1)^{\delta^{-1}}} \right]
\]

has stable distribution of index \(s\).

(c) If \(s = 1\) and \(\Theta_j\) is a Poisson process with intensity \(\theta^{-2}\) then

\[
\lim_{\delta \to 0} \left[ \sum_{j \in \Theta_j} \Theta_j - |\ln \delta| \right]
\]

has stable distribution of index \(1\).

Appendix C. Renewal theorem for a system of contractions

**C.1. Main result**

Let \(M_1\) and \(M_2\) be compact metric spaces, and let \(\omega = \{g_j\}\) be a sequence of iid \(M_2\) valued random variables. Suppose that there is a \(C^0\) map \(G: M_1 \times M_2 \rightarrow M_1\). Here and below \(C^0\) denotes the space of Holder maps, that is

\[
d(G(\Phi_1, g_1), G(\Phi_2, g_2)) < C \left[ d^0(\Phi_1, \Phi_2) + d^0(g_1, g_2) \right].
\]

Note that later we impose a stronger condition on \(\Phi\) dependence (see (C.1) below). Consider Markov process on \(M_1\)

\[
\Phi_j +1 = G_j(\Phi_j), \quad \text{where } G_j(\Phi) = G(\Phi, g_j).
\]
We suppose that the maps $G_j$ are contractions in the sense that there exist constants $C$ and $\theta < 1$ such that if $\Phi_j'$ and $\Phi_j''$ are two realizations of this Markov chain starting from $\Phi_0'$ and $\Phi_0''$ and evolving in the same environment $\omega$ then with probability 1 we have

$$d(\Phi_j', \Phi_j'') < C^{\theta^j}d(\Phi_0', \Phi_0'').$$

(C.1)

Denote $\tilde{\omega} = \{(g_j, \Phi_j)\}$, $\Omega = (M_2)^\mathbb{Z}$, $\tilde{\Omega} = (M_1 \times M_2)^\mathbb{Z}$. Let $P$ be the distribution of $\omega$. In view of (C.1) there is a unique stationary distribution $\pi$ such that if $\theta < 1$

Theorem 9. Suppose that for any numbers $a$, $\bar{a}$ denote

$$\lim_{k \to \infty} (\Phi_0, g).$$

Then there are constants $s > 0, \tilde{s} > s$ and $\tilde{C} > 0$ such that

(a) If $k > \tilde{C} \ln t$ then

$$P(z_n^+ - z_{n,k}^+ > 1) < \tilde{C} t^{-\tilde{s}}, \quad P(z_n^- - z_{n,k}^- > 1) < \tilde{C} t^{-\tilde{s}}.

(b) There exists a function $f(\Phi)$ such that if $k > \tilde{C} \ln t$ then

$$P(z_n^+ > t|\Phi_0 = \Phi) \sim f(\Phi)^{-s}.

(Here and below ‘$\sim$’ means that the estimates are uniform in $k$, that is given $\varepsilon$ there exists $t_0$ such that if $t > t_0$ and $k > \tilde{C} \ln t$ then

$$|t^{-s}P(z_n^+ > t|\Phi_0 = \Phi) - f(\Phi)| < \varepsilon.$$

In particular

$$P(z_n^+ > t|\Phi_0 = \Phi) \sim f(\Phi)^{-s}.

(c) There exists a measure $\tilde{v}$ on $\tilde{\Omega}$ such that for any continuous function $H$ on $\tilde{\Omega}$ the following asymptotic relations hold if $k > \tilde{C} \ln t$ and $t \to \infty$:

$$\tilde{E}(1_{z_n^+ > t} H(\tilde{\omega})) \sim t^{-\tilde{s}} \tilde{v}(H).

In particular

$$\tilde{E}(1_{z_n^+ > t} H(\tilde{\omega})) \sim t^{-\tilde{s}} \tilde{v}(H).
C.2. Renewal theorem and large deviations

We will deduce theorem 9 from a large deviation bound. Let $y_n = \sum_{j=0}^{n-1} a_j$.

**Theorem 10.** Suppose that (C.3) has no solutions. Then there is a number $\alpha^* > 0$ and a strictly convex analytic function $\gamma : [\tilde{E}(\alpha), \alpha^*) \to \mathbb{R}$ such that

$$
\gamma'(\alpha) > 0 \text{ if } \alpha > \tilde{E}(\alpha), \quad \lim_{\alpha \to \alpha^*} \gamma'(\alpha) = +\infty
$$

(C.4)

and

(a) If $\alpha > \alpha^*$ then for each $\beta \in \mathbb{R}$, $\mathbb{P}_\Phi(y_n \geq \alpha n) = O(e^{-\beta n})$;

(b) If $\alpha \in (\tilde{E}(\alpha), \alpha^*)$ then there is a measure $\nu_{\alpha}$ on $\tilde{\Omega}$ and a function $h_{\alpha} : M_1 \to \mathbb{R}$ such that

$$
\mathbb{E}_\Phi \left( 1_{y_n \in J} H \left( t^n \tilde{\omega} \right) \right) \sim \frac{e^{-\gamma(\alpha)n}}{\sqrt{n}} \nu_{\alpha}(\Phi) \int_J e^{-\gamma'(\alpha)t} \gamma'(\alpha) \, dt
$$

(c) If $\alpha \in (\tilde{E}(\alpha), \alpha^*)$ then for any $J \subset [0, \infty]$ and any continuous functions $H, \tilde{H} : \tilde{\omega} \to \mathbb{R}$

$$
\mathbb{E} \left( 1_{y_n \in J} \tilde{H} \left( t^n \tilde{\omega} \right) H \left( t^n \tilde{\omega} \right) \right) \sim \frac{e^{-\gamma(\alpha)n}}{\sqrt{n}} \nu_{\alpha}(\tilde{\Phi}) \nu_{\alpha}(H) \int_J e^{-\gamma'(\alpha)t} \gamma'(\alpha) \, dt
$$

Theorem 10 is proven in section C.3. Here we use this theorem to obtain theorem 9.

Let $\alpha_0 = \arg\min(\gamma(\alpha)/\alpha)$ and $s = \gamma(\alpha_0)/\alpha_0$. Note that by (C.4) the minimum is strictly inside $(0, \alpha^*)$. Indeed from part (b) of theorem 10 it is clear that $\gamma$ is non-negative on $(\tilde{E}(\alpha), \alpha^*)$ and since it is strictly increasing it follows that it is in fact strictly positive on this interval. Accordingly $\lim_{\alpha \to \alpha^*} \gamma(\alpha)/\alpha = +\infty$ and so there exists $\epsilon_1 > 0$ such that the minimum cannot be achieved on $[0, \epsilon_1]$. Next we show that the minimum cannot be achieved near $\alpha^*$. Notice that $\gamma$ is increasing. Now if $\lim_{\alpha \to \alpha^*} \gamma(\alpha) = +\infty$ then the statement is obvious. Otherwise the formula

$$
\left( \frac{\gamma(\alpha)}{\alpha} \right)' = \frac{\gamma'(\alpha)}{\alpha} - \frac{\gamma(\alpha)}{\alpha^2}
$$

shows that there exists some $\epsilon_2 > 0$ such that $\gamma(\alpha)/\alpha$ is increasing on $(\alpha^* - \epsilon_2, \alpha^*)$ proving our claim that $s \in (0, \alpha^*)$.

**Lemma C.1.** There exists constants $C$ and $s > 0$ such that

(a) for each $\Phi$ we have $\mathbb{P}_\Phi(\max y_n \geq Y) \leq C e^{-sY}$;

(b) For each $n$ we have $\mathbb{P}_\Phi(\exists n : y_n \geq Y \text{ and } y_{n+1} \geq Y) \leq \left( C / \sqrt{n} \right) e^{-s Y + \tilde{s} Y}$. In particular,

$$
\mathbb{P}_\Phi(\exists n_1, n_2 : n_2 > n_1 + k \text{ and } y_{n_j} \geq Y \text{ for } j = 1, 2) \leq C e^{-(\alpha + \tilde{s} Y + \tilde{s} k)}.
$$

**Proof.** By theorem 10

$$
\mathbb{P}_\Phi(y_n \geq Y) \leq \frac{C}{\sqrt{n}} \exp \left[ - \left( \frac{\gamma(Y/n)}{Y/n} \right) Y \right] \leq \frac{C}{\sqrt{n}} \exp[-sY] \exp \left[ -c \left( \frac{Y}{n} - \alpha_0 \right)^2 Y \right].
$$

(C.5)

The main contribution to this sum comes from $n \approx Y/\alpha_0$. For those $n$

$$
c \left( \frac{Y}{n} - \alpha_0 \right)^2 Y \leq c (Y - \alpha_0 n)^2.
$$

Since

$$
\sum_n \frac{1}{\sqrt{Y}} \exp \left[ -c (Y - \alpha_0 n)^2 / Y \right] \leq C
$$

part (a) is proved.
To prove (b) observe that by Markov property
\[ P_{\Phi}(y_n \geq Y \text{ and } y_{n+1} \geq Y) < \frac{C}{\sqrt{Y}} e^{-\gamma Y + \varepsilon k} \leq P_{\Phi}(y_n \geq Y) \max_{\Phi(y_l > 0)} \]

The second term is less than \((C/\sqrt{Y})e^{-\gamma Y}\) due to theorem 10 while the first term is less than \((C/\sqrt{Y})e^{-\gamma Y}\) by part (a). Now the first inequality of part (b) follows by summation over \(l > k\) and the second one follows by summation over \(n\) and \(l\).

\[ \square \]

Lemma C.2. Suppose that \(n, Y \to \infty\) so that \((n - (Y/\alpha_0)/\sqrt{Y}) \to \beta\) Denote
\[ \Omega_n = \{y_n \geq Y, y_m < Y \text{ for all } 0 < m < n\}. \]

Then for each \(\Phi\) and a continuous function \(\tilde{H} : \mathbb{R} \times \tilde{\Omega} \to \mathbb{R}\) the following limits exist

(a) \[ \lim_{n \to \infty} \int P_{\Phi}(\Omega_n) \sqrt{Y} e^{\beta Y} ; \]

(b) \[ \lim_{n \to \infty} \mathbb{E}_{\Phi}(1_{\Omega_n} \tilde{H}(y_n - Y, \varepsilon)) \sqrt{Y} e^{Y}. \]

Moreover both limits are bounded by \(\text{Const} e^{-c\beta^2}\).

**Proof.** (a) By lemma C.1(b) it is enough to show that for each \(k\) the following limit exists

\[ \lim_{n \to \infty} P_{\Phi}(y_n \geq Y, y_m < Y \text{ for } n - k < m < n) \sqrt{Y} e^{Y}. \]

The limiting expression equals to

\[ \int \mathbb{P}_{\Phi}(y_n-k \in [Y-z_1, Y-z_2]) dP^\Phi_{z_1, z_2} \sqrt{Y} e^{Y}, \]

where \(P^\Phi_{z_1, z_2}\) is the distribution function of the random vector \((y_k, \max_{0<j<k} y_j)\) for our Markov chain started from \(\Phi\). The integral above is the limit of Lebesgue–Stiltjes sums where each term has form

\[ \mathbb{P}_{\Phi}(y_n-k \in [Y-z_1, Y-z_2]) \mathbb{P}_{\Phi_{z_1, z_2}, \epsilon}(y_k \in [z_1, z_1+\epsilon], \max_{0<j<k} y_j \in [z_2, z_2+\epsilon]) \sqrt{Y} e^{Y}. \]

Since for each \(z_1, z_2\) the last probability is a function of \(\Phi_{n-k}\) part (a) follows from theorem 10.

For part (b) it is sufficient to restrict our attention to a dense set of functions \(\tilde{H}\). In particular we can assume that \(\tilde{H}\) depends only on finitely many coordinates, that is we need to compute the limit

\[ \lim_{n \to \infty} \mathbb{E}_{\Phi}(1_{\Omega_n} \tilde{H}(\Phi_{n-k}, \Phi_{n-k+1}, \ldots, \Phi_{n-1}, \Phi_{n})) \sqrt{Y} e^{Y}. \]

for some \(\tilde{H} : (M_1 \times M_2)^{2k+1} \to \mathbb{R}\). The analysis of the last limit is the same as the analysis of (C.6).

Finally the fact that above limits are \(O(e^{-c\beta^2})\) follows from the estimates in the proof of lemma C.1 (see (C.5)).

\[ \square \]

**Proof of theorem 9.** Let \(Y = \ln t\). Note that \(z^*_n - z^*_{n,k} = \sum_{j=k+1}^{\infty} \Delta j e^{Y}\). Take small \(\varepsilon > 0\). Then if \(t\) is sufficiently large, the inequality \(z^*_n - z^*_{n,k} \geq 1\) implies that there is \(j > k\) such that \(y_j > (1-\varepsilon)Y - \varepsilon j\). Note that \(-\varepsilon j < (1-\varepsilon)Y - \varepsilon j < \varepsilon j\) provided \(\tilde{C}\) is large enough. Hence

\[ \mathbb{P}_{\Phi}(y_j > (1-\varepsilon)Y - \varepsilon j) \leq \frac{C}{\sqrt{j}} e^{-\tilde{C} j}, \]

where \(\tilde{C} = \min_{r>0} \frac{1}{r} Y\). Summing over \(j\) we get

\[ \mathbb{P}_{\Phi}(z^*_n - z^*_{n,k} > 1) \leq C(\varepsilon) e^{-\tilde{C} k}. \]

Finally the fact that above limits are \(O(e^{-c\beta^2})\) follows from the estimates in the proof of lemma C.1 (see (C.5)).
Since $k \geq \bar{\Lambda} \ln t$ and $\bar{\Lambda}$ can be taken sufficiently large this proves the first inequality of part (a). The proof of the second inequality is similar.

To prove part (b) take $M \gg 1$. We claim that terms with $y_m \leq Y - M$ can be ignored. Indeed for terms with $Y - M - 1 < y_m < Y - M$ to make a contribution greater than $e^{-M/2} C$ such terms. By lemma C.1 the probability of such an event is

\[ \mathcal{O}(\exp \left[-s(Y - M) + \tilde{s} \exp(M/2)\right]) \]

which establishes our claim. Therefore for large $M$ and $l$ we can approximate

\[ P_{\Phi,1}(z + n, k \geq t) \]

by

\[ \sum_n P_{\Phi,1}(1_{B_{n,l,\varepsilon}}) t^s \]

where

\[ B_{n,l} = \left\{ \sum_{|m-n| < l} \Delta_m e^{\nu_m} \geq 1 - \varepsilon \right\} . \]

The fact that the last sum approaches the limit as $t \to \infty$ follows from lemma C.2(b). This proves part (b). The proof of part (c) is similar.

C.3. Large deviations

We follow the approach of [32].

To simplify the notation we assume for the rest of this section that (C.1) holds with $C = 1$, the general case can be reduced to this one by considering our Markov chains only at the times which are multiples of a sufficiently large $n_0$.

For $\kappa \geq 0$ consider operators $P_\kappa$ given by

\[ P_\kappa(h)(\Phi) = E_{\Phi}(e^{\kappa a(\Phi, g)} h(\Phi_1)) \]

$P_\kappa$ is a positive operator preserving the space of $C^\eta$ functions. Moreover it has many invariant cones as we describe below. Let

\[ \mathcal{C}_K = \{ h \geq 0 : \text{ for all } \bar{\Phi}, \bar{\Phi} \text{ we have } h(\bar{\Phi}) \leq e^{K^* a(\Phi, \bar{\Phi})} h(\bar{\Phi}) \} . \]

A direct computation (using (C.1) with $C = 1$) shows that $P_\kappa(\mathcal{C}_K) \subset \mathcal{C}_{K^*}$ where $K^* = K^* + \kappa H$ and $H$ is the Holder constant of $a$ with respect to $\Phi$ variable. Now [22] shows that if $K$ is so large that $K > K^*$ then $P_\kappa$ contracts the Hilbert metric on $\mathcal{C}_K$ and so there exist positive eigenfunctions

\[ P_\kappa h = e^{\lambda_\kappa} h \]

in $\mathcal{C}_K$ and moreover for any two elements $h', h''$ of $\mathcal{C}_K$ the directions of $P_\kappa h'$ and $P_\kappa h''$ converge to each other exponentially fast. This in turn implies that the rest of the spectrum of $P_\kappa$ is contained in a disc of radius strictly smaller than $e^{\lambda_\kappa}$. Since $e^{\lambda_\kappa}$ is an isolated eigenvalue of $P_\kappa$, $\lambda_\kappa$ depends analytically on $\kappa$.

We need the fact that the map $\kappa \to \lambda_\kappa$ is strictly convex. To see this we need formulas for derivatives of $\lambda$ with respect to $\kappa$. To this end let $v_\kappa$ be the eigenvector of the adjoint operator

\[ v_\kappa(P_\kappa h) = e^{\lambda_\kappa} v_\kappa(h) . \]

Differentiating (C.7) we get

\[ \mathbb{E}_\Phi (a(\Phi, g) e^{\kappa a(\Phi, g)} h(\Phi_1)) + P_\kappa(h'_\kappa) = \lambda_\kappa e^{\lambda_\kappa} h_\kappa + e^{\lambda_\kappa} h'_\kappa . \]

Applying $v_\kappa$ to both sides we get

\[ \lambda_\kappa = \frac{v_\kappa(\mathbb{E}(a(\Phi, g) e^{\kappa a(\Phi, g)} h_k(\Phi_1)))}{v_\kappa(e^{\lambda_\kappa} h_k)} . \]
Let
\[ \tilde{a}_\kappa = \kappa a - \lambda_\kappa + \ln h_\kappa - \ln h_\kappa'(G(\Phi_\kappa, g)). \]

Then
\[ E_\Phi \left( e^{\tilde{a}_\kappa} \right) = 1 \tag{C.9} \]
so we can consider a Markov chain with generator
\[ \tilde{P}_\kappa (h)(\Phi) = E_\Phi \left( e^{\tilde{a}_\kappa h(\Phi_1)} \right). \]

Observe that
\[ \tilde{P}_\kappa = e^{-\lambda_\kappa} M_\kappa^{-1} P_\kappa M_\kappa \]
where \( M_\kappa \) denotes the multiplication by \( e^{\tilde{a}_\kappa} \) so the eigenvector of the adjoint operator (which is the stationary measure for our Markov process) equals
\[ m_\kappa(h) = \nu_\kappa(h h_\kappa). \]

Normalize \( m_\kappa \) by the condition
\[ m_\kappa(1) = 1. \]

Then \( m_\kappa \) is the invariant measure for the Markov process with transition operator \( P_\kappa \). Denoting by \( m_\kappa \) the corresponding invariant measure on \( \Omega \) we can rewrite (C.8) as
\[ \lambda_\kappa''(a) = m_\kappa(a). \tag{C.10} \]

Next we compute \( \lambda_\kappa''(a) \). Fix a \( a_0 \) and let
\[ \tilde{P}_\kappa h = E_\Phi (e^{\tilde{a}_\kappa h(\Phi_1)}). \]

where
\[ \tilde{a}_\kappa = \kappa (a - m_{a_0}(a)) + \ln h_{a_0} - \ln h_{a_0}(G(\Phi, g)) - \lambda_{a_0} + \kappa_0 m_{a_0}(a). \]

Then the leading eigenvalue of \( \tilde{P}_\kappa \) is
\[ \tilde{\lambda}_\kappa = \lambda_\kappa - (\kappa - \kappa_0) m_{a_0}(a) - \lambda_{a_0} \]
and so \( \tilde{\lambda}_\kappa''(a_0) = \lambda_\kappa''(a_0) \). Let \( \tilde{h}_\kappa \) be the leading eigenvector of \( \tilde{P}_\kappa \) and \( \tilde{\mu}_\kappa \) be the leading eigenvalue of the adjoint operator. Then we have
\[ \tilde{P}_\kappa'' h + 2 \tilde{P}_\kappa' \tilde{h}_\kappa + \tilde{P}_\kappa \tilde{h}_\kappa'' = e^{\tilde{\lambda}_\kappa} \left( \tilde{\lambda}_\kappa \tilde{h}_\kappa \right) + e^{\tilde{\lambda}_\kappa'} \tilde{\mu}_\kappa \tilde{h}_\kappa + 2 e^{\tilde{\lambda}_\kappa'} \tilde{\mu}_\kappa' \tilde{h}_\kappa' + e^{\tilde{\lambda}_\kappa''}. \]

Applying \( \tilde{v}_\kappa \) to both sides and using that
\[ \tilde{\lambda}_{a_0} = 0, \quad \tilde{\lambda}'_{a_0} = 0, \quad \tilde{h}_{a_0} = 1, \quad \tilde{v}_{a_0} = m_{a_0} \]
we get
\[ \lambda_{a_0}'' = m_{a_0}(\hat{a}^2) + 2 m_{a_0}(\hat{a}(\Phi_0, g_0) h_{a_0}(\Phi_1)) \]
where \( \hat{a} = a - m_{a_0}(a) \). Applying the same argument to \( P_\kappa^n \), which has the leading eigenvalue \( e^{n \lambda_\kappa} \) we get
\[ n \lambda_{a_0}' = m_{a_0} \left( \left( \sum_{j=0}^{n-1} \hat{a}(\Phi_j, g_j) \right)^2 \right) + 2 m_{a_0} \left( \left( \sum_{j=0}^{n-1} \hat{a}(\Phi_j, g_j) \right) h_{a_0}(\Phi_n) \right). \]

Since the Markov process with transition operator \( \tilde{P}_\kappa \) has a spectral gap the measure \( m_\kappa \) is ergodic and hence
\[ \frac{1}{n} \sum_{j=0}^{n-1} \hat{a}(\Phi_j, g_j) \to 0 \]
almost surely. Therefore
\[
\lambda''_0 = \lim_{n \to \infty} \frac{1}{n} m_{s_0} \left( \left( \sum_{j=0}^{n-1} \tilde{a}(\Phi_1, g_j) \right)^2 \right).
\] (C.11)

Since the rhs of the last expression is positive we conclude that \( \lambda_x \) is convex. We now show following the argument of theorem 12 of [2] that \( \lambda_x \) is actually strictly convex. Consider the following operator on \( C^\eta(M_1 \times M_2) \)
\[
(\hat{P}_\kappa h)(\Phi_1, g_0) = \tilde{E}_{\Phi_1, g_0} h(\Phi_1, g_1).
\]
Denote \( \Gamma = (1 - \hat{P}_\kappa)^{-1} \hat{a} = \sum_{j=0}^\infty \hat{P}_\kappa^j \hat{a} \). Then a direct computation shows that the rhs of (C.11) equals to \( m_\kappa (\Gamma^2 - (\hat{P}_\kappa \Gamma)^2) \). Hence if \( \lambda''_{s_0} = 0 \) then we have
\[
m_\kappa((\hat{P}_\kappa \Gamma)^2) = m_{s_0}(\Gamma^2).
\]
Since \( m_\kappa \) is stationary for \( \hat{P}_\kappa \) this implies that
\[
m_{s_0}((\hat{P}_\kappa \Gamma)^2) = m_{s_0}(\hat{P}_\kappa (\Gamma^2)).
\]
Now Jensen inequality tells us that \( \Gamma(\Phi, g) \) is actually independent of \( g \), \( \Gamma = \Gamma(\Phi) \). Then
\[
\hat{a}(\Phi, g) = \Gamma(\Phi) - (\hat{P}_\kappa \Gamma)(\Phi, g) = \Gamma(\Phi) - \Gamma(\Phi g)
\]
contradicting (C.2) (as well as (C.3)). This proves that \( \lambda_x \) is strictly convex.

**Proof.** To prove part (a) we iterate (C.9) to get
\[
\tilde{E}_\Phi \left( e^{\tilde{h}_n} \right) = 1,
\]
where
\[
\tilde{y}_n = \xi y_n - n\lambda_x + \ln h(\Phi_0) - \ln h(\Phi_n).
\]
Hence by Markov inequality
\[
P_\Phi(y_n > n\alpha) \leq Ce^{-\pi (\alpha - \lambda_x)} \leq Ce^{\pi (\alpha^* - \alpha)}
\]
where the last inequality uses that \( \lambda_x < \kappa \alpha^* \). This proves part (a).

To prove part (b) suppose that \( \kappa \) is such that \( m_\kappa(a) = \alpha \). Note that \( m_0(a) = \tilde{E}(a) \) so if \( \alpha > \tilde{E}(a) \) then \( \kappa \) is strictly positive.

Let \( \tilde{E}_\kappa \) denote the expectation with respect to the Markov process with generator \( \tilde{P}_\kappa \).

**Lemma C.3.** If (C.3) has no solutions then there exists a function \( \varphi(\Phi) \) such that
\[
\tilde{E}_\kappa \left( 1_{y_n - n\alpha \in I} H(t^n \tilde{\omega}) \right) \sqrt{n} \to \text{Leb}(I) \varphi(\Phi) \nu(H).
\] (C.12)
The proof of this lemma is given in section C.4.

Now take \( I = [t, t + \epsilon] \) then the rhs equals \( e^{\varphi(\Phi)} \nu(H) \) while the lhs equals
\[
\sqrt{n} h_\kappa(\Phi) \tilde{E}_\Phi \left( e^{\pi (\alpha - \lambda_x)} H(t^n \tilde{\omega}) \frac{h_\kappa(\Phi_n)}{h_\kappa(\Phi_0)} e^{\gamma t} \right) (1 + o_{\epsilon \to 0}(1)).
\]
Dividing \( J \) into the segments of length \( \epsilon \ll 1 \) we obtain part (b). Part (c) follows from part (b) and the Markov property. Finally observe that
\[
\gamma'(\alpha) = \alpha \kappa - \lambda_x
\]
where \( \kappa \) satisfies \( \lambda''_x = \alpha \). Thus
\[
\frac{\partial \gamma}{\partial \alpha} = (\alpha - \lambda_x) \frac{\partial \kappa}{\partial \alpha} + \kappa = \kappa.
\]
This proves (C.4). \( \Box \)
C.4. Local limit theorem

Consider
\[ \tilde{P}_{\kappa, u}(h)(\Phi) = E_{\Phi} \left( e^{\tilde{a}_\kappa - i u(a_\kappa - a)} h(\Phi_1) \right). \]

Then
\[ \tilde{P}_{n, \kappa, u}(h)(\Phi) = E_{\Phi} \left( \exp \left[ \sum_{j=0}^{n-1} (\tilde{a}_\kappa - i u(a_\kappa - a))(\tilde{\Phi}_j, g_j) \right] h(\Phi_n) \right). \]

Since
\[ \left| \sum_{j=0}^{n-1} (\tilde{a}_\kappa - i u(a_\kappa - a))(\tilde{\Phi}_j, g_j) - (\tilde{a}_\kappa - i u(a_\kappa - a))(\tilde{\Phi}_j, g_j) \right| \leq C \sum_{j=0}^{n-1} d^n(\tilde{\Phi}_j, \tilde{\Phi}_j) \leq \tilde{C} d^n(\tilde{\Phi}_0, \tilde{\Phi}_0) \]
we get
\[ \| P_{\kappa, u}^n h \|_{C^1} \leq \left[ \| h \|_{C^0} + C(u) \theta^n \| h \|_{C^1} \right]. \] (C.13)

Using the theory of Doeblin–Fortet operators ([22]) we conclude that for each \( \theta' > \theta \) the spectrum of \( P_{\kappa, u} \) outside the disc of radius \( \theta' \) consists of a finite number of eigenvalues with absolute values at most 1. We claim that in fact there are no eigenvalues of absolute value 1. Indeed let \( \tilde{a}^{ui} \) be such an eigenvalue and \( h \) be the corresponding eigenfunction. Then
\[ E_{\Phi} (e^{\tilde{a}^{ui} + i u(\Phi, g)} h(\Phi)) = e^{i \tilde{a}^{ui}} h(\Phi). \] (C.14)

Let \( \Phi^* = \arg \max |h| \). Without loss of generality we can assume that \( |h(\Phi^*)| = 1 \). Now (C.9) implies that (C.14) is only possible if \( |h(G(\Phi^*))| = 1 \) with probability 1. Iterating we see that for all \( n \)
\[ |h(G_n \ldots G_1(\Phi^*))| = 1. \] (C.15)

We claim that this implies that
\[ |h(\Phi)| \equiv 1 \] (C.16)
on the support of \( \tilde{\mu} \). Indeed if \( |h| < 1 - \varepsilon \) on a relatively open subset \( U \) of \( \text{supp}(\tilde{\mu}) \) then there would exist \( \Phi \in M_1 \) and \( n_k \to \infty \) and such that \( G_{n_k} \ldots G_1(\Phi) \in U \) with positive probability. Since \( G_n \ldots G_1 \) contracts with speed \( \theta^n \) for large \( k \) we would have \( |h(G_{n_k} \ldots G_1(\Phi^*))| < 1 - \varepsilon/2 \) with positive probability, contradicting (C.15). Now (C.16) and (C.9) show that
\[ e^{i u(\Phi, g)} h(G(\Phi, g)) = e^{i \tilde{a}^{ui}} h(\Phi) \]
which contradicts (C.3).

We are now ready to prove lemma C.3. Since the lhs of (C.12) is monotone function of \( H \) it suffices to prove the result for a dense set of functions. In particular we may assume that \( H \) depends only on finitely many coordinates
\[ H = H(\omega_{k-1}, \ldots, \omega_0, \ldots, \omega_{k-1}). \]

Then
\[ \sqrt{n} \tilde{\mu}_\Phi \left( 1_{y_n \in \Delta} H(\tau^n \omega) \right) = \sqrt{n} \int E_{\Phi} \left( 1_{y_n \in \Delta} h(\Phi_{n-k}, Z) \right) dP_k(Z) \]
where $h(\Phi, Z) = \mathbb{E}_{\Phi \sim \Phi}(\hat{H}(\tau^{-k} \omega) | y_k = Z)$ and $P_k(z)$ is the distribution function of $y_k$. Observe that for each $Z$ the rhs has the same form as the lhs of (C.12) except that $n$ is replaced to $n - k$ and $h$ depends only on one coordinate. Therefore, it suffices to prove (C.12) in the case where $H = h(\Phi_0)$.

Let $\Gamma_\theta(y) = \frac{1}{\delta}((1 - \cos(\delta y)) / \delta^2) e^{i \theta y}$. Then $\hat{\Gamma}_\theta(u) = (1 - (|u| / \delta))_+$ and $\hat{\Gamma}_\theta(u) = \hat{\Gamma}_0(u + \theta)$. By section 2.5 of [7] it suffices to show that for each $\theta, \delta$ we have

$$\sqrt{n} \mathbb{E}_{\Phi}^n (\Gamma_n(y_n - \alpha n) h(\Phi_n)) \to \phi(\omega) \nu(h) \int \Gamma_n(y) \, dy. \quad (C.17)$$

We have the following inversion formula:

$$\mathbb{P}^n_{\Phi} (\Gamma_\theta(y_n - \alpha n) h(\Phi_n)) = \frac{1}{2\pi} \int_M \hat{\Gamma}_\theta(u) \hat{P}^n_{k,u}(h) \, du$$

where $M$ is such that $\hat{\Gamma}_\theta(u) = 0$ outside $[-M, M]$.

Next, take a small $\epsilon_0$. Then for $|u| < \epsilon_0$ we have the decomposition

$$\hat{P}_{k,u}(h) = \lambda_{k,u} m_{k,u}(h) h_{k,u} + R_{k,u},$$

where

$$R_{k,u}(h_{k,u}) = 0, \quad m_{k,u}(R_{k,u} h) = 0 \quad \text{and} \quad \|R_{k,u}\| \leq K \bar{\theta}^n \text{ for some } \bar{\theta} < 1.$$

It follows that

$$\sqrt{n} \int_{|u| < \epsilon_0} \hat{\Gamma}_\theta(u) \hat{P}^n_{k,u}(h) \, du = \frac{\sqrt{n}}{2\pi} \int_{-\epsilon_0}^{\epsilon_0} \hat{\Gamma}_\theta(u) \lambda^n_{k,u} m_{k,u}(h) h_{k,u} \, du + O(\bar{\theta}^n).$$

Next, letting $u = t / \sqrt{n}$ we can rewrite above integral as

$$\frac{1}{2\pi} \int_{-\epsilon_0 \sqrt{n}}^{\epsilon_0 \sqrt{n}} \Gamma(t / \sqrt{n}) \lambda^n_{k,t / \sqrt{n}} m_{k,t / \sqrt{n}}(h) h_{k,t / \sqrt{n}} \, dt. \quad (C.18)$$

The computations of the previous section give $\lambda_{k,0} = 1$, $\lambda_{k,0}' = 0$ so that $\lambda^n_{k,t / \sqrt{n}} \to e^{-\sigma^2 t^2 / 2}$, where $\sigma^2 = \lambda_{k,t}^{\nu}$, and this convergence is dominated, that is, if $\epsilon_0$ is small enough and $|t| \leq \epsilon_0 \sqrt{n}$ we have $\lambda^n_{k,t / \sqrt{n}} \leq e^{-\sigma^2 t^2 / 4}$. As $u \to 0$ we have $m_{k,u} \to m_k$, $h_{k,u} \to 1$ so that the integral (C.18) converges to $\sigma \hat{\Gamma}_\theta(0)m_k(h)$. On the other hand since for $\epsilon_0 \leq |u| \leq M$ the spectral radius of $\hat{P}_{k,u}$ is strictly less than 1 we have

$$\sqrt{n} \int_{\epsilon_0 \leq |u| < M} \hat{\Gamma}_\theta(u) \hat{P}^n_{k,u}(h) \, du = O(\sqrt{n} \bar{\theta}^n)$$

and (C.17) follows.

Appendix D. Contracting property of the Markov chain $\psi_n$

For two stochastic matrices $\psi'_n$ and $\psi''_n$ define $\psi'_n$ and $\psi''_n$ for $n \geq 1$ using the second formula in (2.1). The sequence $\{P_k, Q_k, R_k\}_{1 \leq k \leq n}$ used in (2.1) is in both cases the same one.

Our goal is to estimate the norm $\|\psi''_n - \psi'_n\|$. We shall prove the following

**Proposition D.1.** Assume that condition (1.5) holds. Then there are constants $K = K(\epsilon) > 0$, $\theta = \theta(\epsilon) < 1$ such that

$$\|\psi''_n - \psi'_n\| \leq K \bar{\theta}^n \|\psi''_0 - \psi'_0\|.$$
Proof. Consider two walkers $X'(t)$ and $X''(t)$ starting at the same site $(n, i)$ in the layer $L_0$ and moving on $[L_0, L_{n+1}]$ with reflecting boundary conditions $\psi'_n$ and $\psi''_n$ at $L_0$ respectively and absorbing boundary conditions at $L_{n+1}$.

We shall show that there exists a coupling between the walkers such that
\[
\mathbb{P}(X'(\tilde{T}'_{n+1}) \neq X''(\tilde{T}''_{n+1})) \leq K\theta^n \|\psi'_n - \psi''_n\|.
\] (D.1)
where $\tilde{T}'_{n+1}$ and $\tilde{T}''_{n+1}$ are the hitting times of $L_{n+1}$ for $X'$ and $X''$ respectively. The statement then follows since, according to the probabilistic meaning of $\psi'_n$ and $\psi''_n$,
\[
\psi'_n(i, j) = \mathbb{P}(X'(\tilde{T}'_{n+1}) = j), \quad \psi''_n(i, j) = \mathbb{P}(X''(\tilde{T}''_{n+1}) = j)
\]
and therefore
\[
\sum_{j=1}^m |\psi'_n(i, j) - \psi''_n(i, j)| \leq 2\mathbb{P}(X'(\tilde{T}'_{n+1}) \neq X''(\tilde{T}''_{n+1})).
\]

The coupling is constructed as follows.

1. The walkers walk together until they either reach $L_0$ for the first time or reach $L_{n+1}$ without visiting $L_0$ (after that they stop). Note that the trajectories which miss $L_0$ do not contribute to the left hand side of (D.1).

2. Let $t_0$ be the first time the walkers reach $L_0$ and $t'_1 - 1, t''_1 - 1$ be the last time $X'$ and $X''$ reach $L_0$ before reaching $L_{n+1}$. Between $t_0$ and $t'_1$ and between $t_0$ and $t''_1$ the walkers move independently.

   Note that after time $t_1$ the walkers again move in the same environment $\{P_k, Q_k, R_k\}_{1 \leq k \leq n}$ but conditioned on reaching $L_{n+1}$ before $L_0$.

3. If $X'(t_1(X')) = X''(t_1(X''))$ then they move together until they reach $L_{n+1}$.

4. Denote by $t'_k, t''_k$ the hitting times of $L_j$ for $X'$ and $X''$, where $l \geq 2$. If $X'(t_1(X')) \neq X''(t_1(X''))$ then they continue to move independently until they reach $L_k$ with the minimal $k$ such that $X'(t'_k) = X''(t''_k)$.

   Note that the only trajectories that contribute to the left hand side of (D.1) are those for which $X'(t'_k) \neq X''(t''_k)$ for all $k \in [2, n+1]$.

We shall use the following estimate

Lemma D.2. Let $X'(t'_1) = (1, J'), X''(t''_1) = (1, J'')$. There exists a constant $C$ such that
\[
d(J', J'') \leq C\|\psi'_n - \psi''_n\|
\]
where $d$ denotes the variational distance between the corresponding distributions.

If at step $k$ of the procedure described above the walkers are uncoupled then condition (1.5) guarantees that the probability that they become coupled at $L_{k+1}$ is at least $\varepsilon m$. Thus the probability that the walkers were uncoupled at time $t_1$ and remain uncoupled until step $n + 1$ is less than
\[
C d(J', J'') \theta^n
\]
and the result follows from lemma D.2. \qed

Proof of lemma D.2. Denote
\[
B(j, k) = \mathbb{P}(X(t_1) = (1, k) \mid X(t_1 + 1) = (1, j)), \quad a(j) = \mathbb{P}(\tilde{T}'_{n+1} < \tilde{T}_1 \mid X(t_0 + 1) = (1, j)),
\]
\[
\Gamma(j, k) = P(X \text{ returns to } (1, k) \text{ after visiting } L_2 \text{ but before } \tilde{T}'_{n+1}) \mid X(t_1 + 1) = (1, j))
\]
Also note that due to ellipticity we have
\[ \min_j a_j \geq \varepsilon \max_j a_j. \]

Let \( M \) be the matrix with \( M(j, k) = a(j) \).

With this notation our goal is to establish Lipshitz dependence of \( B \) on the boundary condition \( \psi \). We have
\[ B = M + Q\psi B + RB + \Gamma B \]
that is
\[ B = (I - Q\psi - R - \Gamma)^{-1} M. \]

Therefore if \( B' \) and \( B'' \) correspond to different boundary conditions \( \psi' \) and \( \psi'' \) respectively then
\[ B' - B'' = (I - Q\psi' - R - \Gamma)^{-1} Q(\psi' - \psi'')(I - Q\psi'' - R - \Gamma)^{-1} M. \]

The estimate we need is a consequence of two inequalities below.
\begin{align*}
|||(\psi' - \psi'')(I - Q\psi'' - R - \Gamma)^{-1}|| &\leq C||\psi' - \psi''|| \quad (D.2) \\
||(I - Q\psi' - R - \Gamma)^{-1}|| &\leq C||M||^{-1} \quad (D.3)
\end{align*}

To prove (D.2) let \( U = Q\psi'' + R + \Gamma \). Given a probability vector \( p \) let \( \pi'_k = p\psi' U^k \), \( \pi''_k = p\psi'' U^k \). Due to ellipticity
\[ ||\pi'_k - \pi''_k|| \leq C(1 - \varepsilon)^k ||\pi'_0 - \pi''_0|| \leq \tilde{C}(1 - \varepsilon)^k ||\psi' - \psi''|| \]
proving (D.2).

To prove (D.3) let \( U = Q\psi' + R + \Gamma \). We have
\[ \sum_k U(j, k) = 1 - a_j. \]
Hence
\[ ||U|| = \max_j \sum_k U(j, k) = 1 - \min_j a_j \leq 1 - \varepsilon \max_j a_j = 1 - \varepsilon ||M||. \]

Accordingly \( ||U^k|| \leq (1 - \varepsilon ||M||)^k \) so that \( ||(1 - U)^{-1}|| \leq (\varepsilon ||M||)^{-1} \) proving (D.3).

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