# Random walks on quasi one dimensional lattices: Large deviations and fluctuation theorems 

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#### Abstract

Several stochastic processes modeling molecular motors on a linear track are given by random walks (not necessarily Markovian) on quasi 1d lattices and share a common regenerative structure. Analyzing this abstract common structure, we derive information on the large fluctuations of the stochastic process by proving large deviation principles for the first-passage times and for the position. We focus our attention on the Gallavotti-Cohen-type symmetry of the position rate function (fluctuation theorem), showing its equivalence with the independence of suitable random variables. In the special case of Markov random walks, we show that this symmetry is universal only inside a suitable class of quasi 1 d lattices.


Résumé. Nous considérons différents processus stochastiques modélisant des moteurs moléculaires : il s'agit de marches aléatoires, non nécessairement markoviennes, le long d'un rail linéaire, presque un réseau unidimensionnel, qui partagent une même structure de régénération. En analysant cette structure abstraite commune nous contrôlons les grandes déviations du processus stochastique, nous établissons des principes de grandes déviations pour les temps de premier passage et pour la variable de position. Nous nous concentrons sur les symétries de type Gallavotti-Cohen de la fonction de taux positionnelle (théorème de fluctuations), en montrant son équivalence avec l'indépendance de certaines variables aléatoires. Dans le cas particulier des marches aléatoires markoviennes, nous montrons que cette symétrie n'est universelle qu'au sein d'une classe particulière de réseaux presque unidimensionnels.

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## 1. Introduction

Molecular motors are proteins working as nanomachines: they usually convert chemical energy coming from ATP hydrolysis to produce mechanical work fundamental e.g. for cargo transport inside the cell, cell division, genetic transcription, muscle contraction. When working in a non-cooperative way, two main modelizations have been proposed. In the so called Brownian ratchet model [16] the dynamics of the molecular motor is given by a one-dimensional diffusion in a switching force field. The other paradigm, on which we concentrate here, is given by continuous time random walks ${ }^{1}$ on quasi linear graphs having a periodic structure [11,12,18-20,30]. We call these graphs quasi $1 d$ lattices, since they are obtained by gluing together several copies of a fundamental cell in a linear fashion. The geometric complexity of the fundamental cell corresponds to the conformational transformations of the molecular motor

[^0]in its mechanochemical cycle. It can be an interval with $N$ sites (thus leading to a random walk on $\mathbb{Z}$ with periodic jump rates) or a more complex graph as in parallel-chain models and divided-chain models [4,17].

In [8] we have derived both the law of large numbers and the invariance principle for random walks on quasi 1 d lattices, and discussed computability issues. We focus here on their large deviations and Gallavotti-Cohen-type symmetries (also called fluctuation theorems), which have received much attention inside the non-equilibrium statistical physics of small systems (cf. $[1,10,22,29]$ and references therein). All relevant information concerning the position of the random walk is encoded in an associated random walk on $\mathbb{Z}$ with nearest neighbor jumps, that we call skeleton process. The holding times of this process are typically non-exponential, and are in general correlated with the location of the jumps. For the skeleton process we derive the LDP for the first-passage times as well as for the position (Theorem 1). We also obtain a detailed qualitative analysis of the rate functions of the above LDPs (Theorem 2 and Proposition 5.3). The tools developed in this part are fundamental to investigate the Gallavotti-Cohen symmetry (shortly, GC symmetry) of the form $I(\vartheta)=I(-\vartheta)+c \vartheta$, where $I$ is the LD rate function for the position of the skeleton process, $\vartheta \in \mathbb{R}$ and $c$ is a suitable constant. The above GC symmetry has been derived in [21,22] for Markov random walks on $\mathbb{Z}$ with periodic rates of period 2 . We show that this symmetry cannot hold for a generic Markov random walk on a quasi 1d lattice. Indeed, there exists a class of quasi 1d lattices (called ( $\underline{v}, \bar{v}$ )-minimal) such that the GC symmetry is verified for any choice of the rates, while outside that class the GC symmetry is violated for Lebesgue any choice of the rates (Theorem 5). In particular, the universality of the above GC symmetry suggested in [22] does not hold. In [9] we will continue our analysis on GC symmetries and we will also consider some examples.

We conclude this introduction with some comments on technical aspects. When considering Markov random walks the proof of the position LD principle is simpler, obtained by the Gärtner-Ellis theorem [15] and by generalizing the matrix approach introduced by [21] (cf. Theorem 3). On the other hand it gives no insight on the mechanism leading to the GC symmetry. The results, presented in Section 2, concerning the LD principles for first-passage times and for the position (Theorems 1 and 2) hold also for non-Markov random walks on quasi 1d lattices. In general (cf. Theorems 6 and 7) they hold for stochastic processes $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$obtained as follows. Consider a sequence $\left(w_{i}, \tau_{i}\right)_{i \geq 1}$ of i.i.d. 2d vectors with values in $\mathbb{R} \times(0,+\infty)$. Defining $W_{m}:=\sum_{i=1}^{m} w_{i}$ and $\mathcal{T}_{m}:=\sum_{i=1}^{m} \tau_{i}$ for $m \geq 0$ integer, set $Z_{t}:=W_{\max \left\{m \geq 0: \mathcal{T}_{m} \leq t\right\}}$. The derivation of the LDP for $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$from the large deviation properties of $\left(W_{m}\right)_{m \geq 0}$ and $\left(\mathcal{T}_{m}\right)_{m \geq 0}$ is in general rather delicate. In [7] a LDP is obtained under the condition that the $\tau_{i}$ 's have finite logarithmic moment generating function. This condition is not satisfied when considering Markov random walks on quasi 1d lattices, hence in our case the results of [7], and the similar ones of [28], cannot be applied. In the context of LDPs for processes under random time changes we also mention the new progresses obtained in $[24,26]$. Restricting to the case $w_{i} \in\{-1,1\}$ (which covers the applications to molecular motors), the process $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$becomes a random walk on $\mathbb{Z}$ with generic holding times (not necessarily exponential). Following the main scheme presented in [5] we derive the LDP for the process $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$. We point out some technical issues making our analysis different from [5]: we allow correlations between $w_{i}$ and $\tau_{i}$ (absent in [5]), moreover the minimum in the support of the law of $\tau_{i}$ can be zero or positive (the first case is excluded in [5]). Hence, although we have no random environment (thus of course simplifying the analysis) in our case there is a richer scenario for the possible behavior of the rate functions of the process $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$and of the associated first-passage times, and this behavior has to be investigated and kept in consideration in order to prove LDPs (see Section 5.2).

The theorems concerning the GC symmetry are the most innovative ones from a mathematical viewpoint. Using the above LD analysis, in Theorems 4 and 8 we prove several characterizations of the GC symmetry for $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$, including the fact that it holds if and only if $w_{i}$ and $\tau_{i}$ are independent. Using the above characterizations, we study the GC symmetry for Markov random walks (Theorem 5). The validity of the GC symmetry for Markov random walks on ( $\underline{v}, \bar{v}$ )-minimal 1d lattices is derived by introducing a special path transformation and comparing the original paths with the transformed ones. On the other hand, the proof of the almost everywhere breaking of the GC symmetry outside the class of $(\underline{v}, \bar{v})$-minimal quasi 1d lattices is based on complex analysis methods.

## 2. Random walks on quasi 1d lattices

We start by defining quasi 1d lattices. Consider first a finite oriented graph $G=(V, E), V$ being the set of vertices and $E$ being the set of oriented edges, $E \subset\{(v, w): v \neq w$ in $V\}$. We fix in $V$ two vertices $\underline{v}, \bar{v}$. We assume that the oriented graph $G$ is connected, i.e. for any $v, w \in V$ there is an oriented path in $G$ from $v$ to $w$. Then the quasi 1d lattice $\mathcal{G}$ associated to the triple ( $G, \underline{v}, \bar{v}$ ) is the oriented graph obtained by gluing together countable copies of $G$ such


Fig. 1. The graph $G=(V, E)$ with vertices $\underline{v}, \bar{v}$ (left) and the associated quasi 1 d lattice $\mathcal{G}=(\mathcal{V}, \mathcal{E})$ (right).
that the point $\bar{v}$ of one copy is identified with the point $\underline{v}$ of the next copy. To give a formal definition, we define $\mathcal{G}$ as $\mathcal{G}=(\mathcal{V}, \mathcal{E})$ with vertex set $\mathcal{V}$ and edge set $\mathcal{E}$ as follows (see Figure 1):

$$
\mathcal{V}:=\left\{v_{n}:=(v, n) \in(V \backslash\{\bar{v}\}) \times \mathbb{Z}\right\}, \quad \mathcal{E}:=\mathcal{E}_{1} \cup \mathcal{E}_{2} \cup \mathcal{E}_{3},
$$

where

$$
\begin{aligned}
& \mathcal{E}_{1}:=\left\{\left(v_{n}, w_{n}\right):(v, w) \in E, n \in \mathbb{Z}\right\}, \\
& \mathcal{E}_{2}:=\bigcup_{n \in \mathbb{Z}}\left\{\left(v_{n}, \underline{v}_{n+1}\right):(v, \bar{v}) \in E\right\}, \\
& \mathcal{E}_{3}:=\bigcup_{n \in \mathbb{Z}}\left\{\left(\underline{v}_{n+1}, v_{n}\right):(\bar{v}, v) \in E\right\} .
\end{aligned}
$$

To simplify notation we set

$$
n_{*}:=\underline{v}_{n}, \quad n \in \mathbb{Z}
$$

On the graph $\mathcal{G}$ we define the shift $\mathcal{T}: \mathcal{V} \rightarrow \mathcal{V}$ as $\mathcal{T}\left(v_{n}\right)=v_{n+1}$. Note that the graph $\mathcal{G}$ is left invariant by the action of $\mathcal{T}$. We can now define the class of stochastic processes on quasi 1 d lattices we are interested in:

Definition 2.1. Given a quasi 1 d lattice $\mathcal{G}$ associated to the triple $(G, \underline{v}, \bar{v})$, we consider a stochastic process $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$ with paths in the Skohorod space $D\left(\mathbb{R}_{+} ; \mathcal{V}\right)$ starting at any site $n_{*}\left(\right.$ we denote by $\mathbb{P}_{n_{*}}$ the associated law on $\left.D\left(\mathbb{R}_{+} ; \mathcal{V}\right)\right)$ and fulfilling the following properties:
(i) for each $n \in \mathbb{Z}, \mathbb{P}_{n_{*}}$-a.s., jumps are possible only along the edges in $\mathcal{E}$,
(ii) for each $n \in \mathbb{Z}$, when $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$is sampled with law $\mathbb{P}_{n_{*}}$ then the law of $\left(\mathcal{T}\left(X_{t}\right)\right)_{t \in \mathbb{R}_{+}}$equals $\mathbb{P}_{(n+1)_{*}}$,
(iii) defining $S$ as the random time

$$
\begin{equation*}
S:=\inf \left\{t \geq 0: X_{t} \in\left\{-1_{*}, 1_{*}\right\}\right\}, \tag{1}
\end{equation*}
$$

it holds $\mathbb{E}_{0_{*}}(S)<\infty$ where $\mathbb{E}_{0_{*}}(\cdot)$ denotes the expectation w.r.t. $\mathbb{P}_{0_{*}}$,
(iv) under $\mathbb{P}_{0_{*}}\left(\cdot \mid X_{S}= \pm 1_{*}\right)$ the random path $\left(X_{S+t}\right)_{t \in \mathbb{R}_{+}}$is independent from $\left(X_{t}\right)_{t \in[0, S]}$ and has law $\mathbb{P}_{ \pm 1_{*}}$.

In the applications, typically $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$is a Markov random walk:
Lemma 2.2. Let $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$be a Markov random walk with state space $\mathcal{V}$ and with positive jump rates $r(x, y),(x, y) \in$ $\mathcal{E}$, such that

$$
\begin{equation*}
r(x, y)=r(\mathcal{T} x, \mathcal{T} y) \tag{2}
\end{equation*}
$$

Then the above random walk is well defined for all times $t$ (no explosion takes place), fulfills the properties of Definition 2.1 and moreover $\mathbb{E}_{0_{*}}\left(e^{\lambda S}\right)<+\infty$ for $\lambda>0$ small enough.

The proof of the above lemma is simple and therefore omitted. The finite exponential moments for $\lambda$ small follow from the exponential decay of hitting probabilities for irreducible Markov chains with finite state space. We point out that in the applications another relevant example is given by a random walk $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$on the graph $\mathcal{G}$ with nonexponential holding times (cf. [18]).

We now introduce the fundamental object of our investigation:
Definition 2.3. Given the stochastic process $X$ as in Definition 2.1, the skeleton process $X^{*}=\left(X_{t}^{*}\right)_{t \in \mathbb{R}_{+}}$is defined as $X_{t}^{*}:=\Phi\left(X_{\iota}\right)$ where $\Phi\left(n_{*}\right)=n$ and

$$
\iota:=\sup \left\{s \in[0, t]: X_{t}=n_{*} \text { for some } n \in \mathbb{Z}\right\} .
$$

$X_{t}^{*}$ has values in $\mathbb{Z}$ and records the last visited state of the form $n_{*}$ up to time $t$.
In the applications to molecular motors, the process $\left(X_{t}^{*}\right)_{t \in \mathbb{R}_{+}}$contains all the relevant information, indeed it allows to determine the position of the molecular motor up to an error of the same order of the monomer size.

## 3. Main results for random walks on quasi 1d lattices

Let $S$ be the random time defined in (1). As proved in [8], one can easily obtained a strong law of large numbers for the skeleton process since $\mathbb{E}_{0_{*}}(S)<+\infty$ (cf. Theorems 1 and 2 in [8]):

$$
\begin{equation*}
\lim _{t \rightarrow \infty} \frac{X_{t}^{*}}{t}=\frac{\mathbb{P}_{0_{*}}\left(X_{S}=1_{*}\right)-\mathbb{P}_{0_{*}}\left(X_{S}=-1_{*}\right)}{\mathbb{E}_{0_{*}}(S)}=: v, \quad \mathbb{P}_{0_{*}} \text {-a.s. } \tag{3}
\end{equation*}
$$

In [8] we study also the Gaussian fluctuations of the skeleton process, proving an invariance principle if $\mathbb{E}_{0_{*}}\left(S^{2}\right)<$ $+\infty$. We concentrate here on large deviations.

### 3.1. Large deviations

From now on, in addition to the requirements in Definition 2.1, we assume that

$$
\begin{equation*}
\mathbb{P}_{0_{*}}\left(X_{S}=1_{*}\right)>0 \quad \text { and } \quad \mathbb{P}_{0_{*}}\left(X_{S}=-1_{*}\right)>0, \tag{4}
\end{equation*}
$$

which holds for molecular motors.
Theorem 1. Consider the process $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$starting at $0_{*}$. Call $T_{n}$ the first time the skeleton process hits $n \in \mathbb{Z}$, i.e.

$$
\begin{equation*}
T_{n}:=\inf \left\{t \in \mathbb{R}_{+}: X_{t}^{*}=n\right\} \in[0,+\infty] . \tag{5}
\end{equation*}
$$

Then the following holds:
(i) As $n \rightarrow \pm \infty$ the random variables $T_{n} /|n|$ satisfy a LDP with speed $|n|$ and convex rate function

$$
\begin{equation*}
J_{ \pm}(u):=\sup _{\lambda \in \mathbb{R}}\left\{\lambda u-\log \varphi_{ \pm}(\lambda)\right\}, \quad u \in \mathbb{R} \tag{6}
\end{equation*}
$$

where

$$
\begin{equation*}
\varphi_{ \pm}(\lambda):=\mathbb{E}_{0_{*}}\left(e^{\lambda T_{ \pm 1}} \mathbb{1}\left(T_{ \pm 1}<\infty\right) \in(0,+\infty], \quad \lambda \in \mathbb{R}\right. \tag{7}
\end{equation*}
$$

The rate function $J_{ \pm}$is good ${ }^{2}$ if and only if $\mathbb{P}_{0_{*}}\left(T_{1}<\infty\right) \neq \mathbb{P}_{0_{*}}\left(T_{-1}<\infty\right)$.

[^1](ii) As $t \rightarrow+\infty$, the random variables $X_{t}^{*} / t$ satisfy a LDP with speed $t$ and good and convex rate function I given by
\[

I(\vartheta)= $$
\begin{cases}\vartheta J_{+}(1 / \vartheta) & \text { if } \vartheta>0,  \tag{8}\\ |\vartheta| J_{-}(1 /|\vartheta|) & \text { if } \vartheta<0,\end{cases}
$$
\]

and $I(0)=\lim _{\vartheta \rightarrow 0} I(\vartheta)$.
Theorem 1 is an immediate consequence of Lemma 4.1 and Theorem 6 in Section 4. Since $T_{n}$ can take value $+\infty$, the meaning of the LDP for $T_{n} /|n|$ as $n \rightarrow \pm \infty$ is the following: for each close subset $\mathcal{C} \subset \mathbb{R}$ and each open subset $\mathcal{O} \subset \mathbb{R}$ it holds

$$
\varlimsup_{n \rightarrow \pm \infty} \frac{1}{|n|} \log \mathbb{P}_{0_{*}}\left(\frac{T_{n}}{|n|} \in \mathcal{C}\right) \leq-\inf _{\mathcal{C}} J_{ \pm}, \quad \lim _{n \rightarrow \pm \infty} \frac{1}{|n|} \log \mathbb{P}_{0_{*}}\left(\frac{T_{n}}{|n|} \in \mathcal{O}\right) \geq-\inf _{\mathcal{O}} J_{ \pm} .
$$

We now collect information on the qualitative behavior of the rate function $I$. The qualitative behavior of the rate functions $J_{-}, J_{+}$is described in Proposition 5.3 in Section 5 . Here we concentrate on the rate function $I$ since the large deviations of $X_{t}^{*} / t$ are more relevant in the applications.

Definition 3.1. We define $\alpha_{ \pm}$as the minimum of the support of the law of $T_{ \pm 1}$.
We point out that $\alpha_{ \pm}$is the minimum of the support of the Borel measure $A \rightarrow \mathbb{P}_{0_{*}}\left(S \in A, X_{S}= \pm 1_{*}\right)$ (see Proposition 4.3 in Section 4). Below $1 / \alpha_{ \pm}$is intended to be $+\infty$ if $\alpha_{ \pm}=0$. Note that $\alpha_{ \pm}=0$ in the case of Markov random walks.

Theorem 2. The following holds:
(i) I is infinite outside $\left[-\frac{1}{\alpha_{-}}, \frac{1}{\alpha_{+}}\right]$, I is finite and $C^{1}$ on $\left(-\frac{1}{\alpha_{-}}, \frac{1}{\alpha_{+}}\right)$, moreover it is smooth on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right) \backslash\{0\}$.
(ii) The following holds:

$$
\begin{align*}
& \lim _{\vartheta \backslash \frac{1}{\alpha_{+}}} I(\vartheta)= \begin{cases}+\infty & \text { if } \mathbb{P}_{0_{*}}\left(T_{1}=\alpha_{+}\right)=0, \\
I\left(\frac{1}{\alpha_{+}}\right)<\infty & \text { otherwise. }\end{cases}  \tag{9}\\
& \lim _{\vartheta \searrow-\frac{1}{\alpha_{-}}} I(\vartheta)= \begin{cases}+\infty & \text { if } \mathbb{P}_{0_{*}}\left(T_{-1}=\alpha_{-}\right)=0, \\
I\left(-\frac{1}{\alpha_{-}}\right)<\infty & \text { otherwise. }\end{cases} \tag{10}
\end{align*}
$$

(iii) The derivative of I satisfies $\lim _{\vartheta \nearrow \frac{1}{\alpha_{+}}} I^{\prime}(\vartheta)=+\infty$ and $\lim _{\vartheta \searrow-\frac{1}{\alpha_{-}}} I^{\prime}(\vartheta)=-\infty$.
(iv) I is lower semicontinuous and convex on $\mathbb{R}$, it is strictly convex on $\left(-\frac{1}{\alpha_{-}}, \frac{1}{\alpha_{+}}\right)$.
(v) I has a unique global minimum, which is given by 0 and is attained at $v \in\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$, where $v$ is the asymptotic velocity defined in (3). Moreover I is strictly decreasing on $\left(-1 / \alpha_{-}, v\right)$ and is strictly increasing on ( $v, 1 / \alpha_{+}$).

Theorem 2 is an immediate consequence of Lemma 4.1 and Theorem 7 in Section 4.
When the process $X$ is a Markov random walk with periodic rates (i.e. satisfying (2)), the derivation of the large deviation principle is simpler. In this case given $x \in \mathcal{V}$ we set $r(x):=\sum_{y:(x, y) \in \mathcal{E}} r(x, y)$ and, given $v \neq w$ in $V \backslash\{\bar{v}\}$, we set (using the convention that $r(y, z)=0$ if $(y, z) \notin \mathcal{E}$ )

$$
r(v):=r\left(v_{n}\right), \quad r_{-}(w, v):=r\left(w_{n-1}, v_{n}\right), \quad r_{0}(w, v):=r(w, v), \quad r_{+}(w, v):=r\left(w_{n+1}, v_{n}\right) .
$$

The above definition is well posed due to (2). Finally, given $\lambda \in \mathbb{R}$ we consider the finite matrix $\mathcal{A}(\lambda)$, with entries parameterized by $(V \backslash\{\bar{v}\}) \times(V \backslash\{\bar{v}\})$, defined as :

$$
\mathcal{A}_{v, w}(\lambda):= \begin{cases}-r(v) & \text { if } v=w,  \tag{11}\\ e^{\lambda} r_{-}(w, v)+r_{0}(w, v)+e^{-\lambda} r_{+}(w, v) & \text { if } v \neq w .\end{cases}
$$

By applying the Gärtner-Ellis Theorem we will derive the following result, where $\mathcal{R}(\gamma)$ denotes the real part of $\gamma \in \mathbb{C}$ :

Theorem 3. Suppose $X$ is a Markov random walk on the quasi 1 d lattice $\mathcal{G}=(\mathcal{V}, \mathcal{E})$ with transition rates satisfying (2). Then, as $t \rightarrow+\infty$, the random variables $X_{t}^{*} / t$ satisfy a large deviation principle with speed $t$ and convex and good rate function $I(\vartheta)$ given by

$$
I(\vartheta)=\sup _{\lambda \in \mathbb{R}}\{\vartheta \lambda-\Lambda(\lambda)\}, \quad \vartheta \in \mathbb{R},
$$

where $\Lambda(\lambda)$ is the finite value $\Lambda(\lambda):=\max \{\Re(\gamma): \gamma$ eigenvalue of $\mathcal{A}(\lambda)\}$ and the matrix $\mathcal{A}(\lambda)$ is defined in (11).

### 3.2. Fluctuation theorems (Gallavotti-Cohen type symmetry)

Theorem 4. The following facts are equivalent:
(i) For some $c \in \mathbb{R}$ the Gallavotti-Cohen type symmetry ${ }^{3} I(\vartheta)=I(-\vartheta)+c \vartheta$ holds for all $\vartheta \in \mathbb{R}$;
(ii) The random variables $X_{S}$ and $S$ are independent.

Moreover, when (i), (ii) hold it must be $c=\log \frac{\mathbb{P}_{0_{*}}\left(X_{S}=-1_{*}\right)}{\mathbb{P}_{0_{*}}\left(X_{S}=1_{*}\right)}$.
Theorem 4 is an immediate consequence of Lemma 4.1 and Theorem 8 in Section 4.
We continue our investigation of the Gallavotti-Cohen type symmetry (shortly, GC symmetry) as in Theorem 4 restricting now to Markov random walks $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$on quasi 1d lattices. Recall that we write $r(x, y),(x, y) \in \mathcal{E}$, for the positive jump rates of the Markov random walk and assume the periodicity (2) to hold. We restrict our discussion to the case of fundamental graphs $G=(V, E)$ such that

$$
\begin{equation*}
(x, y) \in E \quad \text { if and only if } \quad(y, x) \in E, \tag{12}
\end{equation*}
$$

which is the standard setting in the investigation of GC symmetry for Markov chains [23].
In what follows, given an edge $(u, v) \in E$ in the fundamental graph $G=(V, E)$, we define

$$
\begin{equation*}
r(u, v)=r(\pi(u), \pi(v)), \tag{13}
\end{equation*}
$$

where $\pi$ is the map $V \rightarrow \mathcal{V}$ defined as $\pi(u)=u_{0}$ if $u \neq \bar{v}$ and $\pi(\bar{v})=\underline{v}_{1}$. Note that, fixed positive numbers $a(e)$, $e \in E$, it is univocally determined a Markov random walk on $\mathcal{G}$ whose rates satisfy (2) and such that $r(e)=a(e)$ for all $e \in E$. We call it the Markov random walk induced by $a(e), e \in E$.

We introduce a special class of graphs $G$ which includes in particular trees. Recall that $G$ has connected support when disregarding the orientation of the edges.

Definition 3.2. We say that the graph $G=(V, E)$ is $(\underline{v}, \bar{v})$-minimal if it satisfies (12) and moreover there is a unique path $\gamma_{*}=\left(z_{0}, z_{1}, \ldots, z_{n}\right)$ such that (i) $z_{0}=\underline{v}$, (ii) $z_{n}=\bar{v}$, (iii) $\left(z_{i}, z_{i+1}\right) \in E$ and (iv) the points $z_{0}, \ldots, z_{n}$ are all distinct.

Note that, given a generic fundamental graph $G=(V, E)$, there exists at least one path $\gamma=\left(z_{0}, z_{1}, \ldots, z_{n}\right)$ satisfying the above properties (i), $\ldots$, (iv). Indeed, since $G$ is connected, there exists a path from $\underline{v}$ to $\bar{v}$. Take such a path and prune iteratively the loops. Since each time a loop is pruned away the length of the path decreases, after a finite number of prunes one gets a minimal path satisfying the above properties (i), ..., (iv).

Now suppose that $G=(V, E)$ is $(\underline{v}, \bar{v})$-minimal and take two points $z_{i} \neq z_{j}$ (the $z_{k}$ 's are as in the Definition 3.2). Then it cannot exist a path from $z_{i}$ to $z_{j}$ whose intermediate points are in $V \backslash\left\{z_{0}, z_{1}, \ldots, z_{n}\right\}$. In particular, the graph $G$ must be as in Figure 2 (due to property (12) we only draw the support of $G$, disregarding orientation). More precisely, the graph is given by the linear chain $\gamma_{*}$ of Definition 3.2, to which one attaches some subgraphs, in such a way that each attached subgraph has exactly one point in common with $\gamma_{*}$.

[^2]

Fig. 2. A $(\underline{v}, \bar{v})$-minimal graph $G=(V, E)$.

Theorem 5. Suppose that $X$ is a Markov random walk and that $G=(V, E)$ is a graph satisfying (12). If $G$ is $(\underline{v}, \bar{v})$ minimal, then the random variables $S$ and $X_{S}^{*}$ are independent, and in particular the large deviation rate function $I$ associated to the skeleton process $X^{*}$ satisfies the Gallavotti-Cohen type symmetry

$$
\begin{equation*}
I(\vartheta)=I(-\vartheta)-\Delta \vartheta, \quad \forall \vartheta \in \mathbb{R}, \tag{14}
\end{equation*}
$$

where

$$
\begin{equation*}
\Delta=\log \frac{r\left(z_{0}, z_{1}\right) r\left(z_{1}, z_{2}\right) \cdots r\left(z_{n-1}, z_{n}\right)}{r\left(z_{1}, z_{0}\right) r\left(z_{2}, z_{1}\right) \cdots r\left(z_{n}, z_{n-1}\right)} \tag{15}
\end{equation*}
$$

and $\gamma_{*}=\left(z_{0}, z_{1}, z_{2}, \ldots, z_{n-1}, z_{n}\right)$ is the path in Definition 3.2.
Vice versa, if $G$ is not $(\underline{v}, \bar{v})$-minimal then the vectors $(r(e): e \in E) \in(0,+\infty)^{E}$ for which the induced random walk on $\mathcal{G}$ satisfies (14) for some constant $\Delta$ (depending on the numbers $r(e), e \in E$ ) has zero Lebesgue measure in $(0,+\infty)^{E}$.

The proof of the above theorem is given in Section 9.
Remark 3.3. In Section 9 we give a sufficient condition (see Criterion 1) assuring (14) for some constant $\Delta$ and for fixed rates (also outside the class of ( $\underline{v}, \bar{v}$ )-minimal graphs).

## 4. Random time changes of cumulative processes

As already mentioned, the results presented in Section 3 hold in a more general context that we now describe. Consider a sequence $\left(w_{i}, \tau_{i}\right)_{i \geq 1}$ of i.i.d. 2 d vectors with values in $\mathbb{R} \times(0,+\infty)$. For each integer $m \geq 1$ we define

$$
\begin{equation*}
W_{m}:=w_{1}+w_{2}+\cdots+w_{m}, \quad \mathcal{T}_{m}:=\tau_{1}+\tau_{2}+\cdots+\tau_{m} \tag{16}
\end{equation*}
$$

We set $W_{0}=\mathcal{T}_{0}=0$. Note that $\lim _{m \rightarrow \infty} \mathcal{T}_{m}=+\infty$ a.s. As a consequence, we can univocally define a.s. a random process $\{\nu(t)\}_{t \in \mathbb{R}_{+}}$with values in $\{0,1,2,3, \ldots\}$ such that

$$
\begin{equation*}
\mathcal{T}_{\nu(t)} \leq t<\mathcal{T}_{\nu(t)+1}, \quad t \geq 0 \tag{17}
\end{equation*}
$$

Note that $v(t)=\max \left\{m \in \mathbb{N}: \mathcal{T}_{m} \leq t\right\}$. Finally, we define the process $Z:[0, \infty) \rightarrow \mathbb{R}$ as

$$
\begin{equation*}
Z_{t}:=W_{\nu(t)} \tag{18}
\end{equation*}
$$

Note that $Z_{0}=0$. The resulting process $Z=\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$is therefore obtained from the cumulative process $\left(W_{m}\right)_{m \geq 0}$ by a random time change.

Due to Definitions 2.1 and 2.3 the skeleton process $X^{*}$ is indeed a special case of process $Z$ (recall the definition of the random time $S$ given in (1)):

Lemma 4.1. Consider a sequence $\left(w_{i}, \tau_{i}\right)_{i \geq 1}$ of i.i.d. vectors, with the same law of the random vector $\left(X_{S}^{*}, S\right) \in$ $\{-1,1\} \times(0,+\infty)$ when the random walk $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$starts at $0_{*}$. We define $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$as the stochastic process built from $\left(w_{i}, \tau_{i}\right)_{i \geq 1}$ according to (18). Then $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$has the same law of $\left(X_{t}^{*}\right)_{t \in \mathbb{R}_{+}}$with $X_{0}^{*}=0$.

The proof of the above lemma is very simple and therefore omitted. We recall also the LLN discussed in [8, Appendix A]:

Proposition 4.2. If $\mathbb{E}\left(\tau_{i}\right)<\infty$, then almost surely $\lim _{t \rightarrow \infty} \frac{Z_{t}}{t}=v:=\frac{\mathbb{E}\left(w_{i}\right)}{\mathbb{E}\left(\tau_{i}\right)}$.
We now state our main results for $\left(Z_{t}\right)_{t \in \mathbb{R}_{+}}$:
Theorem 6 [LDP]. Assume: (A1) $w_{i} \in\{-1,1\}$ a.s., (A2) $\mathbb{P}\left(w_{1}=+1\right)>0$ and $\mathbb{P}\left(w_{1}=-1\right)>0$. Set $T_{n}:=\inf \{t \in$ $\left.\mathbb{R}_{+}: Z_{t}=n\right\} \in[0,+\infty)$. Define $J_{ \pm}$and $\varphi_{ \pm}$as in (6) and (7). Then the following holds:
(i) As $n \rightarrow \pm \infty$ the random variables $T_{n} /|n|$ satisfy a LDP with speed $|n|$ and convex rate function $J_{ \pm}$. The rate function $J_{ \pm}$is good if and only if $\mathbb{P}\left(T_{1}<\infty\right) \neq \mathbb{P}\left(T_{-1}<\infty\right)$.
(ii) As $t \rightarrow+\infty$, the random variables $Z_{t} / t$ satisfy a LDP with speed $t$ and good and convex rate function I given by

$$
I(\vartheta)= \begin{cases}\vartheta J_{+}(1 / \vartheta) & \text { if } \vartheta>0,  \tag{19}\\ |\vartheta| J_{-}(1 /|\vartheta|) & \text { if } \vartheta<0,\end{cases}
$$

and $I(0)=\lim _{\vartheta \rightarrow 0} I(\vartheta)$.
The proof of the above result is given in Sections 5 and 6 .
We introduce the functions $f_{ \pm}$on $\mathbb{R}$ as

$$
\begin{equation*}
f_{ \pm}(\lambda):=\mathbb{E}\left(e^{\lambda \tau_{1}} \mathbb{1}\left(w_{1}= \pm 1\right)\right) \in(0,+\infty) . \tag{20}
\end{equation*}
$$

Note that $f_{ \pm}(\lambda)>0$ under assumption (A2). We call $\alpha_{ \pm}$the minimum value in the support of the law of $T_{ \pm 1}$.
Proposition 4.3. Under assumptions (A1) and (A2) of the previous theorem the following holds:
(i) The function $\varphi_{ \pm}(\lambda)$ satisfies

$$
\begin{equation*}
\varphi_{ \pm}(\lambda)=\frac{1-\sqrt{1-4 f_{-}(\lambda) f_{+}(\lambda)}}{2 f_{\mp}(\lambda)} \tag{21}
\end{equation*}
$$

for $\lambda \leq \lambda_{c}$, where $\lambda_{c} \in[0,+\infty)$ is the unique value of $\lambda$ such that $f_{-}(\lambda) f_{+}(\lambda)=1 / 4$, while $\varphi_{ \pm}(\lambda)=+\infty$ for $\lambda>\lambda_{c}$.
(ii) Consider the measure $\mu_{ \pm}$on $[0,+\infty)$ such that $\mu_{ \pm}(A)=\mathbb{P}\left(\tau_{1} \in A, w_{1}= \pm 1\right)$ for any Borel $A \subset \mathbb{R}$. Then $\alpha_{ \pm}$is the minimum value in the support of $\mu_{ \pm}$. Moreover $\mathbb{P}\left(T_{ \pm 1}=\alpha_{ \pm}\right)=\mathbb{P}\left(\tau_{1}=\alpha_{ \pm}, w_{1}= \pm 1\right)$.

The proof is given at the beginning of Section 5.2.
The qualitative behavior of the rate function $I(\vartheta)$ is described by the following theorem (for the qualitative behavior of $J_{ \pm}$see Proposition 5.3):

Theorem 7. Theorem 2 remains valid in the present more general context, with $v$ defined as in Proposition 4.2.
We conclude with a result on the presence of a Gallavotti-Cohen type symmetry in the rate function $I$ (the proof is given in Section 7):

Theorem 8. The following facts are equivalent:
(i) There exists a constant $c \in \mathbb{R}$ such that the Gallavotti-Cohen type symmetry

$$
\begin{equation*}
I(\vartheta)=I(-\vartheta)+c \vartheta \tag{22}
\end{equation*}
$$

holds for all $\vartheta \in \mathbb{R}$;
(ii) fixed $i$, the random variables $w_{i}, \tau_{i}$ are independent;
(iii) the functions $\varphi_{+}(\lambda)$ and $\varphi_{-}(\lambda)$ are proportional where finite: $\exists C>0$ such that $\varphi_{+}(\lambda)=C \varphi_{-}(\lambda)$ for all $\lambda \leq \lambda_{c}$.

Moreover, if we let $p:=\mathbb{P}\left(w_{i}=1\right)$ and $q:=\mathbb{P}\left(w_{i}=-1\right)$ (with $p, q>0$ by Assumption (A2)), then $C=p / q$ and $c=\log (q / p)=-\log C$.

We end this section with some comments on the existing literature. Considering [26] with $\mathcal{X}, x_{i}, \tau: \mathcal{X} \rightarrow(0, \infty)$ given by $\mathcal{X}:=\{-1,1\} \times(0,+\infty), x_{i}=\left(w_{i}, \tau_{i}\right)$ and $\tau(x, t)=t$, one gets the large deviation principle of the random probability measure $\pi_{t}=\frac{1}{t} \int_{[0, t)} \delta_{\left(w_{\nu(s)+1}, \tau_{\nu(s)+1}\right)} d s \in \mathcal{P}(\mathcal{X})$ as $t \rightarrow \infty$ (recall (17)). Let us restrict to the case $w_{i}=$ $\pm 1$. If we consider the function $f: \mathcal{X} \rightarrow(0,+\infty)$ defined as $f(x, t)=x / t$, then the expectation $\pi_{t}(f)$ satisfies $\left|\pi_{t}(f)-\frac{Z_{t}}{t}\right| \leq \frac{1}{t} \frac{t-\mathcal{T}_{v(t)}}{\tau_{\nu(t)+1}} \leq \frac{1}{t}$, hence $\pi_{t}(f)$ and $\frac{Z_{t}}{t}$ must satisfy the same large deviation principle. On the other hand, the map $\mathcal{P}(\mathcal{X}) \ni \pi \rightarrow \pi(f) \in(0,+\infty)$ is not a bounded continuous map, and one cannot derive immediately by contraction the LDP of $\pi_{t}(f)$ (and therefore of $Z_{t} / t$ ) from the LDP of the empirical measure $\pi_{t}$ stated in [26, Theorem 1.3, Theorem 1.6]. Of course, one could try to apply some extended version of the contraction principle (cf. [6]) to follow this route. On the other hand, supposing to succeed in such a direction, one would end up with a rate function for $Z_{t} / t$ characterized as the solution of a variational problem on an infinite dimensional space, instead of a finite dimensional space as in Theorem 6.

## 5. Proof of Theorem 6(i) and Theorem 7

In this section we prove item (i) of Theorem 6 and we study the behavior of the functions $I, J_{ \pm}$defined in Theorem 6. In particular, we prove Theorem 7 at the end of this section.

### 5.1. Proof of item (i) in Theorem 6

For $n \geq 1$ the random variable $T_{n}$ has the same law of $\sum_{k=1}^{n} \hat{\tau}_{n}$, where $\hat{\tau}_{n}$ 's are i.i.d. random variables taking value in $(0,+\infty)$, distributed as $T_{1}$. The thesis can then be obtained from Cramér Theorem. We give the proof in the case $n \rightarrow \infty$. Call $\alpha:=\mathbb{P}\left(T_{1}<\infty\right)$ and note that $P\left(T_{n}<\infty\right)=\alpha^{n}$. Then for each subset $\mathcal{A} \subset \mathbb{R}$ we can write $\mathbb{P}\left(T_{n} / n \in\right.$ $\mathcal{A})=\alpha^{n} \mathbb{P}\left(T_{n} / n \in \mathcal{A} \mid T_{n}<\infty\right)$. Now we observe that, conditioning on the event $T_{n}<\infty, T_{n}$ can be represented as $\sum_{k=1}^{n} \hat{\tau}_{n}^{\prime}$, where the real random variables $\tau_{n}^{\prime}$ are i.i.d. and distributed as $T_{1}$ conditioned to be finite. In conclusion $\mathbb{P}\left(T_{n} / n \in \mathcal{A} \mid T_{n}<\infty\right)=P\left(\frac{1}{n} \sum_{k=1}^{n} \hat{\tau}_{n}^{\prime} \in \mathcal{A}\right)$. At this point one only need to apply Cramér Theorem for i.i.d. real random variables (cf. [6, Theorem 2.2.3]) observing that the moment generating function of $\hat{\tau}_{n}^{\prime}$ is $\varphi_{+} / \alpha$. The fact that $J_{ \pm}$is good if and only if $\mathbb{P}\left(T_{1}<\infty\right) \neq \mathbb{P}\left(T_{-1}<\infty\right)$ is proved in the next subsection (see Remark 5.4 below).

### 5.2. Qualitative study of the functions $J_{ \pm}(\vartheta), I(\vartheta)$

In this subsection we first prove some properties of the function $I(\vartheta)$ defined in Theorem 6 by (19) and the identity $I(0)=\lim _{\vartheta \rightarrow 0} I(\vartheta)$. In the next subsection we will indeed prove that $I(\vartheta)$ is the rate function of the LDP for $Z_{t} / t$.

We start by proving Proposition 4.3.
Proof of Proposition 4.3. Let us prove point (i). Recall the definition of the positive functions $f_{ \pm}$given in (20). Distinguishing on the value of $w_{1}$ we can write

$$
\begin{equation*}
T_{1}=\mathbb{1}\left(w_{1}=1\right) \tau_{1}+\mathbb{1}\left(w_{1}=-1\right)\left(\tau_{1}+T_{1}^{\prime}+T_{1}^{\prime \prime}\right) \tag{23}
\end{equation*}
$$

where $T_{1}^{\prime}, T_{1}^{\prime \prime}$ are independent random variables, independent from $w_{1}, \tau_{1}$ and distributed as $T_{1}$ (roughly, $T_{1}^{\prime}$ is the time for $Z$ to go from -1 to 0 and $T_{1}^{\prime \prime}$ is the time for $Z$ to go from 0 to 1 ). The above identity implies that

$$
\begin{equation*}
\varphi_{+}(\lambda)=f_{+}(\lambda)+f_{-}(\lambda) \varphi_{+}(\lambda)^{2} . \tag{24}
\end{equation*}
$$

From this we deduce that $\varphi_{+}(\lambda)<+\infty$ if and only if $f_{+}(\lambda) f_{-}(\lambda) \leq 1 / 4$, and moreover in this case (21) holds. By symmetry, one obtains that the same condition implies $\varphi_{-}(\lambda)<+\infty$. Trivially $f_{ \pm}$is increasing, $\lim _{\lambda \rightarrow-\infty} f_{ \pm}(\lambda)=0$ and $\lim _{\lambda \rightarrow+\infty} f_{ \pm}(\lambda)=+\infty$. Moreover, $f_{ \pm}(\lambda)$ is smooth and strictly increasing on the open set $\left\{f_{ \pm}<+\infty\right\}$. As a consequence there exists a unique value $\lambda_{c} \in \mathbb{R}$ such that $f_{+}(\lambda) f_{-}(\lambda)=1 / 4$ and therefore $f_{+}(\lambda) f_{-}(\lambda) \leq 1 / 4$ for $\lambda \leq \lambda_{c}$. Since trivially $\varphi_{+}(\lambda)<+\infty$ for $\lambda \leq 0$ it must be $\lambda_{c} \geq 0$. This completes the proof of point (i).

We now move to point (ii). To see that $\mathbb{P}\left(\tau_{1}<\alpha_{+}, w_{1}=1\right)=0$ observe that by (23) $\left\{w_{1}=1\right\} \subseteq\left\{T_{1}=\tau_{1}\right\}$ and therefore $\mathbb{P}\left(\tau_{1}<\alpha_{+}, w_{1}=1\right)=\mathbb{P}\left(T_{1}<\alpha_{+}, w_{1}=1\right) \leq \mathbb{P}\left(T_{1}<\alpha_{+}\right)=0$. To get the thesis it remains to show that $\forall \varepsilon>0 \mathbb{P}\left(\tau_{1} \in\left[\alpha_{+}, \alpha_{+}+\varepsilon\right), w_{1}=1\right)>0$. Assume the contrary. Then there exists $\hat{\varepsilon}>0$ such that $\mathbb{P}\left(\tau_{1} \in\left[\alpha_{+}, \alpha_{+}+\right.\right.$ $\left.\hat{\varepsilon}), w_{1}=1\right)=0$. By definition of $\alpha_{+}$, on the other hand, we have $\mathbb{P}\left(T_{1} \in\left[\alpha_{+}, \alpha_{+}+\hat{\varepsilon}\right)\right)>0$. Combining this with the decomposition in (23), we find

$$
\begin{aligned}
0 & <\mathbb{P}\left(T_{1}<\alpha_{+}+\hat{\varepsilon}\right) \\
& =\mathbb{P}\left(\tau_{1}<\alpha_{+}+\hat{\varepsilon}, w_{1}=1\right)+\mathbb{P}\left(\tau_{1}+T_{1}^{\prime}+T_{1}^{\prime \prime}<\alpha_{+}+\hat{\varepsilon}, w_{1}=-1\right) \\
& \leq \mathbb{P}\left(\tau_{1}<\alpha_{+}+\hat{\varepsilon}, w_{1}=-1\right) \mathbb{P}\left(T_{1}<\alpha_{+}+\hat{\varepsilon}\right)^{2} \\
& \leq \mathbb{P}\left(w_{1}=-1\right) \mathbb{P}\left(T_{1}<\alpha_{+}+\hat{\varepsilon}\right)^{2}
\end{aligned}
$$

Dividing both sides by the positive quantity $\mathbb{P}\left(T_{1}<\alpha_{+}+\hat{\varepsilon}\right)$ and recalling that by (A2) $\mathbb{P}\left(w_{1}=-1\right)<1$, we get the contradiction and this concludes the proof.

We now focus on the behavior $\log \varphi_{ \pm}$. Recall the definition of $\lambda_{c}, \alpha_{ \pm}$given in Proposition 4.3.

## Lemma 5.1. The following holds:

(i) $\log \varphi_{ \pm}$is strictly increasing and continuous on $\left(-\infty, \lambda_{c}\right]$, convex and smooth on $\left(-\infty, \lambda_{c}\right)$ and moreover $\lim _{\lambda \rightarrow-\infty} \log \varphi_{ \pm}(\lambda)=-\infty ;$
(ii) the second derivative $\left(\log \varphi_{ \pm}\right)^{\prime \prime}$ is strictly positive on $\left(-\infty, \lambda_{c}\right)$ (in particular $\left(\log \varphi_{+}(\lambda)\right)^{\prime}$ is strictly increasing on $\left.\left(-\infty, \lambda_{c}\right)\right)$ and

$$
\begin{align*}
& \lim _{\lambda \rightarrow-\infty}\left(\log \varphi_{ \pm}(\lambda)\right)^{\prime}=\alpha_{ \pm}  \tag{25}\\
& \lim _{\lambda \nearrow \lambda_{c}}\left(\log \varphi_{ \pm}(\lambda)\right)^{\prime}=+\infty \tag{26}
\end{align*}
$$

Proof. The proof of point (i) is rather standard (see Lemma 2.2.5 in [6], the fact that $\log \varphi_{ \pm}$is strictly increasing on $\left(-\infty, \lambda_{c}\right]$ and convex follows also from point (ii)).

We prove point (ii) restricting to $\varphi_{+}$without loss of generality. Note that, for $\lambda<\lambda_{c}$,

$$
\left(\log \varphi_{+}(\lambda)\right)^{\prime \prime}=\frac{\mathbb{E}\left(T_{1}^{2} e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right)}{\mathbb{E}\left(e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right)}-\frac{\mathbb{E}\left(T_{1} e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right)^{2}}{\mathbb{E}\left(e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right)^{2}}=\operatorname{Var}_{\mathbb{Q}}\left(T_{1}\right)
$$

where $\mathbb{Q}$ is the probability defined as $\mathbb{Q}(A)=\mathbb{E}\left(\mathbb{1}(A) e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right) / \mathbb{E}\left(e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right)$. Since $T_{1}$ is non constant $\mathbb{Q}$-a.s. by assumption (A1), we conclude that $\left(\log \varphi_{+}(\lambda)\right)^{\prime \prime}>0$ for $\lambda<\lambda_{c}$, hence $\left(\log \varphi_{+}(\lambda)\right)^{\prime}$ on $\left(-\infty, \lambda_{c}\right)$ is strictly increasing.

We first derive (25) in the case $\alpha_{+}=0$. It is convenient to prove the thesis for a generic nonnegative random variable $T_{1}$, nonnecessarily defined as in Theorem 6 . Suppose first that $\mathbb{P}\left(T_{1}=0\right)>0$. Since $\varphi_{+}(\lambda) \geq \mathbb{P}\left(T_{1}=0\right)$, while $\lim _{\lambda \rightarrow-\infty} \varphi^{\prime}(\lambda)=\lim _{\lambda \rightarrow-\infty} \mathbb{E}\left(T_{1} e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<\infty\right)\right)=0$ by the monotone convergence theorem, we get (25).

We now consider the case $\mathbb{P}\left(T_{1}=0\right)=0$, thus implying $\mathbb{P}\left(T_{1} \in(0, \varepsilon)\right)>0$ for all $\varepsilon>0$. We fix any $c>0$ and take $\lambda<-1 / c$. By this choice it holds $\sup _{x \geq c} x e^{\lambda x}=c e^{\lambda c}$. Moreover we fix $c_{1}, c_{2}: 0<c_{1}<c_{2}<c$ such that $\mathbb{P}\left(c_{1} \leq T_{1} \leq\right.$ $\left.c_{2}\right)>0$. Define:

$$
\begin{aligned}
& e_{1}(\lambda):=\mathbb{E}\left(T_{1} e^{\lambda T_{1}} \mathbb{1}\left(c \leq T_{1}<\infty\right)\right) \leq c e^{\lambda c} \\
& e_{2}(\lambda):=\mathbb{E}\left(T_{1} e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<c\right)\right) \geq c_{1} e^{\lambda c_{2}} \mathbb{P}\left(c_{1} \leq T_{1} \leq c_{2}\right)
\end{aligned}
$$

$$
\begin{aligned}
& e_{3}(\lambda):=\mathbb{E}\left(e^{\lambda T_{1}} \mathbb{1}\left(c \leq T_{1}<\infty\right)\right) \leq e^{\lambda c}, \\
& e_{4}(\lambda):=\mathbb{E}\left(e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<c\right)\right) \geq e^{\lambda c / 2} \mathbb{P}\left(T_{1}<c / 2\right)>0 .
\end{aligned}
$$

By the previous bounds we have $\lim _{\lambda \rightarrow-\infty} e_{1}(\lambda) / e_{2}(\lambda)=0$ and $\lim _{\lambda \rightarrow-\infty} e_{3}(\lambda) / e_{4}(\lambda)=0$. In conclusion

$$
\begin{aligned}
0 & \leq \lim _{\lambda \rightarrow-\infty}\left(\log \varphi_{+}(\lambda)\right)^{\prime}=\lim _{\lambda \rightarrow-\infty} \frac{\varphi_{+}^{\prime}(\lambda)}{\varphi_{+}(\lambda)}=\lim _{\lambda \rightarrow-\infty} \frac{e_{2}(\lambda)\left(1+e_{1}(\lambda) / e_{2}(\lambda)\right)}{e_{4}(\lambda)\left(1+e_{3}(\lambda) / e_{4}(\lambda)\right)} \\
& =\lim _{\lambda \rightarrow-\infty} \frac{e_{2}(\lambda)}{e_{4}(\lambda)}=\lim _{\lambda \rightarrow-\infty} \frac{\mathbb{E}\left(T_{1} e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<c\right)\right)}{\mathbb{E}\left(e^{\lambda T_{1}} \mathbb{1}\left(T_{1}<c\right)\right)} \leq c .
\end{aligned}
$$

Since $c>0$ is arbitrary we get (25).
To complete the proof of (25) it remains to discuss the case $\alpha_{+}>0$. To this aim note that 0 is the minimum in the support of the law of $\hat{T}_{1}:=T_{1}-\alpha_{+}$. Hence, by what just proven, it holds $\lim _{\lambda \rightarrow-\infty}\left(\log \hat{\varphi}_{+}(\lambda)\right)^{\prime}=0$, where $\hat{\varphi}_{+}(\lambda):=\mathbb{E}\left(e^{\lambda \hat{T}_{1}} \mathbb{1}\left(\hat{T}_{1}<\infty\right)\right)$. Since $\varphi_{+}(\lambda)=e^{\lambda \alpha_{+}} \hat{\varphi}_{+}(\lambda)$, we get (25).

To conclude the proof of point (ii) we need to justify (26). Since by point (i) $\log \varphi_{+}$is smooth and convex on $\left(-\infty, \lambda_{c}\right)$, the derivative $\left(\log \varphi_{+}(\lambda)\right)^{\prime}=\varphi_{+}^{\prime}(\lambda) / \varphi_{+}(\lambda)$ is increasing and therefore the limit in (26) exists. Moreover, since $\lim _{\lambda / \lambda_{c}} \varphi_{+}(\lambda)=\varphi_{+}\left(\lambda_{c}\right)<\infty$, we only need to show that $\lim _{\lambda / \lambda_{c}} \varphi_{+}^{\prime}(\lambda)=+\infty$. To this aim recall (24). Differentiating such identity for $\lambda<\lambda_{c}$ (note that everything is finite and smooth) we get

$$
\begin{equation*}
\left(1-2 f_{-}(\lambda) \varphi_{+}(\lambda)\right) \varphi_{+}^{\prime}(\lambda)=f_{+}^{\prime}(\lambda)+f_{-}^{\prime}(\lambda) \varphi_{+}(\lambda)^{2} . \tag{27}
\end{equation*}
$$

By the monotone convergence theorem we get that $f_{-}(\lambda), \varphi_{+}(\lambda)$ and the derivative $f_{ \pm}^{\prime}(\lambda)=\mathbb{E}\left(\tau_{1} e^{\lambda \tau_{1}} \mathbb{1}\left(w_{1}= \pm 1\right)\right)$ converge to $f_{-}\left(\lambda_{c}\right), \varphi_{+}\left(\lambda_{c}\right)$ and $f_{ \pm}^{\prime}\left(\lambda_{c}\right)$ respectively as $\lambda \nearrow \lambda_{c}$. Observing that

$$
\begin{equation*}
\varphi_{+}\left(\lambda_{c}\right)=\frac{1}{2 f_{-}\left(\lambda_{c}\right)} \tag{28}
\end{equation*}
$$

due to (21) and the identity $f_{+}\left(\lambda_{c}\right) f_{-}\left(\lambda_{c}\right)=1 / 4$, we get that $1-2 f_{-}(\lambda) \varphi_{+}(\lambda)$ converges to zero as $\lambda \nearrow \lambda_{c}$. On the other hand, as $\lambda \nearrow \lambda_{c}$ the r.h.s. of (27) converges to its value at $\lambda_{c}$, which is nonzero. The limit (26) is therefore the only possibility as $\lambda \nearrow \lambda_{c}$ in (27).

Recall the definition of the asymptotic velocity $v$ given in (3).
Remark 5.2. Taking the expectation in (23) and in the analogous expression for $T_{-1}$, one concludes that $\mathbb{E}\left(T_{ \pm 1}\right)<$ $+\infty$ implies that $\mathbb{E}\left(T_{ \pm 1}\right) \mathbb{E}\left(w_{1}\right)= \pm \mathbb{E}\left(\tau_{1}\right)$. Since $\mathbb{E}\left(\tau_{1}\right) \neq 0$, we conclude that if $\mathbb{E}\left(T_{ \pm 1}\right)<\infty$ then $\mathbb{E}\left(w_{1}\right) \neq 0$ and $\mathbb{E}\left(T_{ \pm 1}\right)= \pm \mathbb{E}\left(\tau_{1}\right) / \mathbb{E}\left(w_{1}\right)= \pm 1 / v$.

From Lemma 5.1 we deduce the qualitative behavior of the rate function $J_{ \pm}(\vartheta):=\sup _{\lambda \in \mathbb{R}}\left\{\lambda \vartheta-\log \varphi_{ \pm}(\lambda)\right\}$ :
Proposition 5.3. The following holds:
(i) $J_{ \pm}$is lower semicontinuous, convex and takes values in $[0,+\infty)$.
(ii) $J_{ \pm}$is finite on $\left(\alpha_{ \pm},+\infty\right)$ and infinite on $\left(-\infty, \alpha_{ \pm}\right)$.
(iii) $J_{ \pm}$is smooth on $\left(\alpha_{ \pm},+\infty\right)$ and the derivative $J_{ \pm}^{\prime}$ satisfies $\lim _{\vartheta \rightarrow+\infty} J_{ \pm}^{\prime}(\vartheta)=\lambda_{c}$. In particular, $\lim _{\vartheta \rightarrow+\infty} J_{ \pm}(\vartheta)=+\infty$ if $\lambda_{c}>0$.
(iv) If $\lambda_{c}=0$ then $J_{ \pm}$is strictly decreasing on $\left(\alpha_{ \pm},+\infty\right)$. If $\lambda_{c}>0$ then there exist $\vartheta_{c}^{ \pm} \in\left(\alpha_{ \pm},+\infty\right)$ such that $J_{ \pm}$is strictly decreasing on $\left(\alpha_{ \pm}, \vartheta_{c}^{ \pm}\right)$, strictly increasing on $\left(\vartheta_{c}^{ \pm},+\infty\right)$. Moreover:

- $v=0$ if and only if $\lambda_{c}=0$,
- if $v>0$ then $\vartheta_{c}^{+}=1 / v$ and $J_{+}\left(\vartheta_{c}^{+}\right)=0$,
- if $v<0$ then $\vartheta_{c}^{-}=-1 / v$ and $J_{-}\left(\vartheta_{c}^{-}\right)=0$.
(v) The value $J_{ \pm}\left(\alpha_{ \pm}\right)$admits the following characterization

$$
J_{ \pm}\left(\alpha_{ \pm}\right)=\lim _{\vartheta \searrow \alpha_{ \pm}} J_{ \pm}(\vartheta)= \begin{cases}+\infty & \text { if } \mathbb{P}\left(T_{ \pm 1}=\alpha_{ \pm}\right)=0,  \tag{29}\\ <\infty & \text { otherwise. }\end{cases}
$$

Remark 5.4. Due to the above result, $J_{ \pm}$is a good rate function (i.e. $\left\{J_{ \pm} \leq u\right\}$ is compact for all $u \in \mathbb{R}$ ) if and only if $\lambda_{c}>0$.

Proof of Proposition 5.3. Without loss of generality we prove the above statements only for $J_{+}$.
The proof of point (i) is standard (cf. [6, Chapter 2]) and we omit it. We now prove point (ii). The fact that $J_{+}(\vartheta)=\infty$ for $\vartheta \leq 0$ follows from Lemma 5.1(i). We now show that if $\alpha_{+}>0$ then $J_{+}(\vartheta)=\infty$ also for $\vartheta \in\left(0, \alpha_{+}\right)$. For such $\vartheta$, by (25) in Lemma 5.1 it holds

$$
\lim _{\lambda \rightarrow-\infty} \frac{\log \varphi_{+}(\lambda)}{\lambda \vartheta}=\lim _{\lambda \rightarrow-\infty} \frac{\left(\log \varphi_{+}\right)^{\prime}(\lambda)}{\vartheta}=\frac{\alpha_{+}}{\vartheta}>1
$$

and therefore

$$
J_{+}(\vartheta) \geq \lim _{\lambda \rightarrow-\infty} \lambda \vartheta\left(1-\frac{\log \varphi_{+}(\lambda)}{\lambda \vartheta}\right)=+\infty .
$$

Take now $\vartheta>\alpha_{+}$. Since by Lemma $5.1\left(\log \varphi_{+}\right)^{\prime}(\lambda)$ is a strictly increasing function which takes values in $\left(\alpha_{+},+\infty\right)$, there exists a unique $\tilde{\lambda}_{+}(\vartheta)$ such that

$$
\begin{equation*}
\vartheta=\left(\log \varphi_{+}\right)^{\prime}\left(\tilde{\lambda}_{+}(\vartheta)\right) . \tag{30}
\end{equation*}
$$

Then $J_{+}(\vartheta)=\vartheta \tilde{\lambda}_{+}(\vartheta)-\log \varphi_{+}\left(\tilde{\lambda}_{+}(\vartheta)\right)$ which is finite. This concludes the proof of (ii).
We now move to point (iii). Observe that, by (30) and Lemma 5.1, $\tilde{\lambda}_{+}$is the inverse function of $\left(\log \varphi_{+}\right)^{\prime}$. By Lemma $5.1\left(\log \varphi_{ \pm}\right)^{\prime}$ is smooth on $\left(-\infty, \lambda_{c}\right)$ and $\left(\log \varphi_{ \pm}\right)^{\prime \prime}>0$ on $\left(-\infty, \lambda_{c}\right)$. Hence, by the implicit function theorem, the function $\tilde{\lambda}_{+}$is smooth on $\left(\alpha_{+},+\infty\right)$ and tending to $\lambda_{c}$ as $\vartheta \rightarrow+\infty$ (see (26)). Hence $J_{+}$is smooth on $\left(\alpha_{+},+\infty\right)$ and

$$
\begin{equation*}
J_{+}^{\prime}(\vartheta)=\vartheta \tilde{\lambda}_{+}^{\prime}(\vartheta)+\tilde{\lambda}_{+}(\vartheta)-\left(\log \varphi_{+}\right)^{\prime}\left(\tilde{\lambda}_{+}(\vartheta)\right) \tilde{\lambda}_{+}^{\prime}(\vartheta)=\tilde{\lambda}_{+}(\vartheta), \tag{31}
\end{equation*}
$$

thus implying that $\lim _{\vartheta \rightarrow+\infty} J_{+}^{\prime}(\vartheta)=\lambda_{c}$. This concludes the proof of (iii).
We now prove point (iv). By Lemma 5.1 $\tilde{\lambda}_{+}$is strictly increasing on $\left(\alpha_{+},+\infty\right), \lim _{\vartheta} \tau_{\alpha_{+}} \tilde{\lambda}_{+}(\vartheta)=-\infty$ and $\lim _{\vartheta \rightarrow \infty} \tilde{\lambda}_{+}(\vartheta)=\lambda_{c}$. If $\lambda_{c}=0$, then $\tilde{\lambda}_{+}$must be negative and from (31) we conclude that $J_{+}$is strictly decreasing on $\left(\alpha_{+},+\infty\right)$. If $\lambda_{c}>0$, then there exists a unique $\vartheta_{c}^{+}$such that $\tilde{\lambda}_{+}\left(\vartheta_{c}^{+}\right)=0, \tilde{\lambda}_{+}$is negative on the left of $\vartheta_{c}^{+}$and is positive on the right of $\vartheta_{c}^{+}$. Hence, by (31) we see that $J_{+}$has a unique minimum at $\vartheta=\vartheta_{c}^{+}$. In this case, from (30) we have

$$
\vartheta_{c}^{+}=\left(\log \varphi_{+}\right)^{\prime}(0)=\frac{\varphi_{+}^{\prime}(0)}{\varphi_{+}(0)}=\frac{\mathbb{E}\left(T_{1} \mathbb{\mathbb { 1 }}\left(T_{1}<\infty\right)\right)}{\mathbb{P}\left(T_{1}<\infty\right)}
$$

If $v>0$, then by the LLN in Proposition 4.2 we get that $T_{1}<\infty$ a.s. and $\vartheta_{c}^{+}=\mathbb{E}\left(T_{1}\right)=1 / v$ (cf. Remark 5.2). Hence, recalling that $\tilde{\lambda}_{+}\left(\vartheta_{c}^{+}\right)=0$,

$$
\inf _{\vartheta \in \mathbb{R}} J_{+}(\vartheta)=J_{+}\left(\vartheta_{c}^{+}\right)=\tilde{\lambda}_{+}\left(\vartheta_{c}^{+}\right) \vartheta_{c}^{+}-\log \varphi_{+}\left(\tilde{\lambda}_{+}\left(\vartheta_{c}^{+}\right)\right)=-\log \mathbb{P}\left(T_{1}<\infty\right)=0 .
$$

The case $v<0$ can be treated similarly. We conclude the proof by showing that $v=0 \Leftrightarrow \lambda_{c}=0$. Trivially, $v=0 \Leftrightarrow$ $\mathbb{P}\left(w_{1}= \pm 1\right)=\frac{1}{2} \Leftrightarrow \mathbb{P}\left(w_{1}=1\right) \mathbb{P}\left(w_{1}=-1\right)=\frac{1}{4}$. The last identity can be rewritten as $f_{+}(0) f_{-}(0)=\frac{1}{4}$, where the function $f_{+}, f_{-}$are defined as in (20). Due to the characterization of $\lambda_{c}$ given after (28), we conclude that the last identity is equivalent to $\lambda_{c}=0$.

To derive point (v) we note that by point (iv) the limit in (29) exists. We first assume $\mathbb{P}\left(T_{1}=\alpha_{+}\right)=0$. Taking $\delta>0$ and $\lambda<0$ we can bound

$$
\varphi_{+}(\lambda) \leq e^{\lambda \alpha_{+}} \mathbb{P}\left(T_{1} \leq \alpha_{+}+\delta\right)+e^{\lambda\left(\alpha_{+}+\delta\right)} \mathbb{P}\left(T_{1}>\alpha_{+}+\delta\right),
$$

thus implying

$$
J_{+}\left(\alpha_{+}\right) \geq \lambda \alpha_{+}-\log \varphi_{+}(\lambda) \geq-\log \left[\mathbb{P}\left(T_{1} \leq \alpha_{+}+\delta\right)+e^{\lambda \delta} \mathbb{P}\left(T_{1}>\alpha_{+}+\delta\right)\right] .
$$

To get that $J_{+}\left(\alpha_{+}\right)=+\infty$ it is enough to take first the limit $\lambda \rightarrow-\infty$ and afterwards the limit $\delta \rightarrow 0$. Since $J_{+}$is also 1.s.c. one has $\lim _{\vartheta} \alpha_{+} J_{+}(\vartheta) \geq J_{+}\left(\alpha_{+}\right)$and therefore one gets (29).

Assume, on the other hand, that $\mathbb{P}\left(T_{1}=\alpha_{+}\right)>0$. The fact that $J_{+}\left(\alpha_{+}\right)<\infty$ follows by the LDP for $T_{n}$ (cf. Section 5.1) and the characterization of $\alpha_{+}$given in Proposition 4.3(ii). Indeed we can bound

$$
\begin{aligned}
-J_{+}\left(\alpha_{+}\right) & \geq \varlimsup_{n \rightarrow \infty} \frac{1}{n} \log \mathbb{P}\left(\frac{T_{n}}{n}=\alpha_{+}\right)=\varlimsup_{n \rightarrow \infty} \frac{1}{n} \log \mathbb{P}\left(\tau_{1}=\alpha_{+}, w_{1}=1\right)^{n} \\
& =\log \mathbb{P}\left(\tau_{1}=\alpha_{+}, w_{1}=1\right)=\log \mathbb{P}\left(T_{1}=\alpha_{+}\right)>-\infty .
\end{aligned}
$$

To see that $J_{+}$is right-continuous at $\alpha_{+}$observe that by lower semicontinuity $J_{+}\left(\alpha_{+}\right) \leq \lim _{\vartheta} \searrow_{\alpha_{+}} J_{+}(\vartheta)$. We claim that $J_{+}\left(\alpha_{+}\right) \geq \lim _{\vartheta \backslash \alpha_{+}} J_{+}(\vartheta)$. Indeed, fixed $\alpha_{0}>\alpha_{+}$, by convexity it holds

$$
J_{+}\left(\alpha_{+}\right) \geq \frac{1}{1-\lambda} J_{+}\left((1-\lambda) \alpha_{+}+\lambda \alpha_{0}\right)-\frac{\lambda}{1-\lambda} J_{+}\left(\alpha_{0}\right) .
$$

The claim then follows from the monotonicity of $J_{+}$on the right of $\alpha_{+}$. Combining the last observations we get $\lim _{\vartheta \searrow \alpha_{+}} J_{+}(\vartheta)=J_{+}\left(\alpha_{+}\right)<\infty$ and this concludes the proof of point (v).

We now move to the study of the function $I(\vartheta)$ defined on $\mathbb{R} \backslash\{0\}$ as

$$
I(\vartheta)= \begin{cases}I_{+}(\vartheta):=\sup _{\lambda \in \mathbb{R}}\left\{\lambda-\vartheta \log \varphi_{+}(\lambda)\right\}, & \vartheta>0,  \tag{32}\\ I_{-}(\vartheta):=\sup _{\lambda \in \mathbb{R}}\left\{\lambda+\vartheta \log \varphi_{-}(\lambda)\right\}, & \vartheta<0 .\end{cases}
$$

Lemma 5.5. It holds

$$
\begin{align*}
& \lim _{\vartheta \nearrow 0} I_{-}(\vartheta)=\lim _{\vartheta \searrow 0} I_{+}(\vartheta)=\lambda_{c},  \tag{33}\\
& \lim _{\vartheta>0} I_{-}^{\prime}(\vartheta)=\lim _{\vartheta \searrow 0} I_{+}^{\prime}(\vartheta) . \tag{34}
\end{align*}
$$

In particular, the definition of $I(\vartheta)$ in Theorem 6 is well posed and $I(0)=\lambda_{c}$. Moreover, $I$ is finite and $C^{1}$ on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$, and it is smooth on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right) \backslash\{0\}$.

Proof. For any $\vartheta>0$ we have $I(\vartheta)=\sup _{\lambda \leq \lambda_{c}}\left\{\lambda-\vartheta \log \varphi_{+}(\lambda)\right\}$ since $\varphi_{+}(\lambda)=+\infty$ if $\lambda>\lambda_{c}$. Moreover, always by Lemma 5.1, for $0<\vartheta<1 / \alpha_{+}$the above supremum is attained at the unique value $\lambda_{+}(\vartheta)<\lambda_{c}$ such that

$$
\begin{equation*}
\left(\log \varphi_{+}\right)^{\prime}\left(\lambda_{+}(\vartheta)\right)=1 / \vartheta \tag{35}
\end{equation*}
$$

thus implying that $\lambda_{+}(\vartheta)$ is a strictly decreasing function and $\lim _{\vartheta \searrow 0} \lambda_{+}(\vartheta)=\lambda_{c}$ (due to Lemma 5.1). In particular, $I(\vartheta)=\lambda_{+}(\vartheta)-\vartheta \log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)$ is finite on $\left(0,1 / \alpha_{+}\right)$and moreover

$$
\lim _{\vartheta \searrow 0} I(\vartheta)=\lim _{\vartheta \searrow 0}\left\{\lambda_{+}(\vartheta)-\vartheta \log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)\right\}=\lambda_{c}
$$

since $\lim _{\lambda / \lambda_{c}} \log \varphi_{+}(\lambda)=\log \varphi_{+}\left(\lambda_{c}\right)$ which is finite due to (28). This concludes the proof of (33) for $I_{+}$. By similar arguments one gets that, given $\vartheta \in\left(-1 / \alpha_{-}, 0\right)$ there is a unique value $\lambda_{-}(\vartheta)$ solving the equation

$$
\begin{equation*}
\left(\log \varphi_{-}\right)^{\prime}\left(\lambda_{-}(\vartheta)\right)=-1 / \vartheta . \tag{36}
\end{equation*}
$$

The function $\lambda_{-}$is strictly increasing on $\left(-1 / \alpha_{-}, 0\right)$ where it holds $I(\vartheta)=\lambda_{-}(\vartheta)+\vartheta \log \varphi_{-}\left(\lambda_{-}(\vartheta)\right)$. As above one gets that $\lim _{\vartheta ~}^{70} I_{-}(\vartheta)=\lambda_{c}$, hence (33). Note that (33) implies that $I$ is well defined in Theorem 6 and that $I(0)=\lambda_{c}$. By the previous results we conclude also that $I$ is finite on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$.

Let us now prove (34) and that $I$ is $C^{1}$ on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right) \backslash\{0\}$. By the implicit function theorem and Lemma 5.1, the function $\left(0,1 / \alpha_{+}\right) \ni \vartheta \rightarrow \lambda_{+}(\vartheta) \in\left(-\infty, \lambda_{c}\right)$ is smooth. In particular, using (35), $I_{+}$is smooth on $\left(0,1 / \alpha_{+}\right)$ where it holds

$$
\begin{align*}
I_{+}^{\prime}(\vartheta) & =\frac{\mathrm{d}}{\mathrm{~d} \vartheta}\left(\lambda_{+}(\vartheta)-\vartheta \log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)\right) \\
& =\lambda_{+}^{\prime}(\vartheta)-\log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)-\vartheta \cdot\left(\log \varphi_{+}\right)^{\prime}\left(\lambda_{+}(\vartheta)\right) \cdot \lambda_{+}^{\prime}(\vartheta) \\
& =-\log \varphi_{+}\left(\lambda_{+}(\vartheta)\right) \tag{37}
\end{align*}
$$

Hence, $\lim _{\vartheta \searrow 0} I_{+}^{\prime}(\vartheta)=-\log \varphi_{+}\left(\lambda_{c}\right)$. By similar arguments and definitions we get that $I_{-}$is smooth on $\left(-1 / \alpha_{-}, 0\right)$ where it holds

$$
\lim _{\vartheta \nearrow 0} I_{-}^{\prime}(\vartheta)=\lim _{\vartheta \nearrow 0} \frac{\mathrm{~d}}{\mathrm{~d} \vartheta}\left(\lambda_{-}(\vartheta)+\vartheta \log \varphi_{-}\left(\lambda_{-}(\vartheta)\right)\right)=\log \varphi_{-}\left(\lambda_{c}\right) .
$$

To conclude the proof of (34) it remains to show that $\log \varphi_{-}\left(\lambda_{c}\right)=-\log \varphi_{+}\left(\lambda_{c}\right)$. To this aim we observe that

$$
\log \varphi_{-}\left(\lambda_{c}\right)+\log \varphi_{+}\left(\lambda_{c}\right)=\log \left[\varphi_{+}\left(\lambda_{c}\right) \varphi_{-}\left(\lambda_{c}\right)\right]=\log \left(\frac{1}{4 f_{+}\left(\lambda_{c}\right) f_{-}\left(\lambda_{c}\right)}\right)=0
$$

due to (28), its analogous version for $\varphi_{-}\left(\lambda_{c}\right)$ and since, by definition, $\lambda_{c}$ is the unique solution of $4 f_{-}(\lambda) f_{+}(\lambda)=1$. This concludes the proof of (34) and that $I$ is smooth on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right) \backslash\{0\}$. Due to (34) one easily gets that $I$ is differentiable at 0 and $I^{\prime}(0)$ equals the limits in (34). This implies that $I$ is $C^{1}$ on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$.

Combining Lemmas 5.1, 5.5 and Proposition 5.3 we are finally able to prove Theorem 7 and therefore also Theorem 2 due to Lemma 4.1.

Proof of Theorem 7. Below the labeling of items is as in Theorem 2.
The fact that $I$ is finite and $C^{1}$ on $\left(-\frac{1}{\alpha_{-}}, \frac{1}{\alpha_{+}}\right)$and infinite outside $\left[-\frac{1}{\alpha_{-}}, \frac{1}{\alpha_{+}}\right]$follows from (19) and Proposition 5.3. This proves item (i).

To prove item (ii) note that if $\mathbb{P}\left(T_{1}=\alpha_{+}\right)>0$ then $\alpha_{+}>0$. Hence, by (29) we get

$$
\lim _{\vartheta \nearrow \frac{1}{\alpha_{+}}} I(\vartheta)=\lim _{\vartheta \nearrow \frac{1}{\alpha_{+}}} \vartheta J_{+}\left(\frac{1}{\vartheta}\right)=\frac{J_{+}\left(\alpha_{+}\right)}{\alpha_{+}}<\infty
$$

and the last term equals $I\left(1 / \alpha_{+}\right)$by definition of $I$. If, on the other hand, $\mathbb{P}\left(T_{1}=\alpha_{+}\right)=0$, then by (29) we get

$$
\lim _{\vartheta \nearrow \frac{1}{\alpha_{+}}} I(\vartheta)=\lim _{u \searrow \alpha_{+}} \frac{J_{+}(u)}{u}=+\infty
$$

The correspondent statements for $\vartheta \searrow-1 / \alpha_{-}$are obtained in the same way, and this concludes the proof of item (ii).
To see (iii) recall that $I^{\prime}(\vartheta)=-\log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)$ for $\lambda \in\left(0,1 / \alpha_{+}\right)$(see (37)). Observe now that $\lim _{\vartheta \nearrow \frac{1}{\alpha_{+}}} \lambda_{+}(\vartheta)=$ $-\infty$ (due to Lemma 5.1(ii) and (35)). This implies that

$$
\lim _{\vartheta \nearrow \frac{1}{\alpha_{+}}}\left(-\log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)\right)=-\log \varphi_{+}(-\infty)=+\infty .
$$

Similarly one sees that

$$
\lim _{\vartheta \searrow-\frac{1}{\alpha_{-}}} I^{\prime}(\vartheta)=\lim _{\vartheta \searrow-\frac{1}{\alpha_{-}}} \log \varphi_{-}\left(\lambda_{-}(\vartheta)\right)=\log \varphi_{-}(-\infty)=-\infty .
$$

We now consider item (iv). We observe that $I_{+}, I_{-}$are l.s.c. because they can be expressed as pointwise suprema of continuous functions, and by (33) they attach in 0 in a continuous fashion. We now prove that $I$ is convex. Being suprema of families of linear functions, $I_{+}$and $I_{-}$are convex. Therefore $I$ is convex on $(0,+\infty)$ and $(-\infty, 0)$ separately. To prove the convexity on all $\mathbb{R}$ it remains to show that $I$ is also convex in $\vartheta=0$. Since the left and right branches of $I$ are differentiable, it suffices to show that the left derivative at $\vartheta=0$ is non greater than the right derivative. In fact, they are equal due to (34). Let us now prove that $I$ is strictly convex on the closure of $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$. We know that $I^{\prime}(\vartheta)=\log \varphi_{-}\left(\lambda_{-}(\vartheta)\right)$ on $\left(-1 / \alpha_{-}, 0\right]$ and $I^{\prime}(\vartheta)=-\log \varphi_{+}\left(\lambda_{+}(\vartheta)\right)$ on $\left[0,1 / \alpha_{+}\right)$ (see the proof of Lemma 5.5). By Lemma 5.1(ii) $\log \varphi_{ \pm}$is strictly increasing with positive derivative, while we know that $\lambda_{+}$is a strictly decreasing function on $\left(0,1 / \alpha_{+}\right)$and $\lambda_{-}$is a strictly increasing on $\left(-1 / \alpha_{-}, 0\right)$. Using also that $I^{\prime}$ is continuous at 0 we conclude that $I^{\prime}$ is strictly increasing on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$, hence $I$ is strictly convex on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$.

We conclude with item (v). We know that $I^{\prime}$ is strictly increasing on $\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$. Due to item (iii) it simple to conclude that there exists a unique minimum point $\vartheta_{*} \in\left(-1 / \alpha_{-}, 1 / \alpha_{+}\right)$such that $I$ is strictly decreasing on $\left(-1 / \alpha_{-}, \vartheta_{*}\right)$ and strictly increasing on $\left(\vartheta_{*}, 1 / \alpha_{+}\right)$. It remains to prove that $\vartheta_{*}=v$.

If $v>0$ then by Proposition 5.3(iv) $v=1 / \vartheta_{c}^{+}$and so $I(v)=v J_{+}(1 / v)=\left(1 / \vartheta_{c}^{+}\right) J_{+}\left(\vartheta_{c}^{+}\right)=0$. If $v<0$ then $v=-1 / \vartheta_{c}^{-}$and so $I(v)=-v J_{-}(-1 / v)=\left(1 / \vartheta_{c}^{-}\right) J_{-}\left(\vartheta_{c}^{-}\right)=0$. If, finally, $v=0$ then again by Proposition 5.3(iv) we have $0=\lambda_{c}=I(0)=I(v)$. In all cases $I(v)=0$ and since $I$ is non-negative we conclude that $v=\vartheta_{*}$.

## 6. Proof of Theorem 6(ii)

Below we show how one can deduce the LDP for the process $Z$ itself from the LDP for the hitting times. Due to Theorem 4.1.11 in [6], the LDP for $Z_{t} / t$ holds with rate function $I$ and speed $t$ if we show that

$$
\begin{align*}
& \lim _{\varepsilon \rightarrow 0} \underset{t \rightarrow \infty}{ } \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \geq-I(\vartheta),  \tag{LB}\\
& \lim _{\varepsilon \rightarrow 0} \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \leq-I(\vartheta) . \tag{UB}
\end{align*}
$$

### 6.1. The lower bound

Given $\vartheta \neq 0$ and $\delta, c \in(0,1)$ we define the events

$$
\begin{aligned}
& A_{t}=A_{t}(\delta, \vartheta):=\left\{(1-\delta) t<T_{\lfloor\vartheta t\rfloor}<(1+\delta) t\right\}, \\
& B_{t}=B_{t}(\delta, c):=\{v(t+\delta t)-v(t-\delta t) \leq c t\} .
\end{aligned}
$$

Lemma 6.1. For any $\vartheta \neq 0$ and $\delta, c \in(0,1)$ there exists $\tilde{t}=\tilde{t}(c)>0$ such that

$$
A_{t} \cap B_{t} \subseteq\left\{Z_{t} / t \in(\vartheta-2 c, \vartheta+2 c)\right\}
$$

for all $t>\tilde{t}$.
Proof. Take any $t>0$ and assume $A_{t} \cap B_{t}$ holds. Then, due to assumption (A1),

$$
\begin{aligned}
\left|Z_{t}-\lfloor\vartheta t\rfloor\right| & =\left|Z_{t}-Z_{T_{\lfloor\vartheta t\rfloor}}\right|=\left|W_{v(t)}-W_{v\left(T_{\lfloor\vartheta t\rfloor}\right)}\right| \\
& \leq v\left(t \vee T_{\lfloor\vartheta t\rfloor}\right)-v\left(t \wedge T_{\lfloor\vartheta t\rfloor}\right) \leq v(t+\delta t)-v(t-\delta t) \leq c t .
\end{aligned}
$$

Hence $Z_{t} \in[\lfloor\vartheta t\rfloor-c t,\lfloor\vartheta t\rfloor+c t]$, thus leading to the thesis.
Lemma 6.2. For any $\vartheta \neq 0$ and $\delta \in(0,1)$ it holds


Proof. We give the proof for $\vartheta>0$, the one for $\vartheta<0$ being the same. Note that, fixed $\varepsilon>0$, for $t$ large enough it holds

$$
\frac{1}{t} \log \mathbb{P}\left(A_{t}\right) \geq \frac{1}{t} \log \mathbb{P}\left(\frac{1-\delta}{\vartheta}+\varepsilon<\frac{T_{\lfloor\vartheta t\rfloor}}{\lfloor\vartheta t\rfloor}<\frac{1+\delta}{\vartheta}\right) .
$$

Thanks to the LDP for the hitting times $T_{n}$ this implies that

$$
\begin{aligned}
\varliminf_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(A_{t}\right) & \geq \vartheta \varliminf_{n \rightarrow \infty} \frac{1}{n} \log \mathbb{P}\left(\frac{1-\delta}{\vartheta}+\varepsilon<\frac{T_{n}}{n}<\frac{1+\delta}{\vartheta}\right) \\
& \geq-\vartheta \inf _{\left(\frac{1-\delta}{\vartheta}+\varepsilon, \frac{1+\delta}{\vartheta}\right)} J_{+} \geq-\vartheta J_{+}\left(\frac{1}{\vartheta}\right)
\end{aligned}
$$

as long as $\varepsilon$ is chosen small enough so that $\frac{1}{\vartheta} \in\left(\frac{1-\delta}{\vartheta}+\varepsilon, \frac{1+\delta}{\vartheta}\right)$.
Since $\tau_{i}$ 's are positive i.i.d. random variables, for every $p_{0} \in(0,1)$ we can find some $\eta>0$ such that $p:=\mathbb{P}\left(\tau_{i} \geq\right.$ $\eta)>p_{0}$. In particular, the i.i.d. random variables $r_{i}$ 's with $r_{i}:=\mathbb{1}\left(\tau_{i} \geq \eta\right)$ are Bernoulli of parameter $p$. They are a useful tool to bound the probability of $B_{t}^{c}$ :

Lemma 6.3. For any $\vartheta \neq 0$ and any $c \in(0,1)$, there exists a constant $\delta_{*}=\delta_{*}(\theta, c) \in(0,1)$ depending only on $\theta, c$ such that, for all $\delta \in\left(0, \delta_{*}\right]$, it holds

$$
\lim _{t \rightarrow \infty} \frac{\mathbb{P}\left(B_{t}^{c}\right)}{\mathbb{P}\left(A_{t}\right)}=0
$$

Proof. We restrict to the case $\vartheta>0$, being the proof for $\vartheta<0$ similar. We observe that the event $B_{t}^{c}$ implies the event

$$
\bigcup_{j=0}^{\infty}\left\{v(t-\delta t)=j, \sum_{\ell=1}^{\lceil c t\rceil-1} \tau_{j+1+\ell} \leq 2 \delta t\right\} .
$$

Since the event $\{v(t-\delta t)=j\}$ depends only on $\tau_{1}, \tau_{2}, \ldots, \tau_{j+1}$, by independence we get

$$
\begin{align*}
\mathbb{P}\left(B_{t}^{c}\right) & \leq \sum_{j=0}^{\infty} \mathbb{P}(\nu(t-\delta t)=j) \mathbb{P}\left(\sum_{\ell=1}^{\lceil c t\rceil-1} \tau_{j+1+\ell} \leq 2 \delta t\right) \\
& =\mathbb{P}\left(\sum_{\ell=1}^{\lceil c t\rceil-1} \tau_{\ell} \leq 2 \delta t\right) \leq \mathbb{P}\left(\sum_{\ell=1}^{\lceil c t]-1} r_{\ell} \leq 2 \delta t / \eta\right) . \tag{38}
\end{align*}
$$

Above we have used that $\tau_{\ell} \geq \eta r_{\ell}$. Now we use Cramér Theorem for sums of i.i.d. $p$-Bernoulli r.v.'s. The associated rate function is given by (cf. exercise 2.2.23 in [6])

$$
\mathcal{I}_{p}(x)= \begin{cases}x \log \frac{x}{p}+(1-x) \log \frac{1-x}{1-p} & \text { if } x \in[0,1] \\ +\infty & \text { otherwise }\end{cases}
$$

with the convention that $0 \log 0:=0$. Trivially, $\mathcal{I}_{p}$ is strictly decreasing on $[0, p]$ and strictly increasing on $[p, 1]$, while $\mathcal{I}_{p}(p)=0$. Let $t_{*}:=\lceil c t\rceil-1$. Writing

$$
\frac{1}{t} \log \mathbb{P}\left(\sum_{\ell=1}^{[c t]-1} r_{\ell} \leq 2 \delta t / \eta\right)=\frac{t_{*}}{t} \frac{1}{t_{*}} \log \mathbb{P}\left(\frac{1}{t_{*}} \sum_{\ell=1}^{t_{*}} r_{\ell} \leq \frac{2 \delta t}{\eta t_{*}}\right)
$$

and using that $2 \delta t\left(\eta t_{*}\right)^{-1} \leq 3 \delta(\eta c)^{-1}$ for $t$ large enough, we get

$$
\begin{equation*}
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\sum_{\ell=1}^{\lceil c t\rceil-1} r_{\ell} \leq 2 \delta t / \eta\right) \leq-c \inf _{\left(-\infty, \frac{3 \delta}{\eta c}\right]} \mathcal{I}_{p} . \tag{39}
\end{equation*}
$$

Now we have to choose carefully the constants in order to win. Fix $\vartheta>0$ and $c \in(0,1)$. The function $\mathcal{I}_{p}(0)=\log \frac{1}{1-p}$ is increasing in $p$ and $\lim _{p \rightarrow 1} \mathcal{I}_{p}(0)=+\infty$. In particular, there exists $p_{0}>0$ such that $\mathcal{I}_{p}(0)>\vartheta J_{+}(1 / \vartheta) / c$ for all $p \geq p_{0}$. We fix $\eta$ such that $p:=\mathbb{P}\left(\tau_{i} \geq \eta\right)>p_{0}$.

If $p=1$ then $\tau_{i} \geq \eta$ a.s. In particular, equation (38) gives $\mathbb{P}\left(B_{t}^{c}\right) \leq \mathbb{1}\left(c t-1 \leq \frac{2 \delta t}{\eta}\right)$, so by setting $\delta_{*}=\eta c / 4$ we have that for any $\delta \leq \delta_{*}$ and $t$ large enough $\mathbb{P}\left(B_{t}^{c}\right)=0$. This, combined with Lemma 6.2 , gives the thesis.

Assume, on the other hand, that $p<1$. Recall that $\mathcal{I}_{p}(0)>\vartheta J_{+}(1 / \vartheta) / c$. Since $\lim _{\varepsilon} \backslash \mathcal{I}_{p}(\varepsilon)=\mathcal{I}_{p}(0)$ and $\mathcal{I}_{p}$ is decreasing near 0 , we can fix $\varepsilon_{0}>0$ such that $\mathcal{I}_{p}(\varepsilon)>\vartheta J_{+}(1 / \vartheta) / c$ for all $\varepsilon \in\left[0, \varepsilon_{0}\right]$. Note that the (now fixed) constants $\eta, p, \varepsilon_{0}$ depend only on $\vartheta, c$. To conclude let $\delta_{*}=\left(\eta c \varepsilon_{0} / 4\right) \wedge 1$. Then for each $\delta \in\left(0, \delta_{*}\right]$ we have $3 \delta(\eta c)^{-1} \leq \varepsilon_{0}$ and therefore the last term in (39) is strictly bounded from above by $-\vartheta J_{+}(1 / \vartheta)$. Coming back to (38) and (39) we conclude that

$$
\begin{equation*}
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(B_{t}^{c}\right)<-\vartheta J_{+}(1 / \vartheta) \tag{40}
\end{equation*}
$$

The above bound together with Lemma 6.2 implies the thesis.
Combining Lemmas 6.2 and 6.3 we can prove the following key lower bound:
Lemma 6.4. For any $\vartheta \neq 0$ and $\varepsilon \in(0,1 / 2)$ the following holds

$$
\begin{equation*}
\underline{\lim }_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \geq-I(\vartheta) . \tag{41}
\end{equation*}
$$

Proof. Given $\varepsilon>0$, take $c:=\varepsilon / 2$ and $\delta:=\delta_{*}(\vartheta, c)$ in the definition of $A_{t}, B_{t}$ given in Lemma 6.1, where the constant $\delta_{*}$ is as in Lemma 6.3. Due to Lemma 6.1 for $t$ large enough we have

$$
\mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \geq \mathbb{P}\left(A_{t} \cap B_{t}\right) \geq \mathbb{P}\left(A_{t}\right)-\mathbb{P}\left(B_{t}^{c}\right)=\mathbb{P}\left(A_{t}\right)\left(1-\frac{\mathbb{P}\left(B_{t}^{c}\right)}{\mathbb{P}\left(A_{t}\right)}\right)
$$

which implies

$$
\frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \geq \frac{1}{t} \log \mathbb{P}\left(A_{t}\right)+\frac{1}{t} \log \left(1-\frac{\mathbb{P}\left(B_{t}^{c}\right)}{\mathbb{P}\left(A_{t}\right)}\right) .
$$

Using Lemma 6.2 to control the first term in the r.h.s. and Lemma 6.3 to control the second term in the r.h.s. we get the thesis.

Being (41) uniform in $\varepsilon \in(0,1 / 2)$, one can let $\varepsilon \rightarrow 0$ to conclude that the lower bound (LB) holds for all $\vartheta \neq 0$. If $\vartheta=0$, take any $\varepsilon>0$ and let $u=\varepsilon / 2$. Then by Lemma 6.4 one has

$$
\underline{\lim }_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(-\varepsilon,+\varepsilon)\right) \geq \underline{\lim }_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(u-\varepsilon / 4, u+\varepsilon / 4)\right) \geq-I(u) .
$$

Letting then $\varepsilon \rightarrow 0$ and therefore $u \rightarrow 0$ gives (recall Theorems 2, 7)

$$
\lim _{\varepsilon \rightarrow 0} \lim _{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(-\varepsilon, \varepsilon)\right) \geq-\lim _{u \rightarrow 0} I(u)=-I(0) .
$$

This concludes the proof of (LB) for all $\vartheta \in \mathbb{R}$.

### 6.2. The upper bound

We now move to the proof of (UB). This is rather easy if the asymptotic velocity $v$ vanishes.
Lemma 6.5. If $v=0$ then (UB) holds for all $\vartheta \in \mathbb{R}$.
Proof. If $\vartheta=0$ it is enough to observe that by the LLN $Z_{t} / t \rightarrow 0$ almost surely as $t \rightarrow \infty$, and therefore in probability. Since $I(0)=0$ by Theorems 2, 7, we get the thesis.

To deal with the case $\vartheta \neq 0$, recall that $v=0 \Leftrightarrow \lambda_{c}=0 \Leftrightarrow J_{ \pm}$are strictly decreasing on ( $\alpha_{ \pm}, \infty$ ) (see Proposition 5.3). Since $J_{+} \equiv+\infty$ on $\left(-\infty, \alpha_{+}\right)$and due to (29) we conclude that $J_{+}: \mathbb{R} \rightarrow[0,+\infty)$ is a decreasing extended function. Fix now any $\vartheta>0$ and $\varepsilon>0$ such that $\vartheta-\varepsilon>0$. Then, given any $\tilde{\varepsilon}>0$, for $t$ large it holds

$$
\begin{aligned}
\mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) & \leq \mathbb{P}\left(\frac{Z_{t}}{t}>\vartheta-\varepsilon\right) \leq \mathbb{P}\left(Z_{t} \geq\lfloor(\vartheta-\varepsilon) t\rfloor\right) \\
& \leq \mathbb{P}\left(T_{\lfloor(\vartheta-\varepsilon) t\rfloor} \leq t\right)=\mathbb{P}\left(\frac{T_{\lfloor(\vartheta-\varepsilon) t\rfloor}\lfloor(\vartheta-\varepsilon) t\rfloor}{} \leq \frac{1}{\vartheta-\varepsilon}+\tilde{\varepsilon}\right) .
\end{aligned}
$$

Hence, using the LDP for the hitting times $T_{n}$ and the fact that $J_{+}$is decreasing,

$$
\begin{align*}
& \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \\
& \leq(\vartheta-\varepsilon) \varlimsup_{n \rightarrow \infty} \frac{1}{n} \log \mathbb{P}\left(\frac{T_{n}}{n} \leq \frac{1}{\vartheta-\varepsilon}+\tilde{\varepsilon}\right) \\
& \quad \leq-(\vartheta-\varepsilon)_{\left(-\infty, \frac{1}{\vartheta-\varepsilon}+\tilde{\varepsilon}\right)} J_{+}=-(\vartheta-\varepsilon) J_{+}\left(\frac{1}{\vartheta-\varepsilon}+\tilde{\varepsilon}\right) . \tag{42}
\end{align*}
$$

Letting $\tilde{\varepsilon} \rightarrow 0$ and using that $J_{+}$is l.s.c. (see Proposition 5.3) we get that the first member of (42) is bounded from above by $-(\vartheta-\varepsilon) J_{+}(1 /(\vartheta-\varepsilon))$. Taking now the limit $\varepsilon \rightarrow 0$ and using again that $J_{+}$is 1.s.c. we get the thesis for $\vartheta>0$. The proof of (UB) for $\vartheta<0$ follows by similar arguments.

We now prove that (UB) holds for all $\vartheta \in \mathbb{R}$ assuming $v>0$. The case $v<0$ can be addressed in the same way. The proof we present is based on a method introduced in [5], that we re-adapt to our setting. The strategy consists in reducing the problem to proving the following:

Proposition 6.6. Assume $v>0$ and define $S_{t}:=\inf \left\{s \geq t: Z_{s} \leq 0\right\}$. Then it holds

$$
\begin{equation*}
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(S_{t}<\infty\right) \leq-I(0) . \tag{43}
\end{equation*}
$$

The fact that the above result implies that (UB) holds for all $\vartheta \in \mathbb{R}$ can be seen reasoning as in [5], page 1017, with minor modifications. For completeness we give a sketch of the proof in Appendix A.

A detailed proof of Proposition 6.6 is, on the other hand, given below. This choice is due to the presence of a small gap [13] in the proof presented in [5] (see formula (4.14) on page 1020 there), and to the fact that some additional arguments are necessary since our holding times can be in general arbitrarily small while in [5] they are bounded from below by 1 .

Proof of Proposition 6.6. Due to Proposition 5.3, since $v>0, \lambda_{c}>0$ and the critical point $\vartheta_{c}^{ \pm}$of $J_{ \pm \pm}$is finite and positive. Take any $u \in\left(0,1 / \vartheta_{c}^{+}\right)$and fix $c>1$ integer such that $c / u>\vartheta_{c}^{-}$. Let, in order to simplify the notation, $b_{t}:=\mathbb{P}\left(S_{t}<\infty\right)$ with the convention that $b_{t}=1$ if $t<0$. Recall that $T_{\lfloor t u\rfloor}$ is the hitting time of $\lfloor t u\rfloor$, and define

$$
\tilde{T}_{0}:= \begin{cases}\inf \left\{s \geq T_{\lfloor t u\rfloor}: Z_{s}=0\right\}, & \text { if } T_{\lfloor t u\rfloor}<\infty, \\ +\infty & \text { otherwise }\end{cases}
$$

Then we have

$$
\begin{align*}
b_{t} \leq & \mathbb{P}\left(T_{\lfloor t u\rfloor} \geq t\right)+\mathbb{P}\left(T_{\lfloor t u\rfloor}<t, \tilde{T}_{0}-T_{\lfloor t u\rfloor} \geq c t, S_{t}<\infty\right) \\
& +\mathbb{P}\left(T_{\lfloor t u\rfloor}<t, \tilde{T}_{0}-T_{\lfloor t u\rfloor}<c t, S_{t}<\infty\right) . \tag{44}
\end{align*}
$$

For the first term in the r.h.s. of (44) the LDP for the hitting times $T_{n}, n \rightarrow \infty$, implies that

$$
\begin{equation*}
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{T_{\lfloor t u\rfloor}\lfloor t u\rfloor}{\lfloor } \geq \frac{t}{\lfloor t u\rfloor}\right) \leq \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{\left.T_{\lfloor t u\rfloor}^{\lfloor t u\rfloor} \geq \frac{1}{u}\right) \leq-u J_{+}(1 / u)=-I(u) . . . . . ~}{\lfloor t u}\right) \tag{45}
\end{equation*}
$$

Above we have used that $J_{+}$is increasing on $\left(\vartheta_{c}^{+},+\infty\right)$.
For the second term we apply the strong Markov property at time $T_{\lfloor t u\rfloor}$ (cf. Definition 2.1(iv)) to get

$$
\mathbb{P}\left(T_{\lfloor t u\rfloor}<t, \tilde{T}_{0}-T_{\lfloor t u\rfloor} \geq c t, S_{t}<\infty\right) \leq \mathbb{P}\left(T_{-\lfloor t u\rfloor} \geq c t\right)
$$

Therefore, by the LDP for the hitting times $T_{-n}, n \rightarrow \infty$, and the fact that $J_{-}$is increasing on $\left(\vartheta_{c}^{-},+\infty\right)$, we obtain

$$
\begin{align*}
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(T_{\lfloor t u\rfloor}<t, \tilde{T}_{0}-T_{\lfloor t u\rfloor} \geq c t, S_{t}<\infty\right) & \leq \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{T_{-\lfloor t u\rfloor}}{\lfloor t u\rfloor} \geq \frac{c}{u}\right) \\
& \leq-u J_{-}(c / u)=-c I(-u / c) . \tag{46}
\end{align*}
$$

For the third term in the r.h.s. of (44) one has to deal with the critical points of $J_{ \pm}$, so the idea is to localize things. Fix $m \in \mathbb{N}$ positive. Fix $0<u^{\prime}<u$, hence $1 /\lfloor t u\rfloor \leq 1 / t u^{\prime}$ for $t$ large (as we assume). We take $u^{\prime}$ very near to $u$ such that $1 / u^{\prime}>\vartheta_{c}^{+}$and $c / u^{\prime}>1 / \vartheta_{c}^{-}$. Then

$$
\begin{align*}
& \mathbb{P}\left(\frac{T_{\lfloor t u\rfloor}}{\lfloor t u\rfloor}<\frac{1}{u^{\prime}}, \frac{\tilde{T}_{0}-T_{\lfloor t u\rfloor}}{\lfloor t u\rfloor}<\frac{c}{u^{\prime}}, S_{t}<\infty\right) \\
& =\sum_{k=1}^{m} \sum_{\ell=1}^{m c} \mathbb{P}\left(\frac{T_{\lfloor t u\rfloor}}{\lfloor t u\rfloor} \in\left[\frac{k-1}{m u^{\prime}}, \frac{k}{m u^{\prime}}\right), \frac{\tilde{T}_{0}-T_{\lfloor t u\rfloor}}{\lfloor t u\rfloor} \in\left[\frac{\ell-1}{m u^{\prime}}, \frac{\ell}{m u^{\prime}}\right), S_{t}<\infty\right) \\
& \leq \sum_{k=1}^{m} \sum_{\ell=1}^{m c} \mathbb{P}\left(\frac{T_{\lfloor t u\rfloor}\lfloor t u\rfloor}{} \in\left[\frac{k-1}{m u^{\prime}}, \frac{k}{m u^{\prime}}\right]\right) \mathbb{P}\left(\frac{T_{-\lfloor t u\rfloor}}{\lfloor t u\rfloor} \in\left[\frac{\ell-1}{m u^{\prime}}, \frac{\ell}{m u^{\prime}}\right]\right) b_{t-\frac{(k+\ell) t}{m}} \tag{47}
\end{align*}
$$

where we have applied the strong Markov property at times $T_{\lfloor t u\rfloor}$ and $\tilde{T}_{0}$ and used that if $s \leq t$ then $b_{s} \geq b_{t}$ since $S_{s} \leq S_{t}$. Now we analyze each term separately. Define

$$
\begin{align*}
& w_{+}(r, \delta):=\max \left\{\left|J_{+}(s)-J_{+}(t)\right|: s, t \in\left[\vartheta_{c}^{+}, r\right],|s-t| \leq \delta\right\},  \tag{48}\\
& w_{-}(r, \delta):=\max \left\{\left|J_{-}(s)-J_{-}(t)\right|: s, t \in\left[\vartheta_{c}^{-}, r\right],|s-t| \leq \delta\right\},
\end{align*}
$$

with the convention that $w_{ \pm}(r, \delta)=0$ if $r<\vartheta_{c}^{ \pm}$. The LDP for the hitting times $T_{n}$ then gives

$$
\begin{align*}
& \overline{\lim }_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{T_{\lfloor t u\rfloor}^{\lfloor t u\rfloor}}{} \in\left[\frac{k-1}{m u^{\prime}}, \frac{k}{m u^{\prime}}\right]\right) \\
& \quad \leq-u \inf _{\left[\frac{k-1}{m u^{\prime}}, \frac{k}{m u^{\prime}}\right\rfloor} J_{+} \leq-u J_{+}\left(k / m u^{\prime}\right)+u w_{+}\left(\frac{k}{m u^{\prime}}, \frac{1}{m u^{\prime}}\right)=-\frac{k}{m} I\left(\frac{u^{\prime} m}{k}\right)+u w_{+}\left(\frac{k}{m u^{\prime}}, \frac{1}{m u^{\prime}}\right) . \tag{49}
\end{align*}
$$

Similarly we get

$$
\begin{align*}
& \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{T_{-\lfloor t u\rfloor}^{\lfloor t u\rfloor}}{} \in\left[\frac{\ell-1}{m u^{\prime}}, \frac{\ell}{m u^{\prime}}\right]\right) \\
& \quad \leq-u \inf _{\left[\frac{\ell-1}{m u^{\prime}}, \frac{\ell}{m u^{\prime}}\right]} J_{-} \leq-u J_{-}\left(\ell / m u^{\prime}\right)+u w_{-}\left(\frac{\ell}{m u^{\prime}}, \frac{1}{m u^{\prime}}\right)=-\frac{\ell}{m} I\left(-\frac{u m^{\prime}}{\ell}\right)+u w_{-}\left(\frac{\ell}{m u^{\prime}}, \frac{1}{m u^{\prime}}\right) . \tag{50}
\end{align*}
$$

We set

$$
\begin{aligned}
W_{k, \ell} & :=w_{+}\left(\frac{k}{m u^{\prime}}, \frac{1}{m u^{\prime}}\right)+w_{-}\left(\frac{\ell}{m u^{\prime}}, \frac{1}{m u^{\prime}}\right), \\
W & :=\max \left\{w_{+}\left(\frac{1}{u^{\prime}}, \frac{1}{m u^{\prime}}\right), w_{-}\left(\frac{c}{u^{\prime}}, \frac{1}{m u^{\prime}}\right)\right\} .
\end{aligned}
$$

The above inequalities (49) and (50), and the convexity of $I$, we have for any $\varepsilon>0$ and $t$ large enough that

$$
\begin{align*}
& \mathbb{P}\left(\frac{T_{\lfloor t u\rfloor}}{\lfloor t u\rfloor} \in\left[\frac{k-1}{m u^{\prime}}, \frac{k}{m u^{\prime}}\right]\right) \mathbb{P}\left(\frac{T_{-\lfloor t u\rfloor}^{\lfloor t u\rfloor}}{\lfloor x}\left[\frac{\ell-1}{m u^{\prime}}, \frac{\ell}{m u^{\prime}}\right]\right) \\
& \quad \leq e^{t \varepsilon+u^{\prime} t W_{k, \ell}} e^{-t\left[\frac{k}{m} I\left(\frac{u^{\prime} m}{k}\right)+\frac{\ell}{m} I\left(-\frac{u^{\prime} m}{\ell}\right)\right]} \\
& \quad \leq e^{t \varepsilon+u^{\prime} t W_{k, \ell}} e^{-t \frac{(k+\ell)}{m} I(0)} \leq e^{t \varepsilon-t \frac{(k+\ell)}{m}\left(I(0)-\frac{W}{c_{0}}\right)}, \tag{51}
\end{align*}
$$

where $c_{0}:=\min \left\{\vartheta_{c}^{+}, \vartheta_{c}^{-}\right\}$. We explain the last bound. Note that $W_{k, \ell}=0$ if $k \leq \vartheta_{c}^{+} m u^{\prime}$ and $\ell \leq \vartheta_{c}^{-} m u^{\prime}$. On the other hand, since $\frac{k}{m u^{\prime}} \leq \frac{1}{u^{\prime}}$ and $\frac{\ell}{m u^{\prime}} \leq \frac{c}{u^{\prime}}$, we have $W_{k, \ell} \leq W$. Hence it holds

$$
\begin{equation*}
u^{\prime} W_{k, \ell} \leq u^{\prime} W \mathbb{1}\left(k+\ell>c_{0} m u^{\prime}\right) \leq \frac{(k+\ell)}{c_{0} m} W . \tag{52}
\end{equation*}
$$

Let now

$$
\begin{equation*}
J:=\min \{I(u), c I(-c / u), I(0)\}-W / c_{0} . \tag{53}
\end{equation*}
$$

Then putting (44), (45), (46), (47) and (51) we have

$$
b_{t} \leq e^{-t I(u)+t \varepsilon}+e^{-t c I(-u / c)+t \varepsilon}+\sum_{k=1}^{m} \sum_{\ell=1}^{m c} e^{t \varepsilon-t \frac{(k+\ell)}{m} J} b_{t-\frac{(k+\ell) t}{m}} .
$$

Note that if $k+\ell>m$, then $b_{t-\frac{(k+\ell) t}{m}}=1$. Hence we get

$$
\begin{equation*}
b_{t} \leq\left(2+m^{2} c\right) e^{-t J+t \varepsilon}+e^{t \varepsilon} \sum_{\substack{(k, \ell): 1 \leq k \leq m, 1 \leq \ell \leq m c \\ k+\ell \leq m}} e^{-t \frac{(k+\ell)}{m} J} b_{t-\frac{(k+\ell) t}{m}} . \tag{54}
\end{equation*}
$$

Call $x:=\overline{\lim }_{t \rightarrow \infty} \frac{1}{t} \log b_{t} \in[-\infty, 0]$. Since, given a finite family of functions $\left\{f_{i}(t)\right\}_{i \in I}$, it holds $\overline{\lim }_{t \rightarrow \infty} \frac{1}{t} \times$ $\log \left(\sum_{i} f_{i}(t)\right) \leq \max _{i \in I} \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \left(f_{i}(t)\right)$ from (54) we get

$$
x \leq \varepsilon+\max _{j: 2 \leq j \leq m}\left\{-\frac{j J}{m}+\left(1-\frac{j}{m}\right) x\right\}=x+\varepsilon-\min _{j: 2 \leq j \leq m} \frac{j(J+x)}{m} .
$$

The above bound holds for any $\varepsilon>0$, hence we conclude that $0 \leq-\min _{j: 2 \leq j \leq m} \frac{j(J+x)}{m}$. This implies that $J+x \leq 0$, i.e. $\overline{\lim }_{t \rightarrow \infty} \frac{1}{t} \log b_{t} \leq-J$. Now, let $m \rightarrow \infty$ first, so that $W \rightarrow 0$ due to the fact that $J_{ \pm}$are even $C^{1}$ on $\left(\alpha_{ \pm}, \infty\right)$ and $\alpha_{ \pm}<\vartheta_{c}^{ \pm}$(see Proposition 5.3 and recall that $1 / u^{\prime}>\vartheta_{c}^{+}, c / u^{\prime}>\vartheta_{c}^{-}$). Now we let $u \rightarrow 0$. By Theorems 2,7 and since $v>0, \min \{I(u), c I(-c / u), I(0)\}$ converges to $I(0)$ as $u \rightarrow 0$. This leads to the thesis.

## 7. Proof of Theorem 8 (Gallavotti-Cohen type symmetry)

Due to the definition of $I, I(\vartheta)=I(-\vartheta)+c \vartheta$ for all $\vartheta \in \mathbb{R}$ if and only if

$$
\begin{equation*}
J_{+}(\vartheta)=J_{-}(\vartheta)+c, \quad \forall \vartheta>0 . \tag{55}
\end{equation*}
$$

We now prove that (55) and item (iii) with $c=-\log C$ are equivalent. To this aim, assume that $\varphi_{+}(\lambda)=C \varphi_{-}(\lambda)$ for all $\lambda \leq \lambda_{c}$ and some $C>0$. Then $\log \varphi_{+}(\lambda)=\log \varphi_{-}(\lambda)+\log C=\log \varphi_{-}(\lambda)-c$ for all $\lambda \in \mathbb{R}$. Hence, taking the Legendre transform and recalling the definition (6) of $J_{ \pm}$as Legendre transform of $\log \varphi_{ \pm}$, we get (55).

On the other hand suppose that (55) holds. We claim that $J_{+}(\vartheta)=J_{-}(\vartheta)+c$ also for all $\vartheta \leq 0$. Indeed, since $\alpha_{ \pm} \geq$ 0 , the claim follows from Proposition 5.3(ii) for $\vartheta<0$. If $\alpha_{+}, \alpha_{-}$are both positive then Proposition 5.3(ii) implies the claim also for $\vartheta=0$. If $\alpha_{+}, \alpha_{-}$are both zero, then (55) and the right continuity of $J_{ \pm}$at $\alpha_{ \pm}$(see Proposition 5.3(v)) imply the claim for $\vartheta=0$. We now show that $\alpha_{-}$and $\alpha_{+}$must be either both positive or both zero, thus concluding the proof of our claim. Suppose for example that $\alpha_{-}=0$ and $\alpha_{+}>0$. Then we would have $J_{+}(\vartheta)=\infty$ for $\vartheta \in$ $\left(0, \alpha_{+}\right)$(by Proposition 5.3(ii)). This fact together with (55) implies that $J_{-}(\vartheta)=+\infty$ for $\vartheta \in\left(0, \alpha_{+}\right)$. Applying Proposition 5.3(ii) to $J_{-}$we conclude that $\alpha_{+} \leq \alpha_{-}$thus getting a contradiction.

Due to (55) and the above claim we conclude that $J_{+}(\vartheta)=J_{-}(\vartheta)+c$ for all $\vartheta \in \mathbb{R} . J_{+}(\vartheta), J_{-}(\vartheta)+c$ are the Legendre transforms of $\log \varphi_{+}, \log \varphi_{-}-c$, respectively, thought as extended functions from $\mathbb{R}$ to $(-\infty,+\infty)$. Due to Lemma $5.1 \log \varphi_{+}, \log \varphi_{-}-c$ are convex, 1.s.c. and not everywhere infinite. Hence, by the Fenchel-Moreau Theorem (cf. [2]) we conclude that $\log \varphi_{+}=\log \varphi_{-}-c$, i.e. item (iii) holds with $c=-\log C$.

We now prove that item (ii) implies item (iii). To this aim assume that $\tau_{i}$ and $w_{i}$ are independent. Then $f_{+}(\lambda)=$ $\mathbb{E}\left(e^{\lambda \tau_{i}}\right) p$ and $f_{-}(\lambda)=\mathbb{E}\left(e^{\lambda \tau_{i}}\right) q$ for all $\lambda \leq \lambda_{c}$. Combining this with (21) we get

$$
\frac{\varphi_{+}(\lambda)}{\varphi_{-}(\lambda)}=\frac{f_{+}(\lambda)}{f_{-}(\lambda)}=\frac{p}{q}=: C, \quad \forall \lambda \leq \lambda_{c},
$$

which is item (iii).
Finally we prove that item (iii) implies item (ii). Hence assume $\varphi_{+}(\lambda)=C \varphi_{-}(\lambda)$ for all $\lambda \leq \lambda_{c}$. By (21) we have $\frac{\varphi_{+}(\lambda)}{\varphi_{-}(\lambda)}=\frac{f_{+}(\lambda)}{f_{-}(\lambda)}=C$. Moreover, taking $\lambda=0$ in the previous identity, from the definition of $f_{ \pm}$we deduce that $C=p / q$.

In particular, given $\lambda, \gamma \leq 0$, we can write

$$
\begin{aligned}
\mathbb{E}\left(e^{\lambda \tau_{i}+\gamma w_{i}}\right) & =\mathbb{E}\left(e^{\gamma} e^{\lambda \tau_{i}} \mathbb{1}\left(w_{i}=1\right)\right)+\mathbb{E}\left(e^{-\gamma} e^{\lambda \tau_{i}} \mathbb{1}\left(w_{i}=-1\right)\right) \\
& =e^{\gamma} f_{+}(\lambda)+e^{-\gamma} f_{-}(\lambda)=f_{+}(\lambda)\left(e^{\gamma}+e^{-\gamma} \frac{q}{p}\right) .
\end{aligned}
$$

On the other hand:

$$
\begin{aligned}
& \mathbb{E}\left(e^{\lambda \tau_{i}}\right)=\mathbb{E}\left(e^{\lambda \tau_{i}} \mathbb{1}\left(w_{i}=1\right)\right)+\mathbb{E}\left(e^{\lambda \tau_{i}} \mathbb{1}\left(w_{i}=-1\right)\right)=f_{+}(\lambda)\left(1+\frac{q}{p}\right)=\frac{f_{+}(\lambda)}{p}, \\
& \mathbb{E}\left(e^{\gamma w_{i}}\right)=e^{\gamma} p+e^{-\gamma} q=p\left(e^{\gamma}+e^{-\gamma} \frac{q}{p}\right) .
\end{aligned}
$$

Putting all together, we conclude that

$$
\mathbb{E}\left(e^{\lambda \tau_{i}+\gamma w_{i}}\right)=f_{+}(\lambda)\left(e^{\gamma}+e^{-\gamma} \frac{q}{p}\right)=\left(\frac{f_{+}(\lambda)}{p}\right)\left(p\left(e^{\gamma}+e^{-\gamma} \frac{q}{p}\right)\right)=\mathbb{E}\left(e^{\lambda \tau_{i}}\right) \mathbb{E}\left(e^{\gamma w_{i}}\right),
$$

thus implying the independence of $\tau_{i}, w_{i}$.

## 8. Proof of Theorem 3 (LDP via Gärtner-Ellis theorem for Markov r.w.'s)

We introduce a $\mathbb{Z}$-valued process $\left(N_{t}\right)_{t \in \mathbb{R}_{+}}$given by the cell number of $X_{t}$. More precisely, we set $N_{t}:=n$ if $X_{t}=v_{n}$ for some $v \in V \backslash\{\bar{v}\}$. Note that the cell number process is in general not a Markovian process and that $\left|X_{t}^{*}-N_{t}\right| \leq 1$.

We introduce the constant $\kappa$ defined as $\kappa:=\max \{r(x): x \in \mathcal{V}\}$, where $r(x)=\sum_{y:(x, y) \in \mathcal{E}} r(x, y)$. Due to the periodicity (2), $\kappa$ is a well defined constant in $(0,+\infty)$.

Lemma 8.1. For each $n \in \mathbb{Z} \backslash\{0\}$ and $t \in \mathbb{R}_{+}$it holds $\mathbb{P}\left(N_{t}=n\right) \leq e^{\kappa t|n|-|n| \log |n|}$. In particular, for each $\lambda \in \mathbb{R}$ and $t \in \mathbb{R}_{+}$it holds $\mathbb{E}\left(e^{\lambda N_{t}}\right)<+\infty$.

Proof. The event $\left\{N_{t}=n\right\}$ implies that the r.w. $X$ has performed at least $|n|$ jumps within time $t$. On the other hand, by definition of $\kappa$, the random walk $X$ waits at each $x \in \mathcal{V}$ an exponential time of mean at least $1 / \kappa$. Hence (by a coupling argument) $\mathbb{P}\left(N_{t}=n\right) \leq P\left(\mathcal{Z}_{t} \geq|n|\right)$, where $\mathcal{Z}_{t}$ is a Poisson random variable with intensity $\kappa t$. Since $E\left(e^{a \mathcal{Z}_{t}}\right)=e^{\kappa t\left(e^{a}-1\right)}$, by Chebyshev inequality with $a=\log |n|$ we get

$$
\mathbb{P}\left(N_{t}=n\right) \leq P\left(\mathcal{Z}_{t} \geq|n|\right) \leq e^{-|n| \log |n|} E\left(e^{\mathcal{Z}_{t} \log |n|}\right)=e^{-|n| \log |n|+\kappa t(|n|-1)}
$$

thus proving the bound on $\mathbb{P}\left(N_{t}=n\right)$. As a consequence, we obtain

$$
\mathbb{E}\left(e^{\lambda N_{t}}\right) \leq 1+2 \sum_{n=1}^{\infty} e^{\lambda|n|+\kappa t|n|-|n| \log |n|}<+\infty
$$

We now define a new function $F:(V \backslash\{\bar{v}\}) \times \mathbb{R} \times \mathbb{R}_{+} \ni(v, \lambda, t) \rightarrow F(v, \lambda, t) \in \mathbb{R}_{+}$as

$$
\begin{equation*}
F(v, \lambda, t)=\sum_{n \in \mathbb{Z}} e^{\lambda n} \mathbb{P}\left(X_{t}=v_{n}\right)=\mathbb{E}\left(e^{\lambda N_{t}} \mathbb{1}\left(X_{t}=v_{n} \text { for some } n \in \mathbb{Z}\right)\right) \tag{56}
\end{equation*}
$$

Recall that, given $v \neq w$ in $V \backslash\{\bar{v}\}$, we have set

$$
r(v):=r\left(v_{n}\right), \quad r_{-}(w, v):=r\left(w_{n-1}, v_{n}\right), \quad r_{0}(w, v):=r(w, v), \quad r_{+}(w, v):=r\left(w_{n+1}, v_{n}\right)
$$

Lemma 8.2. Given $\lambda \in \mathbb{R}$ consider the finite matrix $\mathcal{A}(\lambda)$ defined in (11). Consider the vector-valued function $\mathbb{R}_{+} \ni$ $t \mapsto F^{(\lambda)}(t) \in \mathbb{R}^{V \backslash\{\bar{v}\}}$ defined as $F^{(\lambda)}(t)_{v}:=F(v, \lambda, t)$. Then $F^{(\lambda)}(\cdot)$ is $C^{1}$ in $t$ and

$$
\begin{equation*}
\partial_{t} F^{(\lambda)}(t)=\mathcal{A}(\lambda) F^{(\lambda)}(t) \tag{57}
\end{equation*}
$$

Proof. Fixed $v, \lambda$, we write $F(v, \lambda, \cdot)$ as the function series $F(v, \lambda, t)=\sum_{n \in \mathbb{Z}} f_{n}(t)$, where $f_{n}(t)=e^{\lambda n} \mathbb{P}\left(X_{t}=v_{n}\right)$. By [27, Theorem 2.8.2], the function $\mathbb{R}_{+} \ni t \mapsto \mathbb{P}\left(X_{t}=v_{n}\right) \in[0,1]$ is differentiable and moreover

$$
\begin{align*}
\partial_{t} \mathbb{P}\left(X_{t}=v_{n}\right)= & -r\left(v_{n}\right) \mathbb{P}\left(X_{t}=v_{n}\right) \\
& +\sum_{w \in V \backslash\{\bar{v}, v\}}\left[r\left(w_{n-1}, v_{n}\right) \mathbb{P}\left(X_{t}=w_{n-1}\right)\right. \\
& \left.+r\left(w_{n}, v_{n}\right) \mathbb{P}\left(X_{t}=w_{n}\right)+r\left(w_{n+1}, v_{n}\right) \mathbb{P}\left(X_{t}=w_{n+1}\right)\right] . \tag{58}
\end{align*}
$$

Then, by Lemma 8.1, we conclude that, for $M>0$ and $n \in \mathbb{Z}$ with $|n| \geq 2$, it holds

$$
\begin{aligned}
& \left\|f_{n}\right\|_{L^{\infty}[-M, M]} \leq e^{|\lambda| \cdot|n|+\kappa M-|n| \log |n|} \\
& \left\|\partial_{t} f_{n}\right\|_{L^{\infty}[-M, M]} \leq 4 \kappa|V| e^{|\lambda| \cdot|n|+\kappa(|n|+1) M-(|n|-1) \log (|n|-1)}
\end{aligned}
$$

The space $C^{1}[-M, M]$ (of functions $C^{1}$ on $(-M, M)$, such that they and their first derivates have continuous extensions to $[-M, M])$ is a Banach space endowed with the norm $\|f\|:=\|f\|_{L^{\infty}[-M, M]}+\left\|\partial_{t} f\right\|_{L^{\infty}[-M, M]}$. We therefore conclude that $F(v, \lambda, t)=\sum_{n \in \mathbb{Z}} f_{n}(t)$ belongs to $C^{1}(\mathbb{R})$ and $\partial_{t} F(v, \lambda, t)=\sum_{n \in \mathbb{Z}} f_{n}^{\prime}(t)$, i.e.

$$
\begin{aligned}
\partial_{t} F(v, \lambda, t) & =\sum_{n \in \mathbb{Z}} e^{\lambda n} \partial_{t} \mathbb{P}\left(X_{t}=v_{n}\right) \\
& =-r(v) F(v, \lambda, t)+\sum_{w \in V \backslash\{\bar{v}, v\}}\left[e^{\lambda} r_{-}(w, v)+r_{0}(w, v)+e^{-\lambda} r_{+}(w, v)\right] F(w, \lambda, t)
\end{aligned}
$$

This concludes the proof.

We can now conclude the proof of Theorem 3 . Since $\left|X_{t}^{*}-N_{t}\right| \leq 1$, it is enough to prove the same LDP for $N_{t} / t$. Due to Lemma 8.2 we have $F^{(\lambda)}(t)=e^{(t-1) \mathcal{A}(\lambda)} F^{(\lambda)}(1)$. Since the graph $\mathcal{G}=(\mathcal{V}, \mathcal{E})$ is connected, Definition (56) implies that the vector $F^{(\lambda)}(1)$ has strictly positive entries. In particular we can write

$$
\mathbb{E}\left(e^{\lambda N_{t}}\right)=e^{-\kappa(t-1)} \sum_{v, v^{\prime} \in V \backslash\{\bar{v}\}}\left[e^{[\mathcal{A}(\lambda)+\kappa](t-1)}\right]_{v, v^{\prime}} F\left(v^{\prime}, \lambda, 1\right)
$$

where $\kappa:=\max \{r(x): x \in \mathcal{V}\}$, as above. Note that $\mathcal{A}(\lambda)+\kappa$ is an irreducible matrix with nonnegative entries and therefore, by Perron-Frobenius theorem, it has a simple positive eigenvalue $\bar{\gamma}$ and an associated eigenvector with strictly positive entries $(a(v))_{v \in V \backslash\{\bar{v}\}}$, while any other eigenvalue $\bar{\gamma}^{\prime}$ is such that $\left|\bar{\gamma}^{\prime}\right| \leq \bar{\gamma}$ (in particular, $\mathcal{R}\left(\bar{\gamma}^{\prime}\right)<$ $\mathcal{R}(\bar{\gamma})$ ). The above eigenvalue $\bar{\gamma}$ is the so called Perron-Frobenius eigenvalue and equals the spectral radius of $\mathcal{A}(\lambda)+\kappa$ (note that $\bar{\gamma}=\bar{\gamma}(\lambda)$ ). Call

$$
\left\{\begin{array}{l}
C(\lambda):=\max \left\{F(v, \lambda, 1) / a_{v}: v \in V \backslash\{\bar{v}\}\right\} \\
c(\lambda):=\min \left\{F(v, \lambda, 1) / a_{v}: v \in V \backslash\{\bar{v}\}\right\}
\end{array}\right.
$$

Note that $C(\lambda), c(\lambda)$ are positive constants. Then

$$
\begin{aligned}
& \mathbb{E}\left(e^{\lambda N_{t}}\right) \leq C(\lambda) e^{-\kappa(t-1)} \sum_{v, v^{\prime} \in V \backslash\{\bar{v}\}}\left[e^{[\mathcal{A}(\lambda)+\kappa](t-1)}\right]_{v, v^{\prime}} a_{v^{\prime}}=e^{(\bar{\gamma}-\kappa)(t-1)} C(\lambda) \sum_{v \in V \backslash\{\bar{v}\}} a_{v}, \\
& \mathbb{E}\left(e^{\lambda N_{t}}\right) \geq c(\lambda) e^{-\kappa(t-1)} \sum_{v, v^{\prime} \in V \backslash\{\bar{v}\}}\left[e^{[\mathcal{A}(\lambda)+\kappa](t-1)}\right]_{v, v^{\prime}} a_{v^{\prime}}=e^{(\bar{\gamma}-\kappa)(t-1)} c(\lambda) \sum_{v \in V \backslash\{\bar{v}\}} a_{v} .
\end{aligned}
$$

It then follows that the limit $\lim _{t \rightarrow \infty} \frac{1}{t} \ln \mathbb{E}\left(e^{\lambda N_{t}}\right)$ exists and equals $\bar{\gamma}(\lambda)-\kappa$, which corresponds to $\Lambda(\lambda)$ by the previous discussion.

By finite-dimensional perturbation theory [25], the Perron-Frobenius eigenvalue $\bar{\gamma}=\bar{\gamma}(\lambda)$ is differentiable in $\lambda$, thus implying that $\Lambda(\lambda)$ is differentiable in $\lambda$. At this point the thesis follows from Gärtier-Ellis theorem (cf. [15, Lemma V. 4 and Theorem V.6]).

## 9. Proof of Theorem 5 (GC type symmetry for Markov r.w.'s)

We start with a technical result, that is also useful in the applications for the computation of the functions $f_{ \pm}(\lambda)$. Consider a generic stochastic process $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$as in Definition 2.1. Define

$$
J_{1}:=\inf \left\{t>0: X_{t} \in\left\{-1_{*}, 0_{*}, 1_{*}\right\}, \exists s \in(0, t) \text { with } X_{s} \neq X_{0}\right\}
$$

and set

$$
\begin{equation*}
\tilde{f}_{ \pm}(\lambda):=\mathbb{E}_{0_{*}}\left(e^{\lambda J_{1}} \mathbb{1}\left(X_{J_{1}}= \pm 1_{*}\right)\right), \quad \tilde{f}_{0}(\lambda):=\mathbb{E}_{0_{*}}\left(e^{\lambda J_{1}} \mathbb{1}\left(X_{J_{1}}=0_{*}\right)\right) \tag{59}
\end{equation*}
$$

Lemma 9.1. If $\tilde{f}_{0}(\lambda)<1$, then $f_{+}(\lambda)=\frac{\tilde{f}_{+}(\lambda)}{1-\tilde{f}_{0}(\lambda)}, f_{-}(\lambda)=\frac{\tilde{f}_{-}(\lambda)}{1-\tilde{f}_{0}(\lambda)}$. If $\tilde{f}_{0}(\lambda) \geq 1$, then $f_{+}(\lambda)=f_{-}(\lambda)=+\infty$.
Proof. We call $J_{k}$ 's the consecutive times at which the stochastic process $\left(X_{t}\right)_{t \geq 0}$ hits the states of type $n_{*}$ :

$$
\left\{\begin{array}{l}
J_{0}:=0 \\
J_{k}:=\inf \left\{t>J_{k-1}: X_{t} \in\left\{-1_{*}, 0_{*}, 1_{*}\right\}, \exists s \in\left(J_{k-1}, t\right) \text { with } X_{s} \neq X_{J_{k-1}}\right\}, \quad k \geq 1
\end{array}\right.
$$

We can write $S=\sum_{k=0}^{\infty} \mathbb{1}\left(X_{J_{0}}=\cdots=X_{J_{k}}=0_{*}, X_{J_{k+1}} \in\left\{-1_{*}, 1_{*}\right\}\right) J_{k+1}$. Taking the exponential at both sides and multiplying by $\mathbb{1}\left(X_{S}=1_{*}\right)$ we get

$$
e^{\lambda S_{\mathbb{1}}}\left(X_{S}=1_{*}\right)=\sum_{k=0}^{\infty} \mathbb{1}\left(X_{J_{0}}=\cdots=X_{J_{k}}=0_{*}, X_{J_{k+1}}=1_{*}\right) e^{\lambda J_{k+1}}
$$

Note that, by Definition 2.1, w.r.t. the probability measure $\mathbb{P}_{0_{*}}\left(\cdot \mid X_{J_{0}}=\cdots=X_{J_{k}}=0_{*}, X_{J_{k+1}}=1_{*}\right)$, the random variables $e^{\lambda\left(J_{i}-J_{i-1}\right)}, 1 \leq i \leq k+1$, are independent with expectation $\mathbb{E}_{0_{*}}\left(e^{\lambda J_{1}} \mid X_{J_{1}}=0_{*}\right)$ if $1 \leq i \leq k$ and $\mathbb{E}_{0_{*}}\left(e^{\lambda J_{1}} \mid X_{J_{1}}=1_{*}\right)$ if $i=k+1$. Hence,

$$
\begin{aligned}
f_{+}(\lambda) & =\mathbb{E}_{0_{*}}\left(e^{\lambda S_{\mathbb{1}}}\left(X_{S}=1_{*}\right)\right) \\
& =\sum_{k=0}^{\infty} \mathbb{P}_{0_{*}}\left(X_{J_{1}}=0_{*}\right)^{k} \mathbb{P}_{0_{*}}\left(X_{J_{1}}=1_{*}\right) \mathbb{E}_{0_{*}}\left(e^{\lambda J_{1}} \mid X_{J_{1}}=0_{*}\right)^{k} \mathbb{E}_{0_{*}}\left(e^{\lambda J_{1}} \mid X_{J_{1}}=1_{*}\right) \\
& =\sum_{k=0}^{\infty} \tilde{f}_{0}(\lambda)^{k} \tilde{f}_{+}(\lambda) .
\end{aligned}
$$

A similar expression holds for $f_{-}(\lambda)$. At this point it is immediate to derive the thesis.
Let us now come back to the same context of Section 3.2: $\left(X_{t}\right)_{t \in \mathbb{R}_{+}}$is a Markov random walk on the quasi 1d lattice $\mathcal{G}=(\mathcal{V}, \mathcal{E})$, with positive rates $r(x, y),(x, y) \in \mathcal{E}$, such that (2) and (12) hold. For the rest of this section, we refer to Markov random walks without stating explicitly that they are Markov. Due to (12) in the figures of $G, \mathcal{G}$ we draw only unoriented edges with the convention that for each unoriented edge $\{x, y\}$ the graph in consideration presents both the edge $(x, y)$ and the edge $(y, x)$.

Recall that given an edge $(u, v) \in E$ in the fundamental graph $G=(V, E)$, we have defined (cf. (13)) $r(u, v)=$ $r(\pi(u), \pi(v))$ where $\pi$ is the map $V \rightarrow \mathcal{V}$ such that $\pi(u)=u_{0}$ if $u \neq \bar{v}$ and $\pi(\bar{v})=\underline{v}_{1}=1_{*}$. Given $v \in V$ we set

$$
\begin{equation*}
r(v):=\sum_{y:(\pi(v), y) \in \mathcal{E}} r(\pi(v), y) . \tag{60}
\end{equation*}
$$

Note that $r(\underline{v})=r(\bar{v})$. We point out that the map $\pi: V \rightarrow \mathcal{V}$ does not induce a graph embedding of $G$ into $\mathcal{G}$. Indeed, problems come from the neighbors of $\underline{v}, \bar{v}$ in $G$. Consider for example the fundamental graph $G$ in Figure 3. Then $x_{0}, y_{-1}, z_{-1}$ are neighboring points of $\underline{v}_{0}$, while $x_{1}, y_{0}, z_{0}$ are neighboring points of $\underline{v}_{1}$. Despite this phenomenon, the map $\pi$ induces an isomorphism between the family of paths $\left(x_{0}, x_{1}, \ldots, x_{m}\right)$ in $G$ from $\underline{v}$ to $\bar{v}$ with interior points in $V \backslash\{\underline{v}, \bar{v}\}$ and the family of paths $\left(x_{0}^{\prime}, x_{1}^{\prime}, \ldots, x_{m}^{\prime}\right)$ in $\mathcal{G}$ from $\underline{v}_{0}=0_{*}$ to $\underline{v}_{1}=1_{*}$ with interior points in $\mathcal{V} \backslash\left\{0_{*}, 1_{*}\right\}$, moreover it holds $r\left(x_{i}, x_{i+1}\right)=r\left(\pi\left(x_{i}\right), \pi\left(x_{i+1}\right)\right)$ for $0 \leq i<m$ and $r\left(x_{i}\right)=r\left(\pi\left(x_{i}\right)\right)$ for $0 \leq i \leq m$. This property will be used below.

By Theorem 8 the Gallavotti-Cohen type symmetry (14) is satisfied for some constant $\Delta$ if and only if $\varphi_{+}(\lambda) / \varphi_{-}(\lambda)=e^{\Delta}$ for all $\lambda \leq \lambda_{c}$. On the other hand, by (21) and the above Lemma 9.1, it holds

$$
\begin{equation*}
\frac{\varphi_{+}(\lambda)}{\varphi_{-}(\lambda)}=\frac{f_{+}(\lambda)}{f_{-}(\lambda)}=\frac{\tilde{f}_{+}(\lambda)}{\tilde{f}_{-}(\lambda)}, \quad \forall \lambda \leq \lambda_{c} . \tag{61}
\end{equation*}
$$

Given an integer $m \geq 1$, let $\mathcal{A}_{m}$ be the family of sequences $\left(x_{0}, x_{1}, \ldots, x_{m}\right)$ such that $x_{0}=\underline{v}, x_{m}=\bar{v},\left(x_{i}, x_{i+1}\right) \in E$ for all $i: 0 \leq i<m$ and $x_{i} \in V \backslash\{\underline{v}, \bar{v}\}$ for all $0<i<m$. We call $\mathcal{A}_{m}^{*}$ the family of sequences satisfying the same


Fig. 3. Example of fundamental graph $G$.
properties as above when exchanging the role of $\underline{v}$ and $\bar{v}$. Then we can write

$$
\begin{align*}
\tilde{f}_{+}(\lambda) & =\sum_{m=1}^{\infty} \sum_{\left(x_{0}, x_{1}, \ldots, x_{m}\right) \in \mathcal{A}_{m}} \int_{\mathbb{R}_{+}^{m-1}} d t_{1} d t_{2} \cdots d t_{m-1} e^{\sum_{i=0}^{m-1}\left(\lambda-r\left(x_{i}\right)\right) t_{i}} \prod_{i=0}^{m-1} r\left(x_{i}, x_{i+1}\right) \\
& = \begin{cases}\sum_{m=1}^{\infty} \sum_{\left(x_{0}, x_{1}, \ldots, x_{m}\right) \in \mathcal{A}_{m}} \prod_{i=0}^{m-1} r\left(x_{i}, x_{i+1}\right) \prod_{i=0}^{m-1} \frac{1}{r\left(x_{i}\right)-\lambda} & \text { if } \lambda<\min _{x \in V} r(x), \\
+\infty & \text { otherwise. }\end{cases} \tag{62}
\end{align*}
$$

A similar expression holds for $\tilde{f}_{-}(\lambda)$, with $\mathcal{A}_{m}$ replaced by $\mathcal{A}_{m}^{*}$.
Given $\gamma=\left(x_{0}, x_{1}, \ldots, x_{m}\right)$ we set

$$
R_{\gamma}:=\prod_{i=0}^{m-1} r\left(x_{i}, x_{i+1}\right), \quad S_{\gamma}(\lambda)=\prod_{i=1}^{m-1} \frac{1}{r\left(x_{i}\right)-\lambda}
$$

Note that the product in $S_{\gamma}(\lambda)$ starts from $i=1$, hence all sites $x_{i}$ appearing in the product belong to $V \backslash\{\underline{v}, \bar{v}\}$. By mean of this notation we can write

$$
\begin{equation*}
\tilde{f}_{+}(\lambda)=\frac{1}{r(\underline{v})-\lambda} \sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}} R_{\gamma} S_{\gamma}(\lambda), \quad \tilde{f}_{-}(\lambda)=\frac{1}{r(\bar{v})-\lambda} \sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}^{*}} R_{\gamma} S_{\gamma}(\lambda) . \tag{63}
\end{equation*}
$$

Hence, due to (61), the Gallavotti-Cohen type symmetry (14) is satisfied for some constant $\Delta$ if and only if

$$
\begin{equation*}
\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}} R_{\gamma} S_{\gamma}(\lambda)=e^{\Delta} \sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}^{*}} R_{\gamma} S_{\gamma}(\lambda), \quad \lambda \leq \lambda_{c} \tag{64}
\end{equation*}
$$

Given $\gamma=\left(x_{0}, x_{1}, \ldots, x_{m}\right) \in \mathcal{A}_{m}$ we write $\bar{\gamma}$ for the loop-erased version of $\gamma$. We recall that $\bar{\gamma}$ is obtained by erasing all the loops of $\gamma$ in chronological order. More precisely, consider the following algorithm. Set $i_{0}:=0$ and, once defined $i_{0}, i_{1}, \ldots, i_{k}$, set $r:=k$ if $i_{k}=m$, otherwise (if $i_{k}<m$ ) set

$$
i_{k+1}:=\max \left\{j: i_{k} \leq j<m \text { and } x_{j}=x_{i_{k}}\right\}+1
$$

(recall that $\gamma \in \mathcal{A}_{m}$ visits $\bar{v}$ only as last point $\left.x_{m}=\bar{v}\right)$. Then $\bar{\gamma}:=\left(x_{i_{0}}, x_{i_{1}}, \ldots, x_{i_{r}}\right)$ (see Figure 4). Since $\gamma \in \mathcal{A}_{m}$ it must be $\bar{\gamma} \in \mathcal{A}_{r}$. Note that

$$
\begin{equation*}
R_{\gamma}=R_{\bar{\gamma}} R_{\gamma}^{\mathrm{loop}}, \quad R_{\gamma}^{\mathrm{loop}}:=\prod_{k=0}^{r-1} \prod_{i=i_{k}}^{i_{k+1}-2} r\left(x_{i}, x_{i+1}\right) \tag{65}
\end{equation*}
$$

with the convention that $\prod_{i=i_{k}}^{i_{k+1}-2} r\left(x_{i}, x_{i+1}\right)=1$ if $i_{k+1}=i_{k}+1\left(R_{\gamma}^{\text {loop }}\right.$ is the contribution to $R_{\gamma}$ given by factors associated to the edges inside the loops).

We write $\gamma^{\dagger}$ for the path in $\mathcal{A}_{m}^{*}$ going from $\bar{v}$ to $\underline{v}$ and obtained from $\gamma$ by reversing the order out of the loops and keeping the same order in the loops. More precisely, with the notation introduced above, it holds

$$
\gamma^{\dagger}=\left(x_{i_{r}}, \star \star \star, x_{i_{r-1}}, \star \star \star, x_{i_{r-2}}, \star \star \star, x_{i_{1}}, \star \star \star, x_{i_{0}}\right),
$$



Fig. 4. Example of a path $\gamma \in \mathcal{A}_{9} \cdot \bar{\gamma}=\left(\underline{v}, x_{1}, x_{2}, x_{3}, x_{8}, \bar{v}\right), R_{\gamma}^{\text {loop }}=r\left(x_{3}, x_{4}\right) r\left(x_{4}, x_{5}\right) r\left(x_{5}, x_{6}\right) r\left(x_{6}, x_{7}\right), \gamma^{\dagger}=\left(\bar{v}, x_{8}, x_{7}, x_{4}, x_{5}, x_{6}, x_{3}\right.$, $\left.x_{2}, x_{1}, \underline{v}\right)$.
where the pieces marked by $\star \star \star$ are determined as follows. Take $k: 1 \leq k \leq r$. If $i_{k}=i_{k-1}+1$, then the piece " $x_{i_{k}}$, $\star \star \star, x_{i_{k-1}}$ " is indeed simply " $x_{i_{k}}$, $x_{i_{k-1}}$." If $i_{k}>i_{k-1}+1$, then the piece " $x_{i_{k}}$, $\star \star \star, x_{i_{k-1}}$ " is given by " $x_{i_{k}}, x_{i_{k}-1}=$ $x_{i_{k-1}}, x_{i_{k-1}+1}, x_{i_{k-1}+2}, \ldots, x_{i_{k}-1}=x_{i_{k-1}}$ "(see Figure 4).

We note that the transformation $\gamma \mapsto \gamma^{\dagger}$ is an involute bijection from $\mathcal{A}_{m}$ to $\mathcal{A}_{m}^{*}$ and that

$$
\begin{equation*}
\overline{\left(\gamma^{\dagger}\right)}=(\bar{\gamma})^{\dagger}, \quad R_{\gamma}^{\text {loop }}=R_{\gamma^{\dagger}}^{\text {loop }}, \quad S_{\gamma}(\lambda)=S_{\gamma^{\dagger}}(\lambda) \tag{66}
\end{equation*}
$$

The above identities, together with (64) and (65), imply the following criterion:
Criterion 1. The Gallavotti-Cohen type symmetry (14) is satisfied for some constant $\Delta$ if and only if

$$
\begin{equation*}
\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}} R_{\gamma}^{\text {loop }}\left(R_{\bar{\gamma}}-e^{\Delta} R_{\bar{\gamma}^{\dagger}}\right) S_{\gamma}(\lambda)=0, \quad \lambda \leq \lambda_{c} \tag{67}
\end{equation*}
$$

In particular, (14) is satisfied for some constant $\Delta$ if it holds $R_{\gamma}=e^{\Delta} R_{\gamma^{\dagger}}$ for any loop-free path from $\underline{v}$ to $\bar{v}$.
If $G$ is $(\underline{v}, \bar{v})$-minimal, then trivially the above criterion implies that (14) is satisfied with $\Delta$ given by (15).
We now prove the reverse implication, in particular we show the following result:
Proposition 9.2. Suppose that the fundamental graph $G$ is not $(\underline{v}, \bar{v})$-minimal. Call $\mathcal{R} \subset(0,+\infty)^{E}$ the family of vectors $\underline{r}:=(r(e): e \in E)$ for which (64) holds ${ }^{4}$ for $\lambda \leq 0$ (instead of $\lambda \leq \lambda_{c}$ ) and for some constant $\Delta$ which can depend on $\underline{r}$ (i.e. $\Delta=\Delta(\underline{r})$ ). Then $\mathcal{R}$ has zero Lebesgue measure in $(0,+\infty)^{E}$.

Consider the oriented subgraph $\hat{G}=(\hat{V}, \hat{E})$ consisting of the points in $V$ and edges in $E$ that appear in some path $\gamma$ as $\gamma$ varies in $\mathcal{A}_{m}$ and $m$ varies in $\{1,2, \ldots\}$. Note that $G$ is $(\underline{v}, \bar{v})$-minimal if and only $\hat{G}$ is $(\underline{v}, \bar{v})$-minimal. Moreover, equation (64) for $\lambda \leq 0$ (instead of $\lambda \leq \lambda_{c}$ ) is the same for the two graphs $G, \hat{G}$. Hence, to simplify the notation below and without loss of generality, we assume from now on that $G=\hat{G}$.

### 9.1. Proof of Proposition 9.2

Given $\underline{r} \in \mathbb{C}^{E}$ we define $\Re(\underline{r}) \in \mathbb{R}^{E}$ as the vector whose entries are the real part of the entries of $\underline{r}$, i.e.

$$
\mathfrak{R}(r)(e):=\mathfrak{R}(r(e)), \quad e \in E
$$

We define $\Omega_{*} \subset \mathbb{C} \times \mathbb{C}^{E}$ as the set of vectors $(\lambda, \underline{r})$ satisfying the following properties:
(i) $\mathfrak{R}(\underline{r}) \in(0,+\infty)^{E}$,
(ii) $|r(e)|<2 \mathfrak{R}(r(e)) \forall e \in E$,
(iii) $-\mathfrak{R}(\lambda)>3 \sum_{e \in E} \mathfrak{R}(r(e))+1$.

Note that $-\mathfrak{R}(\lambda)>0$ and that $\Omega:=\Omega_{*} \cap\left(\mathbb{R} \times \mathbb{R}^{E}\right)$ is simply given by the elements $(\lambda, \underline{r})$ with $\underline{r} \in(0,+\infty)^{E}$ and $-\lambda>3 \sum_{e \in E} r(e)+1$.

Claim 9.3. There exist holomorphic functions $h_{ \pm}: \Omega_{*} \rightarrow \mathbb{C}$ such that

$$
\begin{equation*}
h_{+}(\lambda, \underline{r})=\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}} R_{\gamma} S_{\gamma}(\lambda), \quad h_{-}(\lambda, \underline{r})=\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}^{*}} R_{\gamma} S_{\gamma}(\lambda) \tag{68}
\end{equation*}
$$

for all $(\lambda, \underline{r}) \in \Omega=\Omega_{*} \cap\left(\mathbb{R} \times(0,+\infty)^{E}\right)$.

[^3]Note that, for $(\lambda, \underline{r}) \in \Omega$, the series in (68) are well posed since given by positive addenda. The proof of Claim 9.3 is postponed to Section 9.1.1.

Since $h_{-}>0$ on $\Omega$, there exists an open subset $\Omega_{* *} \subset \mathbb{C} \times \mathbb{C}^{E}$ with $\Omega \subset \Omega_{* *} \subset \Omega_{*}$ and such that $h_{-} \neq 0$ on $\Omega_{* *}$. At cost to restrict $\Omega_{* *}$ we can assume that the set

$$
\begin{equation*}
\left\{\lambda \in \mathbb{C}:(\lambda, \underline{r}) \in \Omega_{* *}\right\} \tag{69}
\end{equation*}
$$

is connected for any fixed $\underline{r} \in(0,+\infty)^{E}$. Since the function $h_{+} / h_{-}$is well defined and holomorphic on $\Omega_{* *}$, also the derivative $\partial_{\lambda}\left(h_{+} / h_{-}\right)$is holomorphic on $\Omega_{* *}$ (cf. [3, Section IV.2.2]).

Due to the definition of $\mathcal{R}$, fixed $\underline{r} \in \mathcal{R}$, the map $\lambda \mapsto \partial_{\lambda}\left(h_{+} / h_{-}\right)(\lambda, \underline{r})$ is zero for $\lambda$ real and sufficiently small. Being holomorphic, this map must be zero on all the set (69) and we get the following result (see Section 9.1.1 for a detailed proof):

Claim 9.4. $\partial_{\lambda}\left(h_{+} / h_{-}\right)(\lambda, \underline{r})=0$ for any $(\lambda, \underline{r}) \in \Omega_{* *}$ with $\underline{r} \in \mathcal{R}$.
Suppose now, by contradiction, that the set $\mathcal{R}$ has positive Lebesgue measure (here and in what follows we refer to the $|E|$-dimensional Lebesgue measure). Given $\lambda<0$, consider the open and connected set $\Omega_{\lambda}:=\left\{\underline{r} \in(0, \infty)^{E}\right.$ : $(\lambda, \underline{r}) \in \Omega\}$ and define the function $h_{\lambda}: \Omega_{\lambda} \rightarrow \mathbb{R}$ as $h_{\lambda}(\underline{r}):=\partial_{\lambda}\left(h_{+} / h_{-}\right)(\lambda, \underline{r})$. Note that $h_{\lambda}$ is a real analytic function (locally it admits a convergent power series expansion, since restriction of an holomorphic function). Since $\Omega_{\lambda} \subset \Omega_{\lambda^{\prime}}$ if $\lambda^{\prime}<\lambda$ and since $\bigcup_{\lambda<0} \Omega_{\lambda}=(0,+\infty)^{E}$, we can find $\lambda_{0}<0$ such that $\Omega_{\lambda} \cap \mathcal{R}$ has positive Lebesgue measure for $\lambda \leq \lambda_{0}$ (from now on we assume $\lambda \leq \lambda_{0}$ ). This fact and Claim 9.4 imply that (i) the set $\left\{\underline{r} \in \Omega_{\lambda}: h_{\lambda}(r)=0\right\}$ contains $\Omega_{\lambda} \cap \mathcal{R}$ and therefore has positive Lebesgue measure. On the other hand, (ii) $h_{\lambda}$ is the restriction to $\Omega_{\lambda}$ of the holomorphic function $\partial_{\lambda}\left(h_{+} / h_{-}\right)(\lambda, \cdot)$ defined on an open subset of $\mathbb{C}^{E}$ containing $\Omega_{\lambda}$. As a byproduct of (i), (ii) and Lemma B. 1 in Appendix B (the latter is based on the Weierstrass Preparation Theorem) we get the following:

Claim 9.5. For any $\lambda<\lambda_{0}$ it holds $h_{\lambda} \equiv 0$ on the entire connected set $\Omega_{\lambda}$. In particular, for any $\lambda<\lambda_{0}$ and $\underline{r} \in$ $(0,+\infty)^{E}$ such that $\lambda<-3 \sum_{e \in E} r(e)-1$, it holds $\partial_{\lambda}\left(h_{+} / h_{-}\right)(\lambda, \underline{r})=0$ and therefore the ratio $h_{+}(\lambda, \underline{r}) / h_{-}(\lambda, \underline{r})$ depends only on $\underline{r}$.

We now show that this is in contradiction with the assumption that the fundamental graph $G$ is not $(\underline{v}, \bar{v})$ minimal. Indeed, since $G$ is not $(\underline{v}, \bar{v})$-minimal, there exist at least two distinct paths $\gamma^{(1)}=\left(z_{0}, z_{1}, \ldots, z_{M}\right)$ and $\gamma^{(2)}=\left(z_{0}^{\prime}, z_{1}^{\prime}, \ldots, z_{M^{\prime}}^{\prime}\right)$ in $\mathcal{A}_{M}$ and $\mathcal{A}_{M^{\prime}}$ respectively, such that the points $z_{i}$ are all distinct and the points $z_{i}^{\prime}$ are all distinct. It is simple to check that we can always reduce to the following case: $\gamma^{(1)}$ and $\gamma^{(2)}$ have in common the first $\kappa_{1}+1$ points and the last $\kappa_{2}+1$ points (for some non-negative integers $\kappa_{1}$, $\kappa_{2}$ with $\kappa_{1}+\kappa_{2}+2 \leq M \wedge M^{\prime}$ ), while they divide in their interior part (see Figure 5).

We define $\Gamma_{1}$ and $\Gamma_{2}$ as the edge set of $\gamma^{(1)}$ and $\gamma^{(2)}$, i.e.

$$
\Gamma_{1}:=\left\{\left(z_{i}, z_{i+1}\right),\left(z_{i+1}, z_{i}\right): 0 \leq i<M\right\}, \quad \Gamma_{2}:=\left\{\left(z_{j}^{\prime}, z_{j+1}^{\prime}\right),\left(z_{j+1}^{\prime}, z_{j}^{\prime}\right): 0 \leq j<M^{\prime}\right\} .
$$

We also set $E_{\star}:=\Gamma_{1} \cup \Gamma_{2}$ and $V_{\star}:=\left\{z_{1}, \ldots, z_{M}\right\} \cup\left\{z_{1}^{\prime}, \ldots, z_{M^{\prime}}^{\prime}\right\}$. As explained in Section 9.1.1, from the above results and tuning to zero the rates on edges in $E \backslash E_{\star}$ one easily gets the following result:

Claim 9.6. Fix $\underline{r}^{\star} \in(0,+\infty)^{E_{\star}}$ and define $\tilde{f}_{ \pm}^{\star}$ as in (59) referred now to the fundamental graph $G_{\star}=\left(V_{\star}, E_{\star}\right)$ and weights $E_{\star} \ni e \rightarrow r^{\star}(e) \in(0,+\infty)$. Then the ratio $\tilde{f}_{+}^{\star}(\lambda) / \tilde{f}_{-}^{\star}(\lambda)$ does not depend on $\lambda<0$ small enough.


Fig. 5. Example of chains $\gamma^{(1)}, \gamma^{(2)}$. Note that they share the first 4 points $\left(\kappa_{1}=3\right)$ and the last 3 points $\left(\kappa_{2}=2\right)$.

At this point it is enough to exhibit $\underline{r}^{\star} \in(0,+\infty)^{E_{\star}}$ violating the conclusion of the last claim. In order to avoid heavy notation, and without loss of generality, from now on we assume that $G=(V, E)$ equals $G_{\star}=\left(V_{\star}, E_{\star}\right)$, i.e. $G$ is simply the graph in Figure 5.

For each path $\gamma \in \mathcal{A}_{m}$ let $J_{(x, y)}(\gamma)$ denote the current of $\gamma$ along $(x, y)$, i.e. $J_{(x, y)}(\gamma)$ equals the number of times the oriented edge $(x, y)$ appears in $\gamma$ minus the number of times the oriented edge $(y, x)$ appears in $\gamma$. Note that, given $\gamma \in \mathcal{A}_{m}$ and a site $x \neq \underline{v}, \bar{v}$, the number of times that $\gamma$ arrives at $x$ equals the number of times that $\gamma$ leaves $x$. Since moreover $\gamma$ leaves $\underline{v}$ exactly once and it never arrives at $\underline{v}$, and it does the opposite with $\bar{v}$, we get that $J_{z_{i}, z_{i+1}}(\gamma)=1$ for any edge $\left(z_{i}, z_{i+1}\right)$ in $\Gamma_{1} \cap \Gamma_{2}$. In addition it holds either

$$
\begin{cases}J_{\left(z_{i}, z_{i+1}\right)}(\gamma)=k+1 & \forall\left(z_{i}, z_{i+1}\right) \in \Gamma_{1} \backslash \Gamma_{2}  \tag{70}\\ J_{\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right)}(\gamma)=-k & \forall\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right) \in \Gamma_{2} \backslash \Gamma_{1}\end{cases}
$$

for some $k \geq 0$, or

$$
\begin{cases}J_{\left(z_{i}, z_{i+1}\right)}(\gamma)=-k & \forall\left(z_{i}, z_{i+1}\right) \in \Gamma_{1} \backslash \Gamma_{2}  \tag{71}\\ J_{\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right)}(\gamma)=k+1 & \forall\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right) \in \Gamma_{2} \backslash \Gamma_{1},\end{cases}
$$

for some $k \geq 0$. It is simple to check that the above situations can indeed take place.
We can therefore define the two disjoint sets

$$
\begin{aligned}
& \mathcal{P}_{1}:=\left\{\text { paths in } \bigcup_{m \geq 1} \mathcal{A}_{m} \text { such that (70) holds }\right\}, \\
& \mathcal{P}_{2}:=\left\{\text { paths in } \bigcup_{m \geq 1} \mathcal{A}_{m} \text { such that (71) holds }\right\} .
\end{aligned}
$$

Note that $\mathcal{P}_{1} \cup \mathcal{P}_{2}=\bigcup_{m \geq 1} \mathcal{A}_{m}$. Recall the bijection $\mathcal{A}_{m} \ni \gamma \mapsto \gamma^{\dagger} \in \mathcal{A}_{m}^{*}$ defined after (65). Splitting, then, the sum in (63) into the contribution of paths in $\mathcal{P}_{1}$ and $\mathcal{P}_{2}$, we find

$$
\begin{equation*}
\frac{\tilde{f}_{+}(\lambda)}{\tilde{f}_{-}(\lambda)}=\frac{\tilde{f}_{1,+}(\lambda)+\tilde{f}_{2,+}(\lambda)}{\tilde{f}_{1,-}(\lambda)+\tilde{f}_{2,-}(\lambda)} \tag{72}
\end{equation*}
$$

where, for $s=1,2$,

$$
\tilde{f}_{s,+}(\lambda)=\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{P}_{s} \cap \mathcal{A}_{m}} R_{\gamma} S_{\gamma}(\lambda), \quad \tilde{f}_{s,-}(\lambda)=\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{P}_{s} \cap \mathcal{A}_{m}} R_{\gamma^{\dagger}} S_{\gamma^{\dagger}}(\lambda)
$$

We observe now that $G$ has only two loop-free paths from $\underline{v}$ to $\bar{v}: \gamma^{(1)}$ and $\gamma^{(2)}$. By the geometric properties of the loops, it then follows that $\mathcal{P}_{j}$ is given by the paths $\gamma \in \bigcup_{m \geq 1} \mathcal{A}_{m}$ whose loop-erased version equals $\gamma^{(j)}$, for $j=1,2$. Then by (65) and (66) we have $S_{\gamma}(\lambda)=S_{\gamma^{\dagger}}(\lambda)$ and, for $\gamma \in \mathcal{P}_{s} \cap \mathcal{A}_{m}, R_{\gamma}=R_{\gamma^{\dagger}} \Delta_{s}$, where

$$
\begin{equation*}
\Delta_{1}:=\prod_{i=0}^{M-1} \frac{r\left(z_{i}, z_{i+1}\right)}{r\left(z_{i+1}, z_{i}\right)}, \quad \Delta_{2}:=\prod_{i=0}^{M^{\prime}-1} \frac{r\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right)}{r\left(z_{i+1}^{\prime}, z_{i}^{\prime}\right)} . \tag{73}
\end{equation*}
$$

It follows that $\tilde{f}_{s,+}(\lambda)=\tilde{f}_{s,-}(\lambda) \cdot \Delta_{s}, s=1,2$. Combining (72) and (73) we have

$$
\begin{equation*}
\frac{\tilde{f}_{+}(\lambda)}{\tilde{f}_{-}(\lambda)}=\frac{\tilde{f}_{1,+}(\lambda)+\tilde{f}_{2,+}(\lambda)}{\Delta_{1}^{-1} \tilde{f}_{1,+}(\lambda)+\Delta_{2}^{-1} \tilde{f}_{2,+}(\lambda)} \tag{74}
\end{equation*}
$$

If $\Delta_{1}=\Delta_{2}$ then the ratio $\tilde{f}_{+}(\lambda) / \tilde{f}_{-}(\lambda)$ is independent of $\lambda$, but $\Delta_{1}=\Delta_{2}$ holds only for a set of rates of Lebesgue measure 0 . Assume now $\Delta_{1} \neq \Delta_{2}$. Then dividing by $\tilde{f}_{1,+}(\lambda)$ we conclude that the l.h.s. of (74) does not depend on $\lambda$ for $\lambda<0$ with $|\lambda|$ large if and only if the same holds for the ratio $\tilde{f}_{2,+}(\lambda) / \tilde{f}_{1,+}(\lambda)$.

We point out that each path in $\mathcal{P}_{1}$ belongs to $\mathcal{A}_{m}$ for some $m \geq M$. The only path in $\mathcal{P}_{1}$ belonging to $\mathcal{A}_{M}$ is $\gamma^{(1)}$, while there is no path in $\mathcal{P}_{1}$ belonging to $\mathcal{A}_{M+1}$. Moreover, $\left|\mathcal{A}_{m}\right| \leq 3^{m}$ since, when constructing a path $\gamma \in \mathcal{A}_{m}$ vertex by vertex, at each step we can choose only among 2 or 3 neighbors. Hence

$$
\begin{equation*}
\left|\sum_{m=M+2}^{\infty} \sum_{\gamma \in \mathcal{P}_{1} \cap \mathcal{A}_{m}} R_{\gamma} S_{\gamma}(\lambda)\right| \leq \sum_{m=M+2}^{\infty} \frac{c^{m}}{|\lambda|^{m-1}}=\frac{c}{1-c /|\lambda|} \cdot\left(\frac{c}{|\lambda|}\right)^{M+1}, \tag{75}
\end{equation*}
$$

where $c:=3 \max \{r(e): e \in E\}$, and $|\lambda|$ is large enough so that $c /|\lambda|<1$. In particular, isolating the contribution of $\gamma^{(1)}$ in the definition of $\tilde{f}_{1,+}(\lambda)$, we have that

$$
\tilde{f}_{1,+}(\lambda)=c_{1} \prod_{i=1}^{M-1} \frac{1}{r\left(z_{i}\right)-\lambda}+O\left(\frac{1}{|\lambda|^{M+1}}\right), \quad c_{1}:=\prod_{i=0}^{M-1} r\left(z_{i}, z_{i+1}\right) .
$$

Above $O\left(\frac{1}{|\lambda|^{M+1}}\right)$ means that the term in consideration is bounded in modulus by $C /|\lambda|^{M+1}$ for some constant $C>0$ as $\lambda \rightarrow-\infty$. Note that for $\lambda<0$ with $|\lambda|$ large we have

$$
\frac{1}{r\left(z_{i}\right)-\lambda}=\frac{1}{|\lambda|} \cdot \frac{1}{1+r\left(z_{i}\right) /|\lambda|}=\frac{1}{|\lambda|}\left(1-\frac{r\left(z_{i}\right)}{|\lambda|}+\mathcal{E}_{i}(\lambda)\right),
$$

where $\lim _{\lambda \rightarrow-\infty}|\lambda| \mathcal{E}_{i}(\lambda)=0$. The same arguments hold for $\tilde{f}_{2,+}$ with $c_{2}:=\prod_{i=0}^{M^{\prime}-1} r\left(z_{i}^{\prime}, z_{i+1}^{\prime}\right)$. In conclusion we have

$$
\begin{aligned}
& \tilde{f}_{1,+}(\lambda)=\frac{c_{1}}{|\lambda|^{M-1}}-\frac{c_{1}}{|\lambda|^{M}} \sum_{i=1}^{M-1} r\left(z_{i}\right)+o\left(\frac{1}{|\lambda|^{M}}\right), \\
& \tilde{f}_{2,+}(\lambda)=\frac{c_{2}}{|\lambda|^{M^{\prime}-1}}-\frac{c_{2}}{|\lambda|^{M^{\prime}}} \sum_{i=1}^{M^{\prime}-1} r\left(z_{i}^{\prime}\right)+o\left(\frac{1}{|\lambda|^{M^{\prime}}}\right) .
\end{aligned}
$$

If $\tilde{f}_{1,+}(\lambda), \tilde{f}_{2,+}(\lambda)$ are proportional for $\lambda<0$ with $|\lambda|$ large enough, it must be $M=M^{\prime}$ and $\sum_{i=1}^{M-1} r\left(z_{i}\right)=$ $\sum_{i=1}^{M^{\prime}-1} r\left(z_{i}^{\prime}\right)$. These identities cannot be true in general. If $M \neq M^{\prime}$ trivially we have a contradiction. Otherwise take $r(e)=1$ for all $e \in \Gamma_{1}$ and $r(e)=a>0$ for all $a \in \Gamma_{2} \backslash \Gamma_{1}$. If $a$ is large then the identity $\sum_{i=1}^{M-1} r\left(z_{i}\right)=\sum_{i=1}^{M^{\prime}-1} r\left(z_{i}^{\prime}\right)$ fails.

### 9.1.1. Proof of Claims 9.3, 9.4 and 9.6

Proof of Claim 9.3. Given $v \in V$ we introduce the map $\phi_{v}: \mathbb{C}^{E} \rightarrow \mathbb{C}$ as

$$
\phi_{v}(\underline{r}):=\sum_{(v, y) \in E} r(v, y) .
$$

For each $\gamma=\left(x_{0}, x_{1}, \ldots, x_{m}\right) \in \mathcal{A}_{m}, m \geq 1$, we consider the holomorphic function (cf. [14]) $g_{\gamma}: \Omega_{*} \rightarrow \mathbb{C}$ defined as

$$
\begin{equation*}
g_{\gamma}(\lambda, \underline{r}):=R_{\gamma} S_{\gamma}(\lambda)=\prod_{e \in E} r(e)^{N_{e}(\gamma)} \prod_{i=1}^{m-1} \frac{1}{\phi_{x_{i}}(\underline{r})-\lambda}, \tag{76}
\end{equation*}
$$

where $N_{e}(\gamma)$ counts the number of times the edge $e$ appears along the path $\gamma$.
Fix $\left(\lambda^{*}, \underline{r}^{*}\right) \in \Omega_{*}$. Consider the open subset $U\left(\lambda^{*}, \underline{r}^{*}\right) \subset \Omega_{*}$ given by the vectors $(\lambda, \underline{r}) \in \Omega_{*}$ such that $-\mathfrak{R}\left(\lambda^{*}\right) / \sqrt{2}<-\mathfrak{R}(\lambda)<-\sqrt{2} \mathfrak{R}\left(\lambda^{*}\right)$ and $\mathfrak{R}\left(\underline{r}^{*}\right) / \sqrt{2}<\mathfrak{R}(\underline{r})<\sqrt{2} \mathfrak{R}\left(\underline{r}^{*}\right)$ (componentwise). Trivially, $\left(\lambda^{*}, \underline{r}^{*}\right) \in$
$U\left(\lambda^{*}, \underline{r}^{*}\right)$. If $\gamma \in \mathcal{A}_{m}$ and $(\lambda, \underline{r}) \in U\left(\lambda^{*}, \underline{r}^{*}\right)$ we can bound

$$
\begin{align*}
\left|g_{\gamma}(\lambda, \underline{r})\right| & =\prod_{e \in E}|r(e)|^{N_{e}(\gamma)} \prod_{i=1}^{m-1} \frac{1}{\left|\phi_{x_{i}}(\underline{r})-\lambda\right|} \\
& \leq 2^{m} \prod_{e \in E} \Re(r(e))^{N_{e}(\gamma)} \prod_{i=1}^{m-1} \frac{1}{\phi_{x_{i}}(\Re(r))-\Re(\lambda)} \\
& \leq 4^{m} \prod_{e \in E} \Re\left(r^{*}(e)\right)^{N_{e}(\gamma)} \prod_{i=1}^{m-1} \frac{1}{\phi_{x_{i}}\left(\Re\left(\underline{r}^{*}\right)\right)-\Re\left(\lambda^{*}\right)} \\
& \leq \prod_{e \in E} \Re\left(4 r^{*}(e)\right)^{N_{e}(\gamma)} \prod_{i=1}^{m-1} \frac{1}{\phi_{x_{i}}\left(\Re\left(4 \underline{*}^{*}\right)\right)+1} . \tag{77}
\end{align*}
$$

Indeed the first bound follows from the assumptions in the definition of $\Omega_{*}\left(\right.$ note that $\left.\phi_{x_{i}}(\Re(\underline{r}))-\Re(\lambda)>0\right)$, the second bound follows from the definition of $U\left(\lambda^{*}, \underline{r}^{*}\right)$, while the last bound follows from assumption (iii) in the definition of $\Omega_{*}$ since we can bound

$$
\begin{equation*}
\phi_{x_{i}}\left(\Re\left(\underline{r}^{*}\right)\right)-\Re\left(\lambda^{*}\right) \geq 4 \phi_{x_{i}}\left(\Re\left(\underline{r}^{*}\right)\right)+1=\phi_{x_{i}}\left(\Re\left(4 \underline{r}^{*}\right)\right)+1 . \tag{78}
\end{equation*}
$$

We are now interested to the series of holomorphic functions $\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}} g_{\gamma}(\lambda, \underline{r})$. By (77) for any $(\lambda, \underline{r}) \in$ $U\left(\lambda^{*}, \underline{r}^{*}\right)$ we have

$$
\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}}\left|g_{\gamma}(\lambda, \underline{r})\right| \leq \sum_{m=1}^{\infty} \sum_{\gamma=\left(x_{0}, x_{1}, \ldots, x_{m}\right) \in \mathcal{A}_{m}} \prod_{e \in E} \Re\left(4 r^{*}(e)\right)^{N_{e}(\gamma)} \prod_{i=1}^{m-1} \frac{1}{\phi_{x_{i}}\left(\Re\left(4 \underline{r}^{*}\right)\right)+1}
$$

Comparing with (62), the above r.h.s. multiplied by $\frac{1}{\phi_{\underline{v}}\left(\Re\left(4 r^{*}\right)\right)+1}$ equals the function $\tilde{f}_{+}^{*}(-1)$ with $\tilde{f}_{+}^{*}$ defined as the function $\tilde{f}_{+}$referred to the random walk on $\mathcal{G}$ induced by weights $e \rightarrow \mathfrak{R}\left(4 r^{*}(e)\right)$. We finally observe that the value $\tilde{f}_{+}^{*}(-1)$ is finite by definition of $\tilde{f}_{+}^{*}$ (cf. (59)), hence the above r.h.s. is finite.

Since each compact subset of $\Omega_{*}$ can be covered by the union of a finite family of sets of the form $U\left(\lambda^{*}, \underline{r}^{*}\right)$ we conclude that series $\sum_{m=1}^{\infty} \sum_{\gamma \in \mathcal{A}_{m}} g_{\gamma}(\lambda, \underline{r})$ converges uniformly on compact subsets of $\Omega_{*}$. By a classical theorem in complex analysis (see e.g. [14, Chapter I, Lemma 11]), we conclude that the limiting function $h_{+}$is holomorphic. By construction, it satisfies (68) on $\Omega$. By a similar argument one can treat the function $h_{-}$.

Proof of Claim 9.4. To simplify the notation we set $h:=\partial_{\lambda}\left(h_{+} / h_{-}\right)$. Fix $\underline{r} \in \mathcal{R}$. By Claim 9.3 and the definition of the set $\mathcal{R}$, since $\lambda<0$ if $(\lambda, \underline{r}) \in \Omega$, it holds

$$
\begin{equation*}
h_{+}(\lambda, \underline{r})=e^{\Delta(\underline{r})} h_{-}(\lambda, \underline{r}) \quad \forall \lambda:(\lambda, \underline{r}) \in \Omega . \tag{79}
\end{equation*}
$$

By definition of $\Omega_{*}$, we have $(\lambda, \underline{r}) \in \Omega$ for $\lambda$ small enough (depending on $\underline{r}$ ). In particular, (79) implies that the map $\lambda \mapsto\left[h_{+} / h_{-}\right](\lambda, \underline{r})$ is constant for $\lambda<0$ sufficiently small. As a consequence, the map $\lambda \mapsto h(\lambda, \underline{r})$ is zero for $\lambda<0$ sufficiently small. Since this map is also holomorphic on the open connected set (69), we get that it is zero on all the set (69) (see [3, Chapter IV.2.3]).

Proof of Claim 9.6. Given a positive integer $k$, define $\underline{r}_{k} \in(0,+\infty)^{E}$ as $r_{k}(e):=r^{\star}(e)$ for $e \in E_{\star}$ and $r_{k}(e):=1 / k$ for $e \in E \backslash E_{\star}$. Let $\tilde{f}_{ \pm, k}(\lambda)$ be the function defined in (59), for the Markov chain with weights $\underline{r}_{k}$ referred to the graph $G=(V, E)$. Then, for $\lambda<0$, it holds $\lim _{k \rightarrow \infty} \tilde{f}_{ \pm, k}(\lambda)=\tilde{f} \neq(\lambda)$ since the probability to have a jump along an edge in $E \backslash E_{\star}$ goes to zero as $k \rightarrow \infty$ (use the graphical construction for Markov chains). Now observe that, due to (63) and (68),

$$
\begin{equation*}
\frac{\tilde{f}_{+, k}(\lambda)}{\tilde{f}_{-, k}(\lambda)}=\frac{h_{+}\left(\lambda, \underline{r}_{k}\right)}{h_{-}\left(\lambda, \underline{r}_{k}\right)} . \tag{80}
\end{equation*}
$$

From now on we restrict to $\lambda<\min \left\{\lambda_{0},-1-3 \sum_{e \in E_{\star}} r^{\star}(e)\right\}$ without further mention (recall that the constant $\lambda_{0}$, introduced before Claim 9.5 , depends only on the geometry of the graph $G$ and not on a specific family of rates). Trivially, there exists $k_{0}$ such that, for all the above values $\lambda$, if $k \geq k_{0}$ then $\lambda<\min \left\{\lambda_{0},-1-3 \sum_{e \in E} r_{k}(e)\right\}$. By applying Claim 9.5 we get that, for $k \geq k_{0}$, the ratio in the r.h.s. of (80) is $\lambda$-independent, i.e. it is a constant $c_{k}$ depending only on $k$ and $\underline{r}_{k}$. As a byproduct with the first part of the proof we conclude that $\tilde{f}_{+}^{\star}(\lambda) / \tilde{f}_{-}^{\star}(\lambda)=\lim _{k \rightarrow \infty} c_{k}$, which therefore does not depend on $\lambda$.

## Appendix A: Proposition 6.6 implies (UB)

For completeness, following similar arguments as in [5], we explain how one can deduce from Proposition 6.6 the upper bound (UB) for all $\vartheta \in \mathbb{R}$, assuming $v>0$. Recall that $S_{t}:=\inf \left\{s \geq t: Z_{s} \leq 0\right\}$, and observe that for all $u>0$ it holds

$$
\begin{equation*}
\mathbb{P}\left(\inf _{s \geq t} Z_{s} \leq u t\right) \leq q^{-u t} \mathbb{P}\left(S_{t}<\infty\right), \quad q:=\mathbb{P}\left(w_{1}=-1\right) \tag{81}
\end{equation*}
$$

To prove the above bound observe that one possible way of realizing the event $\left\{\inf _{s \geq t} Z_{s} \leq 0\right\}$ is the following. If $Z_{t}>\lfloor u t\rfloor$ then the process hits $\lfloor u t\rfloor$ after time $t$ and then performs $\lfloor u t\rfloor$ consecutive steps to the left. If $Z_{t} \leq\lfloor u t\rfloor$ then after time $t$ the process performs $\lfloor u t\rfloor$ consecutive steps to the left. In particular we get

$$
\mathbb{P}\left(S_{t}<\infty\right)=\mathbb{P}\left(\inf _{s \geq t} Z_{s} \leq 0\right) \geq \mathbb{P}\left(\inf _{s \geq t} Z_{s} \leq\lfloor u t\rfloor\right) q^{\lfloor u t\rfloor} \geq \mathbb{P}\left(\inf _{s \geq t} Z_{s} \leq u t\right) q^{u t}
$$

From (81) and Proposition 6.6 we readily get (UB) for $\vartheta=0$ :

$$
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(-\varepsilon, \varepsilon)\right) \leq \varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\inf _{s \geq t} Z_{s} \leq \varepsilon t\right) \leq-\varepsilon-I(0) \xrightarrow{\varepsilon \rightarrow 0}-I(0) .
$$

Fix, now, any $\vartheta>0$ and take $\varepsilon$ small enough so that $u:=\vartheta-\varepsilon>0$ and fix $u^{\prime} \in(0, u)$. Let $m$ be any positive integer. Then we have for $t$ large (as we assume)

$$
\begin{aligned}
\mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) & \leq \mathbb{P}\left(Z_{t} \in[u t, u t+2 \varepsilon t]\right)=\mathbb{P}\left(T_{\lfloor u t\rfloor} \leq t, Z_{t} \in[u t, u t+2 \varepsilon t]\right) \\
& \leq \mathbb{P}\left(\frac{T_{\lfloor u t\rfloor}\lfloor u t\rfloor}{} \leq \frac{1}{u^{\prime}}, \inf _{s \geq t} Z_{s} \leq u t+2 \varepsilon t\right) \\
& \leq \sum_{k=1}^{m} \mathbb{P}\left(\frac{T_{\lfloor u t\rfloor}}{\lfloor u t\rfloor} \in\left[\frac{(k-1)}{u^{\prime} m}, \frac{k}{u^{\prime} m}\right]\right) \mathbb{P}\left(\inf _{s \geq t-\frac{k t}{m}} Z_{s} \leq 2 \varepsilon t\right) \\
& \leq q^{-2 \varepsilon t} \sum_{k=1}^{m} \mathbb{P}\left(\frac{T_{\lfloor u t\rfloor}}{\lfloor u t\rfloor} \in\left[\frac{(k-1)}{u^{\prime} m}, \frac{k}{u^{\prime} m}\right]\right) \mathbb{P}\left(S_{t-\frac{k t}{m}}<\infty\right) .
\end{aligned}
$$

We point out that the third inequality above follows from the strong Markov property applied at time $T_{\lfloor u t\rfloor}$ and the fact that the probability $\mathbb{P}\left(\inf _{s \geq t-a} Z_{s} \leq 2 \varepsilon t\right)$ is increasing in $a$. The last inequality follows from (81).

Reasoning as in (49) and using Proposition 6.6, we get for $1 \leq k<m$

$$
\begin{aligned}
\mathbb{P}\left(\frac{T_{\lfloor u t\rfloor}}{\lfloor u t\rfloor} \in\left[\frac{(k-1)}{u^{\prime} m}, \frac{k}{u^{\prime} m}\right]\right) \mathbb{P}\left(S_{t-\frac{k t}{m}}<\infty\right) & \leq e^{t \varepsilon+t u w_{+}\left(\frac{1}{u^{\prime}}, \frac{1}{m u^{\prime}}\right)} e^{-t\left[\frac{k}{m} I\left(\frac{m u^{\prime}}{k}\right)+\left(1-\frac{k}{m}\right) I(0)\right]} \\
& \leq e^{t \varepsilon+t u w_{+}\left(\frac{1}{u^{\prime}}, \frac{1}{m u^{\prime}}\right)} e^{-t I\left(u^{\prime}\right)}
\end{aligned}
$$

where $t$ is taken large enough and $w_{+}$is defined as in (48). Note that the last inequality follows from the convexity of $I$. When $k=m$, on the other hand, $\mathbb{P}\left(S_{t-\frac{k t}{m}}<\infty\right)=\mathbb{P}\left(S_{0}<\infty\right)=1$ and, as in (49), for $t$ large we have

$$
\mathbb{P}\left(\frac{T_{\lfloor u t\rfloor}\lfloor u t\rfloor}{\lfloor u}\left[\frac{(m-1)}{u^{\prime} m}, \frac{1}{u^{\prime}}\right]\right) \leq e^{t \varepsilon+t u w_{+}\left(\frac{1}{u^{\prime}}, \frac{1}{m u^{\prime}}\right)} e^{-t I\left(u^{\prime}\right)} .
$$

Putting all together, we have shown that for any $\varepsilon$ small and $t$ large enough it holds

$$
\mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \leq m \cdot q^{-2 \varepsilon t} \cdot e^{t \varepsilon+t u w_{+}\left(\frac{1}{u^{\prime}}, \frac{1}{m u^{\prime}}\right)} e^{-t I\left(u^{\prime}\right)}
$$

with $u=\vartheta-\varepsilon$, and therefore

$$
\varlimsup_{t \rightarrow \infty} \frac{1}{t} \log \mathbb{P}\left(\frac{Z_{t}}{t} \in(\vartheta-\varepsilon, \vartheta+\varepsilon)\right) \leq-2 \varepsilon \log q+\varepsilon+u w_{+}\left(\frac{1}{u^{\prime}}, \frac{1}{m u^{\prime}}\right)-I\left(u^{\prime}\right) .
$$

Letting, now, $m \rightarrow \infty$ and then $\varepsilon \rightarrow 0$ (so that also $u \rightarrow \vartheta$ ) and taking $u^{\prime} \rightarrow \vartheta$ gives (UB).
The proof of the same bound for $\vartheta<0$ follows by similar arguments.

## Appendix B: Zeros set of holomorphic functions in several variables

The following result is based on the Weierstrass Preparation Theorem.
Lemma B.1. Fix $n \geq 1$ integer. Let $V$ be an open set of $\mathbb{C}^{n}$ such that $U:=V \cap \mathbb{R}^{n}$ is connected. Let $f: V \rightarrow \mathbb{C}$ be an holomorphic function. Then either $f \equiv 0$ on $U$ or the set $\{z \in U: f(z)=0\}$ has zero $n$-dimensional Lebesgue measure.

Proof. Note that $U$ is open. Below Lebesgue measure is considered as $n$-dimensional. It is enough to prove the following claim:

Claim B.2. For any $z \in U$ there is a neighborhood $B_{z}$ of $z$ in $U$ such that the set $\left\{y \in B_{z}: f(y)=0\right\}$ has nonempty open part or has zero Lebesgue measure.

Let us first assume the above claim and show how to conclude. If for all $z \in U$ the set $\left\{y \in B_{z}: f(y)=0\right\}$ has zero Lebesgue measure, then each compact subset $K \subset U$ can be covered by a finite family $B_{z_{1}}, B_{z_{2}}, \ldots, B_{z_{r}}$, thus implying that $\{y \in K: f(y)=0\}$ has zero Lebesgue measure. This trivially leads to the fact that $\{z \in U: f(y)=0\}$ has zero Lebesgue measure. On the other hand, if for some $z \in U$ the set $\left\{y \in B_{z}: f(y)=0\right\}$ has nonempty open part, then the analytic function given by $f$ restricted to $U$ is zero on a ball inside the open connected set $U$ and therefore is zero on all $U$ (see [3, Chapter IV.2.3]).

At this point we only need to prove the above Claim B.2. If $f(z) \neq 0$ then for $B_{z}$ small the set $\left\{y \in B_{z}: f(y)=0\right\}$ is empty and we are done. Suppose that $f(z)=0$ and $f$ is not identically zero around $z$. By Weierstrass preparation theorem [14, Chapter II.B], there exists $\varepsilon>0$ such that for all $y=\left(y_{1}, y_{2}, \ldots, y_{n}\right) \in \mathbb{C}^{n}$ with $\left|y_{i}-z_{i}\right|<\varepsilon$ for all $i$ it holds

$$
\begin{align*}
f(y)= & h(y)\left[\left(y_{n}-z_{n}\right)^{k}+a_{1}\left(y_{1}, \ldots, y_{n-1}\right)\left(y_{n}-z_{n}\right)^{k-1}\right. \\
& \left.+\cdots+a_{k-1}\left(y_{1}, \ldots, y_{n-1}\right)\left(y_{n}-z_{n}\right)+a_{k}\left(y_{1}, \ldots, y_{n-1}\right)\right], \tag{82}
\end{align*}
$$

where $k$ is a suitable nonnegative integer, $a_{1}, \ldots, a_{k}$ are holomorphic functions, and $h$ is a never-zero holomorphic function. It then follows that, fixed $\left(y_{1}, \ldots, y_{n-1}\right)$ with $\left|y_{i}-z_{i}\right|<\varepsilon$, the set $\left\{y_{n} \in \mathbb{C}:\left|y_{n}-z_{n}\right|<\right.$ $\left.\varepsilon, f\left(y_{1}, \ldots, y_{n-1}, y_{n}\right)=0\right\}$ has cardinality at most $k$ (in particular, it has zero Lebesgue measure when intersected with $\mathbb{R}$ ). The thesis follows by taking $B_{z}:=\left\{y \in \mathbb{R}^{n}:\left|y_{i}-z_{i}\right|<\varepsilon\right\}$ and applying Fubini theorem.

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[^0]:    ${ }^{1}$ By random walks we mean stochastic jump processes on a given graph. When we restrict to random walks given by Markov chains [27] (hence with exponential waiting times), we call them Markov random walks

[^1]:    ${ }^{2}$ We use the same terminology of [6], in particular $J$ is good if it has compact level sets.

[^2]:    ${ }^{3}$ Sum is thought in $[0,+\infty)$

[^3]:    ${ }^{4}$ Recall that $\lambda_{c} \geq 0$.

