# Reconstructing a piece of 2-color scenery 

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# Reconstructing a piece of 2-color scenery 

Jüri Lember and Heinrich Matzinger<br>EURANDOM<br>P.O. Box 513-5600 MB Eindhoven, The Netherlands

## Contents

1 Introduction and Result ..... 2
1.1 Introduction ..... 2
1.2 Main notations and assumptions ..... 2
1.3 The theorem ..... 4
1.4 Preview ..... 5
1.4.1 Simplified selection rule ..... 6
1.4.2 Avoiding non-ladder words ..... 7
1.4.3 The names ..... 8
1.4.4 Getting selected ..... 10
1.4.5 Avoiding mistakes ..... 11
1.4.6 Final selection ..... 14
2 Iteration ..... 15
2.1 OK cells ..... 15
2.2 Iterated $g$-functions ..... 16
2.3 Counting blocks ..... 18
2.4 Block at origin ..... 23
3 Reconstruction at level $l_{1}$ ..... 24
3.1 Some definitions ..... 24
3.2 Stopping-time events ..... 25
3.2.1 Right side ..... 25
3.2.2 Left side ..... 26
3.2.3 Attributes ..... 27
3.3 Algorithm ..... 28
3.4 Combinatorics for main theorem ..... 31
3.5 Probabilities for main theorem ..... 32
3.5.1 Scenery-dependent events ..... 32
3.5.2 Random-walk depending events ..... 34
3.6 Tuning parameters ..... 39
3.7 Proof of the main theorem ..... 40
4 Appendix ..... 41
4.1 Proof of Theorem 2.1 ..... 41
4.2 Proof of Proposition ..... 43

## 1 Introduction and Result

### 1.1 Introduction

A (one dimensional) scenery $\xi$ is a coloring of the integers $\mathbb{Z}$ with $C_{0}$ colors $\left\{1, \ldots, C_{0}\right\}$. Two sceneries $\xi, \xi^{\prime}$ are called equivalent, $\xi \approx \xi^{\prime}$, if one of them is obtained from the other by a translation or reflection. Let $(S(t))_{t \geq 0}$ be a recurrent random walk on the integers. Observing the scenery $\xi$ along the path of this random walk, one sees the color $\xi(S(t))$ at time $t$. The scenery reconstruction problem is concerned with trying to retrieve the scenery $\xi$, given only the sequence of observations $\chi:=(\xi(S(t)))_{t \geq 0}$. Quite obviously retrieving a scenery can only work up to equivalence. For an overview about scenery reconstruction we refer the reader to an excellent survey in [13].

The research in scenery reconstruction was first motivated by the work on the properties of the color record $\chi$ by Keane and den Hollander [11], [3]. They investigated the ergodic properties of $\chi$, this study was motivated (among others) by the work of Kalikow [10] and den Hollander, Steif [4] in ergodic theory. In particular, the research on scenery reconstruction started with the scenery distinguishing problem. The question was raised independently by Benjamini and Kesten in [1] and [12] as well as by den Hollander and Keane in [11]. These questions motivated many researchers to work in the areas concerning randomly observed scenery, let us just mention Harris [5], Heicklen [6], Burdzy [2], Hoffman [6], Howard [9], [8], [7], Kesten and Spitzer [14], Levin [17], Lindenstauss [18], Rudolph [6], Pemantle [17], Peres [17].

In [12], Kesten asked whether one can recognize a single defect in a random scenery. In order to provide an answer to this question, Matzinger in his Ph.D. thesis [21] proved a somewhat stronger result: typical sceneries can be reconstructed a.s. up to equivalence. The sceneries in Matzinger's setup are independent uniformly distributed random variables. He showed that almost every scenery can be almost surely reconstructed. In [13], Kesten noticed that Matzinger's proof in [21] heavily relies on the skip-free property of the random walk. He asked whether the result might still hold in the case of a random walk with jumps. Merkl, Matzinger and Loewe in [20] gave a positive answer to Kesten's question under a particular assumption: there are strictly more colors than possible single steps for the random walk.

In the present paper we consider the following problem: can a two-color scenery be reconstructed, if it is observed along a random walk with jumps. Among others, this question was asked by H. Kesten in [13]. It turns out that the two color case $\left(C_{0}=2\right)$ is more difficult than the case investigated by Merkl, Matzinger and Loewe in [20]. Although several arguments in [20] do not use the fact that there are more than two colors, the central idea hopelessly fails in the two-color case. To overcome the problem, the existence of certain test becomes crucial. The aim of the tests is to provide some information about the localization of random walk. As explained later, this kind of information makes the scenery reconstruction possible.
The existence of such kind of test was proved in [15]. This was the first important step towards the whole two-color scenery reconstruction. The present paper provides the second step of two-color scenery reconstruction. We construct an algorithm that, given some general information about the origin (stopping times) as well as a small piece of original scenery, retrieves a (long) piece of scenery with exponentially small error. With this result in hand, one can use the method described in [20] to reconstruct the whole scenery. In the terminology of [20], the constructed algorithm provides the "zag"-procedure of overall scenery reconstruction; in fact, "zag"-procedure is the core of scenery reconstruction. The whole scenery reconstruction shall be given in a follow-up paper.

### 1.2 Main notations and assumptions

We define the main concepts of the paper: scenery, random scenery random walk and observations. Also, some general notations will be introduced.

* Scenery is an element of $\{0,1\}^{\mathbb{Z}}$.

For every $I \subseteq \mathbb{Z}$, the elements of $\{0,1\}^{I}$ are called pieces of scenery. Given a piece of scenery $\phi \in\{0,1\}^{I}$, and a subset $I^{\prime} \subseteq I$, the piece of scenery $(\phi(i))_{i \in I^{\prime}}$ is denoted by $\phi \mid I^{\prime}$.
Two pieces of scenery $\phi \in\{0,1\}^{I}$ and $\phi^{\prime} \in\{0,1\}^{I^{\prime}}$ are equivalent, $\phi \approx \phi^{\prime}$, if $\phi$ is obtained by some translation and reflection of $\phi^{\prime}$, i.e. $I^{\prime}=a I+b$, for some $a \in\{-1,+1\}, b \in \mathbb{Z}$ and $\phi(i)=\phi^{\prime}(a i+b)$, $\forall i \in I$. If $\phi$ is obtained from $\phi^{\prime}$ by translation, i.e. $\phi(i)=\phi^{\prime}(b+i)$, then $\phi$ and $\phi^{\prime}$ are called strongly equivalent, we denote this $\phi \equiv \phi^{\prime}$. If $\phi$ is obtained from $\phi^{\prime}$ by reflection i.e. $\phi(i)=\phi^{\prime}(-i)$, $\forall i \in I$, we write $\phi=\phi^{\prime-}$. By definition, $\phi \sqsubseteq \phi^{\prime}$ means that $\phi \approx \phi^{\prime} \mid J$ for some $J \subseteq I^{\prime}$. If, in addition, the equivalence is strong, we write $\phi \sqsubseteq \phi^{\prime}$. In this case $\phi$ is equal to $\phi^{\prime} \mid J$ up to the translation, only.
For a piece of scenery $\phi \mid[x, y]$, where $[x, y]=(x, \ldots, y) \subset \mathbb{Z}$ is an integer interval, we often write $\phi_{x}^{y}$. If $x=0$, then it is skipped, i.e. $\phi \mid[x, y]$ is written as $\phi^{y}$.

* Random scenery $\xi=\{\xi(z)\}_{z \in \mathbb{Z}}$ is a family of i.i.d. Bernoulli random variables with parameter $1 / 2$. We use $\psi$ for a realization of $\xi$, i.e. a scenery $\psi$ is of random element $\xi$.
The notations defined above is valid for random sceneries. For example, $\xi_{x}^{y}$ stands for random piece of scenery $\xi \mid[x, y]$, $\xi^{y}$ means $\xi \mid[0, y]$ etc. etc.
* In this paper, $S=\{S(t)\}_{t \in \mathbb{N}}$ is a recurrent random walk that visits every integer $z$ with positive probability. We assume $S$ starts at origin, i.e. $S(0)=0$. For a $z \in \mathbb{Z}$ we denote $S_{z}=S+z$. An important assumption is that $S$ has only a finite number of steps ("bounded jumps"). More precisely, we assume that the set $\{z: P(S(1)-S(0)=z)>0\}$ is finite. Throughout this paper we denote

$$
L:=\max \{z: P(S(1)-S(0)=z)>0\}
$$

Thus $L$ stands for length of the maximum jump.
We also define

$$
p_{L}:=P(S(L)-S(0)), \quad p_{\text {min }}:=\min _{i}\{P(S(1)-S(0)=i)>0\}
$$

To simplify some proofs we also assume that $S$ is symmetric (however, we do not believe that the symmetricity is necessary).

* We realize $(\xi, S)$ as canonical projections of $\Omega=\{0,1\}^{\mathbb{Z}} \times \Omega$ endowed with product $\sigma$-algebra and probability measure $B\left(1, \frac{1}{2}\right)^{\mathbb{Z}} \times Q_{o}$, where $\Omega_{2} \subseteq \mathbb{Z}^{\mathbb{N}}$ is the set of all possible paths $S, Q$ denotes the law of $S$ and $B\left(1, \frac{1}{2}\right)$ is the Bernoulli $\frac{1}{2}$-distribution. Hence, the random walk $S$ and scenery $\xi$ are independent. For a fixed scenery $\psi \in\{0,1\}^{\mathbb{Z}}$ (a realization of $\xi$ ), we write $P_{\psi}=\delta_{\psi} \times Q=P(\cdot \mid \xi=\psi)$.
We define the filtrations $\mathcal{F}:=\left(\mathcal{F}_{n}\right)_{n \in \mathbb{N}}$, where $\mathcal{F}_{n}:=\sigma(\xi, S(k): k=0, \ldots, n)$ and $\mathcal{G}:=\left(\mathcal{G}_{n}\right)_{n \in \mathbb{N}}$, where $\mathcal{G}_{n}=\sigma(\chi(1), \ldots \chi(n))$.
* We denote by $\chi$ the observations:

$$
\chi:=\xi(S(0)), \xi(S(1)), \xi(S(2)), \ldots
$$

and we interpret $\chi$ as a random piece of scenery $\{0,1\}^{\mathbb{N}}$, so that $\chi(k):=\xi(S(k))$ for all $k \in \mathbb{N}$.
For any $z \in \mathbb{Z}$, we denote $\chi_{z}(k)=\xi\left(S_{z}(k)\right)$. The notation introduced in connection with sceneries are used with observations; in particular, for time interval $[x, y] \subset \mathbb{N}$ we denote

$$
\chi_{z} \mid[x, y]:=: \chi_{z, x}^{y}:=\left(\chi_{z}(x), \chi_{z}(x+1), \ldots, \chi_{z}(y)\right), \quad \chi_{z}^{y}:=\chi_{z, 0}^{y}, \quad \chi^{y}:=\chi_{0,0}^{y}
$$

* Words are the binary vectors $(w(1), \ldots, w(n)), w(i) \in\{0,1\}, n \in \mathbb{N}$. Formally, words are just the pieces of sceneries $\phi_{1}^{N}$. Therefore, all definitions introduced in connection with sceneries hold for words as well. In particular, two words $w$ and $w^{\prime}$ can be equivalent (requires the same length) or they can satisfy the relation $w \sqsubseteq w^{\prime}$. We shall also use the reflected words $w^{-}$. Hence, for a word $w=\left(w_{1}, \ldots, w_{N}\right)$, $w^{-}=\left(w_{N}, \ldots, w_{1}\right)$.

Let $I=[x, y]$. The piece of scenery $\phi \mid I$ ( where $\phi$ is usually $\xi$ or $\chi$ ) as a mapping consists of domain $I$ as well as from the image. The term "word" is usually used in connection with images only. So, we consider a piece of scenery as a word, if the domain is not important or needs not to be specified (although, formally every word has a domain $(1, \ldots, N))$. Hence, we can state that "the piece $\phi_{x}^{y}$ is the word $w$ ", meaning that the image of $\phi \mid I$ is $w$ or, equivalently, $\psi_{x}^{y} \equiv w$. Depending on $\phi$, we shall call $w$ as the observation- or scenery-word.

### 1.3 The theorem

The aim of the paper is to show that that, for every natural number $l_{1}$ that is big enough, there exists an algorithm $\mathcal{A}^{1}$ which is capable with high probability to reconstruct a finite piece of $\xi$ of length $4 e^{l_{1}}$ around the origin. For that, the algorithm $\mathcal{A}^{1}$ uses first $\exp ^{12 \alpha l_{1}}+1$ observations, $\chi^{12 \alpha l_{1}}$, only. Throughout the paper $\alpha>0$ is a fixed constant that does not depend on $l_{1}$. We need $\alpha$ to be big enough and we specify it in Subsection 3.6. Since $\mathcal{A}^{1}$ is supposed to reconstruct the scenery around the origin, it becomes necessary to get some additional information about the location of $S$ around the origin. In other words, besides the observations, the algorithm $\mathcal{A}^{1}$ should receive some signals telling him that a particular observation was generated when $S$ was sufficiently close to the origin. To get such information, $\mathcal{A}^{1}$ is given $\exp \left(\alpha l_{1}\right)$ $\mathcal{G}$-adapted stopping times $\tau=\left(\tau(1), \ldots, \tau\left(\exp \left(\alpha l_{1}\right)\right)\right)$ as an additional input. The stopping times are assumed to satisfy the conditions:

$$
\begin{equation*}
\tau(k)-\tau(k-1) \geq 2 \exp \left(2 l_{1}\right), \quad k=2,3, \ldots, \exp \left(\alpha l_{1}\right)+1, \quad \text { where } \tau\left(\exp \left(\alpha l_{1}\right)+1\right):=\exp \left[12 \alpha l_{1}\right] \tag{1.1}
\end{equation*}
$$

The aim of $\tau$ is to show when $S$ is at $\operatorname{most} \exp \left(l_{1}\right)$ from origin. Thus, they do well, if the following event holds

$$
E_{\text {stop }}^{1}(\tau):=\left\{|S(\tau(k))| \leq \exp \left(l_{1}\right), \quad k=1, \ldots, \exp \left(\alpha l_{1}\right)\right\}
$$

The condition (1.1) states that all stopping times are sufficiently far from each other and they depend on first $\exp \left(12 \alpha l_{1}\right)$ observation $\chi^{\exp \left[12 \alpha l_{1}\right]}$, only. In particular, for each $\tau(k)$, the algorithm $\mathcal{A}^{1}$ can use $2 \exp \left(2 l_{1}\right)$ observations starting from $\tau(k)$. On $E_{\text {stop }}^{1}(\tau)$, all these observations are generated by $S$ being at most $\exp \left(l_{1}\right)+2 \exp \left(2 l_{1}\right)$ from origin. These are the observations that are actually used by $\mathcal{A}^{1}$. The information provided by $\tau$ is essential for the algorithm $\mathcal{A}^{1}$, which is supposed to work on $E_{\text {stop }}^{1}(\tau)$, only. We shall not define the stopping times in this paper. The construction of $\tau$ such that the probability of $E_{\text {stop }}^{1}(\tau)$ is sufficiently big is the so-called zig-step of overall scenery reconstruction (see Chapter 3 in [20]).

Besides the observations and the stopping times, $\mathcal{A}^{1}$ is given the third input: a (small) piece $\psi^{o}$ of original scenery. Formally, $\psi^{o}=\psi \mid I^{o}$, where $I^{o}$ is an integer interval and $\psi$ is the underlying scenery (the realization of $\xi$.) The length of $\psi^{o}$ (i.e. the length of $I^{o}$ ) is at least $l_{1} c_{1} L$, moreover, we assume $I^{o} \subseteq\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]$. Here $c_{1}$ is a fixed constant not depending on $l_{1}$ (see Section 3.6).

The output of $\mathcal{A}^{1}$ is a word of length $4 \exp \left(l_{1}\right)$. Hence, formally $\mathcal{A}^{1}$ is the mapping

$$
\mathcal{A}^{1}:\{0,1\}^{\left[0, \exp \left(12 \alpha l_{1}\right)\right]} \times\left[0, \exp \left(12 \alpha l_{1}\right)\right]^{\left[1, \exp \left(\alpha l_{1}\right)\right]} \times\left(\bigcup_{k=2 c_{1} l_{1} L+1}^{2 \exp \left(l_{1}\right)+1}\{0,1\}^{k}\right) \mapsto\{0,1\}^{\left[-2 \exp \left(l_{1}\right), 2 \exp \left(l_{1}\right)\right]}
$$

where the first input stands for observations $\chi^{12 \alpha_{1} l_{1}}$, the second for stopping times $\tau$ and the third for $\psi^{o}$.
The aim of $\mathcal{A}^{1}$ is to produce a piece of original scenery that lies between $\psi \mid\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]$ and $\psi \mid\left[-3 \exp \left(l_{1}\right), 3 \exp \left(l_{1}\right)\right]$. Recall that $\psi$ is the realization of $\xi$. Thus, $\mathcal{A}^{1}$ does well, if the following event holds

$$
\begin{equation*}
E_{\text {alg works }}^{1}\left(\tau, I^{o}\right):=\left\{\xi\left|\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right] \sqsubseteq \mathcal{A}^{1}\left(\chi^{\exp \left(12 \alpha l_{1}\right)}, \tau, \xi \mid I^{o}\right) \sqsubseteq \xi\right|\left[-3 \exp \left(l_{1}\right), 3 \exp \left(l_{1}\right)\right]\right\} . \tag{1.2}
\end{equation*}
$$

Obviously the event (1.2) depends on $\tau$ as well as on the chosen interval $I^{o}$. In the following we do not know exactly the interval $I^{o}$. Hence, we want that $\mathcal{A}^{1}$ works with any given interval $I^{o}$. The corresponding event is

$$
E_{\text {alg works }}^{1}(\tau):=\bigcap_{I^{\circ} \subset\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]} E_{\text {alg works }}^{1}\left(\tau, I^{o}\right)
$$

The description and formal definition of $\mathcal{A}^{1}$ is given in Subsection 3.3. The main result of the paper, Theorem 1.1 states that the definition of $\mathcal{A}^{1}$ is successful: given $E_{\text {stop }}^{1}(\tau)$ holds, the conditional probability of $E_{\text {alg works }}^{1}(\tau)$ is big.
Theorem 1.1 There exists a constant $k>0$ not depending on $l_{1}$ such that, for $l_{1}$ big enough

$$
\begin{equation*}
P\left(E_{\text {stop }}^{1}(\tau) \cap\left(E_{\text {alg works }}^{1}(\tau)\right)^{c}\right) \leq e^{-k l_{1}} \tag{1.3}
\end{equation*}
$$

The use of $\tau$ and $\psi^{o}$ might seem unrealistic - one would like to reconstruct (a piece of) scenery without any additional help. In Chapter 3 of [20], a general description of such a scenery reconstruction procedure is given. This procedure is based on repeated use of algorithms $\mathcal{A}^{1}$, where in every stage a longer and longer piece of scenery around origin is constructed ( $l_{1}$ is increasing). In this procedure, the output of $\mathcal{A}^{1}$ in a lower level (for small $l_{1}$ ) is used to define stopping times $\tau$ in higher level (for big $l_{1}$ ) such that with high probability the event $E_{\text {stop }}^{1}(\tau)$ holds. Also the output in lower level is used as an input $\psi^{o}$ for $\mathcal{A}^{1}$ in higher level. In the perspective of such a feedback, the result of the present paper becomes necessary; in fact, this is the core of the overall scenery reconstruction.

### 1.4 Preview

Let us briefly introduce some main ideas behind the construction of $\mathcal{A}^{1}$. We begin with the description of a ladder word. Let $x, y \in \mathbb{Z}$ be two location points such that $y=x+c_{1} l_{1} L$, where $c_{1}$ is a fixed constant, specifies in Section ??. A ladder word $w$ is the piece of observations that $S$ generates by moving from $x$ to $y$ as quickly as possible. Since the length of the maximum step of $S$ is $L$, then for $\xi=\psi$ the described ladder word is obviously the vector

$$
\begin{equation*}
\left.\left(\psi(x), \psi(x+L), \ldots, \psi\left(x+\left(c_{1} l_{1}-1\right) L\right)\right), \psi(y)\right) \tag{1.4}
\end{equation*}
$$

The importance of the ladder words in scenery reconstruction comes form the fact that they can be sometimes recognized (with high probability). Indeed, suppose we "see $x$ and $y$ in $\chi$ ", i.e. looking at the observations, we know exactly when $S$ is in location $x$ and in location $y$. In this case, we can almost surely identify (1.4): just look at all occurrences of $x$ and $y$ in $\chi$ with minimal distances. The words occurring in $\chi$ between $x$ any $y$ are (a.s.) always the same and equal to (1.4). The formal definition of ladder words is given in Section 3.1.
The algorithm $\mathcal{A}^{1}$ consists of two phases. In the first phase, $\mathcal{A}^{1}$ builds a collection of ladder words, $\mathcal{W}^{1}$. For this, we introduce a selection rule: an observation-word $w$ passes the selection and will be collected as a ladder word, if it satisfies certain criterions. In the second phase, $\mathcal{A}^{1}$ assembles the words of $\mathcal{W}^{1}$ to produce a word of length $4 \exp \left(2 l_{1}\right)$ as the output. The assembling-rule of the second phase is straightforward: we start with the given piece $\psi^{o}$, and we attach a ladder word $w \in \mathcal{W}^{1}$ with it only if $w$ has an overlap with $\psi^{o}$ at least $\frac{c_{1} l_{1}}{4}$. Thus, the second phase looks like a puzzle playing. The role of $\psi^{o}$ becomes now obvious $-\psi^{o}$ is the starting piece (the "seed") for our puzzle. For the second phase to works, it is clearly necessary that every ladder word of length $\frac{c_{1} l_{1}}{4}$ occurs only once in $\xi \mid\left[-e^{3 l_{1}}, e^{3 l_{1}}\right]$. It turns out that for $c_{1}$ big enough, the latter holds with high probability (Proposition 3.1). Clealry, it is necessary that $\mathcal{W}^{1}$ contains enough ladder words. On the other hand, for $\mathcal{A}^{1}$ to work, it is also necessary that $\mathcal{W}^{1}$ contains only ladder words. This means that the selection rule for $\mathcal{W}^{1}$ must be balanced - it cannot be neither too strict nor to weak. To construct such a selection rule is the most difficult part of the scenery reconstruction.

### 1.4.1 Simplified selection rule

The selection rule is based on the fact that (with high probability) some location pairs ( $x, y$ ) such that $y=c_{1} l_{1} L+x$ can be seen from observations. This is done by the location tests. Roughly speaking, a location test for $y$ is the procedure that allows us to take decision, whether a particular observation $\chi(t)$ was generated on $y$ (i.e. $S(t)=y$ ) or not. As explained before, with such information in hand, one can easily "collect" the ladder word (1.4).
Let us briefly introduce the main ideas behind the location test for $y$. For tutorial reason, we start with a very unrealistic and oversimplified version of the tests and then, step by step, we approach to the real tests.
Let $\xi=\psi$. We consider a long piece of scenery $\psi \mid[y, y+l m]$, where $l, m$ are sufficiently big constants; and we aim to define a (name) function $g(\psi \mid[y, y+l m])=: g_{y}(\psi)$ as well as a (name reading) function $\hat{g}(w), w \in\{0,1\}^{l m^{2}+1}$ such that the following holds

1 If $S(t) \geq y$, then $\hat{g}\left(\chi \mid\left[t, t+l m^{2}\right]\right)$ is able to reproduce $g_{y}(\xi)$ with certain positive probability;
2 If $S(t)<y$, then the probability that $\hat{g}\left(\chi \mid\left[t, t+l m^{2}\right]\right)$ reproduces $g_{y}(\xi)$ is negligible.
In other words, we try to define the name function $g$ and the name-reader $\hat{g}$ such that $\hat{g}\left(\chi \mid\left[t, t+l m^{2}\right]\right)$ reads $g_{y}(\psi)$ only if the piece of observation $\chi \mid\left[t, t+l m^{2}\right]$ satisfies $S(t) \geq y$.
Similarly, to get a location test for $x$, we define the name function $g^{*}(\psi \mid[x-\operatorname{lm}, x])=: g_{x}^{*}(\psi)$ and the (name reading) function $\hat{g}^{*}(w), w \in\{0,1\}^{l m^{2}+1}$ such that the following holds

1* If $S(t) \leq x$, then $\hat{g}^{*}\left(\chi \mid\left[t-l m^{2}, t\right]\right)$ is able to reproduce $g_{x}^{*}(\xi)$ with certain positive probability;
2* If $S(t)>x$, then the probability that $\hat{g}^{*}\left(\chi \mid\left[t-l m^{2}, t\right]\right)$ reproduces $g_{x}^{*}(\xi)$ is negligible.
It is easy to see that $g^{*}$ and $\hat{g}^{*}$ can be deduced from $g$ and $\hat{g}-$ just define $g^{*}(w):=g\left(w^{-}\right)$and $\hat{g}^{*}(w):=\hat{g}\left(w^{-}\right)$.

Suppose, for a moment, that we have a working location tests for a pair $(x, y)$, with $y=x+c_{1} l_{1} L$. Moreover, suppose that "being able to reproduce" above just means equalities $\hat{g}\left(\chi \mid\left[t, t+l m^{2}\right]\right)=g_{y}(\psi)$, $\hat{g}^{*}\left(\chi \mid\left[t, t+l m^{2}\right]\right)=g_{x}^{*}(\psi)$ and "is negligible" means being zero. In this case, the reconstruction (or collecting) of the word (1.4) is rather straightforward. Indeed, for each $t \geq 0$ define the observation words

$$
\begin{equation*}
w^{1}(t):=\chi\left|[t-l m, t], \quad w^{2}(t):=\chi\right|\left[t, t+c_{1} l_{1}\right], \quad w^{3}(t):=\chi \mid\left[t+c_{1} l_{1}, t+c_{1} l_{1}+l m^{2}\right] \tag{1.5}
\end{equation*}
$$

and apply the name-reading functions $\hat{g}^{*}\left(w^{1}(t)\right)$ and $\hat{g}\left(w^{3}(t)\right)$. Because $S$ is recursive, a.s. there exists a $t$ such that $\hat{g}^{*}\left(w^{1}(t)\right)=g_{x}^{*}(\psi)$ and $\hat{g}\left(w^{3}(t)\right)=g_{y}(\psi)$. In particular, this implies that

$$
\begin{equation*}
S(t) \leq x \quad \text { and } \quad S\left(t+c_{1} l_{1}\right) \geq y \tag{1.6}
\end{equation*}
$$

On the other hand, during $c_{1} l_{1}$ steps, the random walk $S$ cannot move more than $c_{1} l_{1} L$. But this is exactly the distance between $x$ and $y$. Hence, the only possibility for (1.6) to hold is that both inequalities are equalities. In this case, $w^{2}(t)$ equals the ladder word (1.4).

The example above is unrealistic in many respect. It is obvious that a necessary condition for the location test to work is that there is no $z<y$ such that $\psi|[z, z+l m]=\psi|[y, y+l m]$. But from the definition of $\xi$ it follows that for almost all realizations such a $z$ exists (any finite pattern occurs infinitely many times in $\xi$ ). Therefore, it is more realistic to assume that the word $\psi \mid[y, y+l m]$ is unique in a certain piece of $\psi \mid I_{1}$, only. Since we are interested in reconstructing the scenery around the origin, from now on, we define

$$
I_{1}:=\left[-\exp \left(3 l_{1}\right), \exp \left(3 l_{1}\right)\right]
$$

and we consider the pairs $(x, y)$ in $I_{1}$, only. Thus the conditions 2 and $2^{*}$ are replaced by

$$
\begin{array}{cl}
P\left(\hat{g}\left(\chi \mid\left[t, t+l m^{2}\right]\right)=g_{y}(\psi),\right. & \left.S(t) \in\left[-\exp \left(3 l_{1}\right), y\right]\right)=0 \\
P\left(\hat{g}^{*}\left(\chi \mid\left[t-l m^{2}, t\right]\right)=g_{x}^{*}(\psi),\right. & \left.S(t) \in\left[x, \exp \left(3 l_{1}\right)\right]\right)=0 \tag{1.8}
\end{array}
$$

Since the above-described selection rule now works only on $I_{1}$, we have to modify the construction of (2.1) such that $S(t), S\left(t+c_{1} l_{1}\right) \in I_{1}$. For this we use the stopping times $\tau(j)$. Define times

$$
\begin{equation*}
T^{1}(j):=\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}, \quad T^{3}(j):=T^{1}(k)+c_{1} l_{1}, \quad j=1, \ldots, \exp \left(\alpha l_{1}\right) \tag{1.9}
\end{equation*}
$$

Note that on $E_{\text {stop }}(\tau)$ it holds $S\left(T^{1}(j)\right), S\left(T^{3}(j)\right) \in I_{1}$, provided $l_{1}$ is big enough. Now the words defined by $T^{1}(j)$ and $T^{3}(j)$ can be used. More precisely, we define

$$
\begin{aligned}
w^{1}(j) & :=\chi \mid\left[T^{1}(j)-l m^{2}, T^{1}(j)\right] \\
w^{2}(j) & :=\chi \mid\left[T^{1}(j), T^{3}(j)\right] \\
w^{3}(j) & :=\chi \mid\left[T^{3}(j), T^{3}(j)+l m^{2}\right]
\end{aligned}
$$

and we use the same selection rule as previously, with $w^{1}(j), w^{2}(j), w^{3}(j)$ instead of $w^{1}(t), w^{2}(t), w^{3}(t)$. Note that a necessary condition for this rule is that the probability in $\mathbf{1}$ and $\mathbf{1}^{*}$ is so big that among $\exp \left(\alpha l_{1}\right)$ stopping times most likely there is at least one $j$ such that $\hat{g}^{*}\left(w^{1}(j)\right)=g_{x}^{*}(\psi)$ and $\hat{g}\left(w^{3}(j)\right)=$ $g_{y}(\psi)$. Also note that the $T^{1}(j)$ is not defined right after $\tau(j)$, but after $\tau(j)+\exp \left(2 l_{1}\right)$, instead. The reason for this is following: we are interested in reconstructing the a piece of scenery with length $4 \exp \left(l_{1}\right)$ around origin (recall the definition of $E_{\text {alg works }}^{1}$ ). This means that we have to collect also these ladder words that are about $2 \exp \left(l_{1}\right)$ from origin. The stopping times $\tau(j)$ stop $S$ at most $\exp \left(l_{1}\right)$ from origin (on $E_{\text {stop }}(\tau)$ ). Hence, for $S$ to reach to the ladder words that are are about $2 \exp \left(l_{1}\right)$ from origin, some additional time is needed.

The rule in the previous example requires that we know the names $g_{x}^{*}:=g_{x}^{*}(\psi)$ and $g_{y}:=g_{y}(\psi)$. They depend on $\psi$ that is unknown. However, by conditions 1 and $1^{*}$, the names $g_{x}^{*}$ and $g_{y}$ can be red with positive probability. We now modify the selection rule to take into consideration that $g_{x}^{*}$ and $g_{y}$ are not known. The modification is based on the fact that the probability to read $g_{x}^{*}$ and $g_{y}$ is so big that among $\exp \left(\alpha l_{1}\right)$ pairs $\hat{g}^{*}\left(w^{1}(j)\right), g\left(w^{3}(j)\right)$ there is at least $\exp \left(\gamma l_{1}\right)$ pairs such that $\hat{g}^{*}\left(w^{1}(j)=g_{x}^{*}\right.$ and $g\left(w^{3}(j)\right)=g_{y}$ (with high probability, of course). Here $0<\gamma<\alpha$ is a properly chosen proportion. If the latter holds, then there exists a pair of names $g_{1}^{*}, g_{3}$ such that the number of stopping times satisfying $\hat{g}^{*}\left(w^{1}(j)=g_{1}^{*}\right.$ and $g\left(w^{3}(j)\right)=g_{3}$ is more than $\exp \left(\gamma l_{1}\right)$. Unfortunately, there can be many pairs having the same property. To choose the right pair, we riep benefit from the conditions (1.7) and (1.8). Due to these condition, the right pair of names $g_{x}^{*}, g_{y}$ has an important characteristic - for every $j$ such that $\hat{g}^{*}\left(w^{1}(j)\right)=g_{x}^{*}$ and $\hat{g}\left(w^{3}(j)\right)=g_{y}$, the word $w^{2}(j)$ must be (1.4) and, therefore, the same. Our modified rule is the following:

Simplified selection: The word $w$ is taken as (1.4), if there exists a pair of names $g_{1}^{*}, g_{3}$ such that the following holds:
a) there exists more than $\exp \left(\gamma l_{1}\right)$ stopping times such that

$$
\begin{equation*}
\hat{g}^{*}\left(w^{1}(j)\right)=g_{1}^{*}, \quad \hat{g}\left(w^{3}(j)\right)=g_{3} \tag{1.10}
\end{equation*}
$$

b) for every $j$ satisfying (1.10), it holds $w^{2}(j)=w$.

### 1.4.2 Avoiding non-ladder words

In the selection rule above, the right choice of $\gamma$ is crucial: if $\gamma$ is too big, then the probability that the true ladder word passes the criterion a) becomes too small. On the other hand, if $\gamma$ is too small, then
the probability that a non-ladder word passes the selection rule becomes too big. Let us briefly introduce the basic argument used to find a suitable lower bound for $\gamma$.
Suppose $z, z^{\prime} \in I_{1}$ such that $\left|z-z^{\prime}\right|<L c_{1} l_{1}$. Consider the possible observation-words that $S$ generates by going from $z$ to $z^{\prime}$ in $c_{1} l_{1}$ steps. If $c_{1}$ is big enough, then the probability that all these words are the same, is small (Proposition 3.1). In Section 3.1 we define the event $B_{\text {recon straight }}^{1}$ which states that for every $z, z^{\prime} \in I_{1}$ there are at least two possible observation-words that $S$ can generate during its way from $z$ to $z^{\prime}$ with $c_{1} l_{1}$ steps. Any path of $S$ that consists of $c_{1} l_{1}$ steps has the probability at least $\left(p_{\text {min }}\right)^{c_{1} l_{1}}$. Suppose $w$ passes the selection rule. Hence, there exists a set $J \subseteq\left\{1, \ldots, \exp \left(\alpha l_{1}\right)\right\}$ such that at least $|J| \geq \exp \left(\gamma l_{1}\right)$ and for each $j \in J$ the following holds : $\left|S\left(T^{3}(j)\right)-S\left(T^{1}(j)\right)\right|<L l_{1} c_{1}$ and $w^{2}(j)=w$. Let $Y_{k}:=1-I_{w^{2}\left(j_{1}\right)}\left(w^{2}\left(j_{k}\right)\right)$, where $j_{1}, j_{2}, \ldots$ are the elements of $J$. This means that $\sum_{k=2}^{\exp \left(\gamma l_{1}\right)} Y_{k}=0$. Suppose now that $w$ is a non-ladder word. If the event $E_{\text {stop }}^{1} \cap B_{\text {recon straight }}^{1}$ holds, then, for each $k \geq 2$, the probability that $Y_{k}=1$ cannot be smaller than $\left(p_{\min }\right)^{c_{1} l_{1}}$. Given $S\left(T^{1}\left(j_{k}\right)\right)$ and $S\left(T^{3}\left(j_{k}\right)\right)$ the random variables $Y_{k}$ are independent. Now the Höffding's inequality can be used to estimate (see (3.32))

$$
P\left(\sum_{k=2}^{\exp \left(\gamma l_{1}\right)} Y_{k}=0 \mid E_{\text {stop }}^{1} \cap B_{\text {recon straight }}^{1}\right) \leq \exp \left[-2 \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right]
$$

The right side of the previous display is exponentially small in exponentially small quantity of $l_{1}$, if $\gamma>-2 c_{1} \ln p_{\min }$ (see 3.42). Using the obtained bound, it is not hard to see that the probability that a non-ladder word passes the selection rule is exponentially small in $l_{1}$ (Proposition 3.2)..
Note that in the foregoing argument we did not use any properties of $g$ and $\hat{g}$. Hence, the argument applies also for the final selection rule given in Subsection 1.4.6.

### 1.4.3 The names

In this subsection, we explain the nature of the functions $g$ and $\hat{g}$ (recall that $\hat{g}^{*}$ and $g^{*}$ are practically the same). The construction of these function is based on the following theorem proved in [15]
Theorem 1.2 There exists constants $c>0$ (not depending on $n$ ), $N<\infty, m(n)>n$, the maps

$$
\begin{aligned}
& g:\{0,1\}^{m+1} \mapsto\{0,1\}^{n^{2}+1} \\
& \hat{g}:\{0,1\}^{m^{2}+1} \mapsto\{0,1\}^{n^{2}}
\end{aligned}
$$

and the sequence of events $B_{\text {cell_OK }}(n) \in \sigma(\xi(z) \mid z \in[-c m, c m])$ such that:

1) $P\left(B_{\text {cell_OK }}(n)\right) \rightarrow 1$
2) For all $n>N$ and $\psi_{n} \in B_{\text {cell_OK }}(n)$ :

$$
P\left(\hat{g}\left(\chi_{0}^{m^{2}}\right) \sqsubseteq g\left(\psi_{0}^{m}\right) \mid S\left(m^{2}\right)=m, \xi=\psi_{n}\right)>3 / 4 .
$$

3) $g\left(\xi_{0}^{m}\right)$ is an i.i.d. binary vector where the components are Bernoulli with parameter $1 / 2$.
(Note the abuse of notation: in [15] the sign "ฬ" was used instead of " $\sqsubseteq$ ".)
From now on we assume that $n>N$ and $m(n)$ are fixed constant. We specify them in Section 3.6. Theorem 1.2 provides a test that uses $m^{2}$ observations $\chi_{t}^{t+m^{2}}$ to test the hypotheses:

$$
\begin{aligned}
& H_{o}: S(t)=y, \\
& H_{1}: S(t)<y-L m^{2}
\end{aligned}
$$

given $S\left(t+m^{2}\right)=S(t)+m$ and $\xi \in B_{\text {cell_OK }}(n)$. Indeed, it $S(t)<y-L m^{2}$, then $\chi_{t}^{t+m^{2}}$ is independent of $g\left(\xi_{y}^{y+m}\right)$. By the properties of $\xi$,

$$
P\left(\hat{g}\left(\chi_{t}^{t+m^{2}}\right) \sqsubseteq g\left(\xi_{y}^{y+m}\right)\right)=\left(\hat{g}\left(\chi_{t}^{t+m^{2}}\right) \sqsubseteq g\left(\xi_{0}^{m}\right)\right) \leq\left(\frac{1}{2}\right)^{n^{2}-1}
$$

On the other hand, if $\psi \in E_{\mathrm{OK}}$, then conditional on $A:=\left\{\xi \in B_{\text {cell_OK }}(n), S\left(t+m^{2}\right)=m, S(t)=y\right\}$ it holds

$$
P\left(\hat{g}\left(\chi_{t}^{t+m^{2}}\right) \sqsubseteq g\left(\xi_{y}^{y+m}\right) \mid A\right)>\frac{3}{4} .
$$

The functions $g$ and $\hat{g}$ look like the desired name and name-reading procedures. Indeed, there is certainly a positive probability that $\hat{g}\left(w^{3}(j)\right)$ "reproduces" $g(\psi \mid[y, y+m]$, where "reproducing" now means the relation $\hat{g}\left(w^{3}(j)\right) \sqsubseteq g_{y}$ (note that in this case "■" actually means the equality to the first or last bit). On the other hand, the following modification of the (1.7) holds

$$
\begin{equation*}
P\left(\hat{g}\left(\chi \mid\left[t, t+m^{2}\right]\right) \sqsubseteq g(\psi \mid[y, y+m]), \quad S(t) \in\left[-\exp \left(3 l_{1}\right), y-L m^{2}\right)\right)=\left(\frac{1}{2}\right)^{n^{2}-1} \tag{1.11}
\end{equation*}
$$

So, taking $n$ big enough, we can make the right side of (1.11) as small as we want.
Unfortunately, for several reasons, the functions from Theorem 1.2 is not good enough. Recall that we want the mistake (1.3) to be exponentially small in $l_{1}$. The right side of (1.11) does not depend on $l_{1}$. To handle this, we apply Theorem 1.2 repeatedly. This procedure is called iteration and it is the subject of Section 2. Let us briefly introduce the main ideas behind the iteration.

From now on, we define

$$
l:=l_{1} \cdot l_{2}, \quad \text { where } l_{2} \text { is fixed positive integer, specified in Section 3.6. }
$$

We shall apply the functions $g$ and $\hat{g}$ from Theorem $1.2 l$ times consecutively. Let $w=(w(0), \ldots, w(l m)) \in$ $\{0,1\}^{l m+1}$. We define $l$ sub-words, called cells

$$
w_{i}=(w((i-1) m), \cdots, w(i m)), \quad i=1, \ldots, l
$$

Note that $w_{i}$ and $w_{i+1}$ are not disjoint. Using the sub-words $w_{i}$, we naturally extend the definition of $g$ to the words in $\{0,1\}^{l m+1}$. We define

$$
g:\{0,1\}^{l m+1} \mapsto\{0,1\}^{l\left(n^{2}+1\right)}, \quad g(w)=\left(g\left(w_{1}\right), \ldots, g\left(w_{l}\right)\right)
$$

Note that we denote by $g$ the function in Theorem 1.2 as well as its extension (they coincide if $l=1$ ).
Similarly, let $v=\left(v(0), \ldots, v\left(l m^{2}\right)\right) \in\{0,1\}^{l m^{2}+1}$. We define cells

$$
v_{i}=\left(v\left((i-1) m^{2}\right), \ldots, v\left(i m^{2}\right)\right), \quad i=1, \ldots, l .
$$

Using the sub-words $v_{i}$, we extend the definition of $\hat{g}$ to the words in $\{0,1\}^{l m^{2}+1}$. We define

$$
\hat{g}:\{0,1\}^{l m^{2}+1} \mapsto\{0,1\}^{l n^{2}}, \quad \hat{g}(v)=\left(\hat{g}\left(v_{1}\right), \ldots, \hat{g}\left(v_{l}\right)\right)
$$

We now give a more accurate interpretation to the phrase "to reproduce" in the description 1. Since the "name-reading" or "reproducing" procedure is based on Theorem 1.2, it is natural to expect that $\hat{g}\left(\chi \mid\left[t, t+m^{2} l\right]\right)$ reproduces $g(\psi \mid[y, y+m l])$, if the relation $\sqsubseteq$ holds cell-wise, i.e. $\hat{g}\left(\chi \mid\left[t+(i-1) m^{2}, t+i m^{2}\right]\right) \sqsubseteq$ $g(\psi \mid[y+(i-1) m, y+i m])$ for each $i=1, \ldots, l$. Note that Theorem 1.2 gives lower bound to the probability

$$
P_{\psi}\left(\hat{g}\left(\chi \mid\left[t+(i-1) m^{2}, t+i m^{2}\right]\right) \sqsubseteq g(\psi \mid[y+(i-1) m, y+i m])\right),
$$

only if the piece of scenery $\psi \mid[y+(i-1) m-c m, y+(i-1) m+c m]$ belongs to the set $E_{\text {cell_ok }}^{n}$. If this is the case, we say that the cell $\psi \mid[y+(i-1) m, y+i m]$ is $O K$.
For each (long) piece of scenery $\psi \mid[y, y+l m]$ we now correspond the index $\operatorname{set} \mathcal{I}(\psi \mid[y, y+l m])=: \mathcal{I}_{y}(\psi) \subset$ $\{1, \ldots, l\}$ of OK-cells. Similarly, we define $\mathcal{I}^{*}(\psi \mid[x-l m, x]):=\mathcal{I}\left((\psi \mid[x-l m, x])^{-}\right)$(the reader should be warned that now we only give a simplified definition of $\mathcal{I}$ and $\mathcal{I}^{*}$; the final definition is given in Section
2.1).

Although $E_{\text {cell_OK }}^{n}$ has the probability close to one, since $l$ is big, we expect a proportion of cells not to be OK, i.e $\mathcal{I}_{y} \neq\{1, \ldots, l\}$. We say that $\psi \mid[y, y+l m]$ is OK, if at least $l(1-3 \epsilon)$ cells are OK, i.e $\left|\mathcal{I}_{y}(\psi)\right| \geq l(1-3 \epsilon)$. We say that $\psi \mid[x-l m, x]$ is $\mathrm{OK}^{*}$, if $(\psi \mid[x-l m, x])^{-}$is OK. Equivalently, $\psi^{-} \mid[-x,-x+l m]$ is OK. We denote by $B_{\text {intervals OK }}^{1}$ the set of sceneries that satisfy: $\psi \|[y, y+l m]$ is OK and $\psi \mid[x-l m, x]$ is $\mathrm{OK}^{*}$ for every pair $(x, y) \in I_{1}$. In particular, if $\psi \in B_{\text {intervals } \mathrm{OK}}^{1}$, then $\left|\mathcal{I}_{y}(\psi)\right|,\left|\mathcal{I}_{x}^{*}(\psi)\right| \geq(1-\epsilon) l$. The proportion $\epsilon$ is chosen such that $P\left(B_{\text {intervals } \mathrm{OK}}^{1}\right)$ is sufficiently big (Theorem 2.1 and the estimation (3.21)).
For not OK cells, the statement 2) of Theorem 1.2 needs not hold, and the cell-wise reproducing might fail. Hence, we relax the requirement of the full cell-wise reproducing to the requirement that the OK cells are reproduced. More formally, for any subset $I \subseteq\{1, \ldots, l\}$, we define $\hat{g}(w) \sqsubseteq_{I} g(v)$, if $\hat{g}\left(w_{i}\right) \sqsubseteq g\left(v_{i}\right)$, $\forall i \in I$. Now we say that $g\left(\chi \mid\left[t, t+m^{2} l\right]\right)$ reproduces $g_{y}(\psi)$, if

$$
g\left(\chi \mid\left[t, t+m^{2} l\right]\right) \sqsubseteq_{\mathcal{I}(\psi)} g_{y}(\psi) .
$$

If $\psi \in B_{\text {intervals OK }}^{1}$, then the latter means that cell-wise reproduction holds for at least $l(1-3 \epsilon)$ cells.

### 1.4.4 Getting selected

Let us now give some insight, how do we show that the probability for a ladder word (1.4) to pass the selection is sufficiently high. What follows, is a simplified version of Proposition 3.2. Let

$$
E_{j}(x, y):=\left\{\begin{array}{c}
S\left(T^{1}(j)-l m^{2}\right)=x-l m \\
S\left(T^{1}(j)\right)=x, S\left(T^{3}(j)\right)=y, \\
\hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{\mathcal{I}_{x}^{*}(\xi)} g_{x}^{*}(\xi), \\
\hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}_{y}(\xi)} g_{y}(\xi)
\end{array}\right\}, \quad Y_{j}:=I_{E_{j}}, \quad j=1, \ldots, e^{\alpha l_{1}} .
$$

Clearly (1.4) passes the selection if

$$
\left\{\sum_{j=1}^{e^{\alpha l_{1}}} Y_{j}>e^{\gamma l_{1}}\right\}
$$

Now, by the Markov property of $S$, for each $\psi$

$$
\begin{aligned}
P_{\psi}\left(Y_{j}=1 \mid E_{\text {stop }}(\tau)\right) & =P_{\psi}\left(S\left(T^{1}(j)-l m^{2}\right)=x-l m \mid E_{\text {stop }}(\tau)\right) \\
& \times P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{\mathcal{I}_{x}^{*}(\psi)} g_{x}^{*}(\psi) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m\right) \\
& \times P_{\psi}\left(S\left(T^{3}(j)\right)=y \mid S\left(T^{1}(j)\right)=x\right) \\
& \times P_{\psi}\left(\hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}_{y}(\psi)} g_{y}(\psi) \mid S\left(T^{3}(j)\right)=y\right) .
\end{aligned}
$$

Recall that $T^{1}(j)-l m^{2}=\tau(j)+\exp \left(2 l_{1}\right)$. By $E_{\text {stop }}(\tau),|S(\tau(j))| \leq \exp \left(l_{1}\right)$. Now, the local central limit theorem (LCLT) can be used to see that for $l_{1}$ big enough

$$
P_{\psi}\left(S\left(\tau(j)+e^{2 l_{1}}\right)=x-\operatorname{lm} \mid E_{\text {stop }}(\tau)\right) \geq \exp \left(-1.5 l_{1}\right)
$$

By the definitions of $w^{1}(j), \hat{g}^{*}$ and $\mathcal{I}^{*}$, we have
$P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{\mathcal{I}_{x}^{*}(\psi)} g_{x}^{*}(\psi) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m\right)=$
$P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{g}^{*}\left(\chi \mid\left[T^{1}(j)-l m^{2}, T^{1}(j)\right]\right) \sqsubseteq_{\left.\mathcal{I}(\psi \mid[x-l m, x])-g\left((\psi \mid[x-l m, x])^{-}\right) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m\right)=}=\right.$
$P_{\psi}\left(S\left(l m^{2}\right)=x, \hat{g}^{*}\left(\chi_{x-l m} \mid\left[0, l m^{2}\right]\right) \sqsubseteq_{\mathcal{I}\left(\psi^{-} \mid[-x,-x+l m]\right)} g\left(\psi^{-} \mid[-x,-x+l m]\right)\right)=$
$P_{\psi}\left(S\left(l m^{2}\right)=x, \hat{g}\left(\left(\chi_{x-l m} \mid\left[0, l m^{2}\right]\right)^{-}\right) \sqsubseteq_{\mathcal{I}_{-x}\left(\psi^{-}\right)} g_{-x}\left(\psi^{-}\right)\right)$

By symmetricity of $S$, for each set $\mathcal{V} \subseteq\{0,1\}^{l n^{2}}$, we have

$$
P_{\psi}\left(S\left(l m^{2}\right)=x, \hat{g}\left(\left(\chi_{x-l m} \mid\left[0, l m^{2}\right]\right)^{-}\right) \in \mathcal{V}\right)=P_{\psi}\left(S\left(l m^{2}\right)=x-l m, \hat{g}\left(\chi_{x} \mid\left[0, l m^{2}\right]\right) \in \mathcal{V}\right)
$$

The right side of the previous display equals

$$
P_{\psi^{-}}\left(S\left(l m^{2}\right)=-x+l m, \hat{g}\left(\chi_{-x} \mid\left[0, l m^{2}\right]\right) \in \mathcal{V}\right)
$$

Hence,

$$
\begin{aligned}
& P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{\mathcal{I}_{x}^{*}(\psi)} g_{x}^{*}(\psi) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m\right)= \\
& P_{\psi^{-}}\left(S\left(l m^{2}\right)=-x+l m, \hat{g}\left(\chi_{-x} \mid\left[0, l m^{2}\right]\right) \sqsubseteq_{\mathcal{I}_{-x}\left(\psi^{-}\right)} g_{-x}\left(\psi^{-}\right)\right)= \\
& P_{\psi^{-}}\left(S\left(T^{3}(j)+l m^{2}\right)=-x+l m, \hat{g}\left(\chi \mid\left[T^{3}(j), T^{3}(j)+l m^{2}\right]\right) \sqsubseteq_{\mathcal{I}_{-x}\left(\psi^{-}\right)} g_{-x}\left(\psi^{-}\right) \mid S\left(T^{3}(j)=-x\right)\right)= \\
& P_{\psi^{-}}\left(S\left(T^{3}(j)+l m^{2}\right)=-x+l m, \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}_{-x}\left(\psi^{-}\right)} g_{-x}\left(\psi^{-}\right) \mid S\left(T^{3}(j)=-x\right)\right) .
\end{aligned}
$$

Suppose $\psi \in B_{\text {intervals ok }}^{1}$. Then the probability in the previous display has the lower bound

$$
\begin{equation*}
\inf _{\psi: \psi \mid[y, y+l m] \text { is OK }} P_{\psi}\left(S\left(T^{3}(j)+l m^{2}\right)=y+l m, \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}_{y}(\psi)} g_{y}(\psi) \mid S\left(T^{3}(j)\right)=y\right) . \tag{1.12}
\end{equation*}
$$

Indeed, (1.12) does not depend on $y$ any more. It is not very hard to see now that by 2 ) of Theorem 1.2, (1.12) can be bounded below by

$$
\inf _{\psi: \psi \mid[y, y+l m] \text { is OK }} \prod_{i \in I(\psi)} P_{\psi}\left(\hat{g}\left(\chi_{(i-1) m^{2}}^{i m^{2}}\right) \sqsubseteq g\left(\psi_{(i-1) m}^{i m}\right) \mid S\left(i m^{2}\right)=S\left((i-1) m^{2}\right)+m\right) \geq\left(\frac{3}{4}\right)^{l} .
$$

Finally, for every $\psi$,

$$
P_{\psi}\left(S\left(T^{3}(j)\right)=y \mid S\left(T^{1}(j)\right)=x\right)=\left(p_{L}\right)^{c_{1} l_{1}}
$$

Hence, if $\psi \in B_{\text {intervals OK }}^{1}$, we have

$$
\begin{equation*}
P_{\psi}\left(Y_{j}=1 \mid E_{\text {stop }}(\tau)\right) \geq \exp \left(-1.5 l_{1}\right)\left(\frac{3}{4}\right)^{l}\left(p_{L}\right)^{c_{1} l_{1}}\left(\frac{3}{4}\right)^{l}=\exp \left[-\left(1.5-2 \ln \left(\frac{3}{4}\right) l_{2}-c_{1} \ln \left(p_{L}\right)\right) l_{1}\right] \tag{1.13}
\end{equation*}
$$

Conditional on $E_{\text {stop }}$ and $\psi$, the random variables $Y_{j}$ are independent. Using Höffding's inequality, it is now not difficult to show that $\alpha$ and $\gamma$ can be chosen such that

$$
P\left(\sum_{j=1}^{e^{\alpha l_{1}}} Y_{j} \leq e^{\gamma l_{1}}, B_{\text {intervals OK }}^{1} \cap E_{\text {stop }}(\tau)\right)
$$

is exponentially small in $l_{1}$. Since $P\left(B_{\text {intervals OK }}^{1}\right)$ is $\operatorname{big}(3.21)$, we obtain that the the probability of selecting (1.4) is sufficiently big.

### 1.4.5 Avoiding mistakes

In the previous subsections we saw how the selection rule works if "being negligible" in $\mathbf{2}$ means "equal to zero". The latter is unrealistic and cannot be guaranteed. We now modify the selection rule such that the the probability in $\mathbf{2}$ is considerably small in comparison with the (modified version of the) right side of (1.13) (which also goes to zero as $l_{1}$ grows). To explain the meaning of the additional modification, we consider the events

$$
\begin{equation*}
E_{z, I}:=\left\{\forall i \in I \text { we have that } S_{z}(m(i-1))<m(i-1)-L m^{2}\right\}, \quad I \subseteq\{1, \ldots, l\} \tag{1.14}
\end{equation*}
$$

Suppose $E_{z, I}$ holds. Then, for each cell $i \in I$, the random variables $\chi_{z} \mid[(i-1) m, i m]$ and $\xi \mid[(i-1) m, i m]$ are independent. By $\mathbf{3}$ of Theorem 1.2 , we then have $P\left(\chi_{z}|[(i-1) m, i m] \sqsubseteq \xi|[(i-1) m, i m]\right)=(0.5)^{n^{2}-1}$. This implies $P\left(\hat{g}\left(\chi_{z}^{l m^{2}}\right) \sqsubseteq_{I} g_{y}(\xi)\right) \leq(0.5)^{\left(n^{2}-1\right)|I|}$ and, for $l$ big enough the latter yields

$$
\begin{equation*}
P\left(B_{\text {intervals OK }}^{1} \cap\left\{\hat{g}\left(\chi_{z}^{l m^{2}}\right) \sqsubseteq_{\mathcal{I}_{y}(\xi)} g_{y}(\xi)\right\} \cap E_{z, \mathcal{I}(\xi)}\right) \leq \exp [-(0.3 n) l] . \tag{1.15}
\end{equation*}
$$

(Corollary 2.1). Recall that on $B_{\text {intervals ok. }}^{1}$. Since $n$ can be chosen very big, the right side of (1.13) can be as many times bigger than $\exp [-(0.3 n) l]$ as we want. This property together with the fact that $P\left(B_{\text {intervals }}^{1} \mathrm{OK}\right)$ is big makes the selection rule work.

We now define an additional characteristic of $\psi \mid[y, y+l m]$, denoted by $q(\psi \mid[y, y+l m])=: q_{y}(\psi)$, and corresponding "reading function" $\hat{q}(w), w \in\{0,1\}^{l m^{2}+1}$ such that for each $j$, we have

3 If $S\left(T^{3}(j)\right) \geq y$, then $\hat{q}\left(w^{3}(j)\right)$ reproduces $q_{y}(\xi)$ with certain probability,
4 If $S\left(T^{3}(j)\right)<y$, then $\hat{q}\left(w^{3}(j)\right)$ reproduces $q_{y}(\xi)$ only if $E_{z, \mathcal{I}(\xi)}$ holds.
Denote $z=T^{3}(j)$. Note the difference with $\mathbf{1}$ and 2: if $z \geq y$, then $\hat{q}$ and $q$ must fulfill the requirement like 1. Of course, the meaning of "reproduction" is now different, we shall call it $q$-reproduction. For $z<y$, the requirements for $q$ and $\hat{q}$ are different from that one in $\mathbf{2}$ - we do not require that the probability for $q$-reproduction is small. We require instead that the $q$-reproducing always implies $E_{z, \mathcal{I}(\xi)}$. And then, as we just saw, the probability that $\hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}(\xi)} g_{y}(\xi)$ (the $g$-reproduction, in the sequel) is exponentially small (at least for $y=0$, but the case for general $y$ is not different). Hence, we consider $g$ and $q$ together. For a ladder word to be selected, both $q$-and $g$-reproduction must simultaneously hold (for $\exp \left(\gamma l_{1}\right)$ stopping times, as usually). In the case $z \geq y$, the additional requirement obviously reduces the probability (1.13); however, if the $q$-reproduction has a relatively big probability, then the lower bound like (1.13) might still hold. In the case $z<y$, the $q$-reproduction of $q_{y}(\xi)$ (which might hold with rather big probability) implies $E_{z, \mathcal{I}(\xi)}$, and then the probability of $g$-reproduction is very small.

The idea of $q$-reproduction is partially based on the fact that we do not need every ladder word (1.4) with $x, y \in I_{1}$ do be collected. So far, we have not restricted our choice of $x$ ( $y$ is obviously uniquely determined by $x$ ). Now we consider pairs $(x, y)$ that satisfy pair $(x, y)$ that

$$
\begin{aligned}
& \psi(y-L)=\cdots=\psi(y-1) \neq \psi(y)=\cdots=\psi\left(y+m^{3} L\right) \neq \psi\left(y+m^{3} L+1\right)=\cdots=\psi\left(y+m^{3} L+L\right) \\
& \psi(x+L)=\cdots=\psi(x+1) \neq \psi(x)=\cdots=\psi\left(x-m^{3} L\right) \neq \psi\left(x-m^{3} L-1\right)=\cdots=\psi\left(x-m^{3} L-L\right) .
\end{aligned}
$$

Such pairs are called a barriers. The barriers are random, they depend on $\xi$. The event $B_{\text {enough barriers }}^{1}$, formally defined in Section 3.1 states that we have sufficiently many barriers. In Proposition 3.1 we show that this event has high probability if $l_{1}$ is big enough.
To the end of this section we assume $y=0$ and we skip $y$ from the notation.
Let $\psi \mid\left[\left(2 L m^{2}-1\right) m,\left(2 L m^{2}\right) m\right]$ be the first OK cell of $\psi$. In terms of cell indexes, $2 L m^{2}=i_{1}:=\min \mathcal{I}(\psi)$. Let $z<y$. We consider now the random walk $S_{z}$, and we want to be able to see from the observations $\chi_{z} \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ whether $S_{z}\left(\left(i_{1}-1\right) m^{2}\right)<\left(i_{1}-1\right) m-L m^{2}$, i.e. $E_{z, i_{1}}$ holds. The number $m(n)$ is certainly so big that $\left(2 L m^{2}-1\right) m-L m^{2}>L m^{3}$. Hence $E_{z, i_{1}}$ holds, if $S_{z}\left(\left(i_{1}-1\right) m^{2}\right) \leq m^{3} L$. The latter obviously holds $S_{z}(t) \leq m^{3} L \forall t \leq\left(i_{1}-1\right) m^{2}$, which, in turn, holds if the observation-word $\chi_{z} \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ has the following property: $\chi_{z} \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ does not contain at least $m^{3}$ consecutive same colors followed by the different color. Indeed, in order to reach a point $z^{\prime}>m^{3} L$, the random walk $S_{z}$ must generate at least $m^{3}$ consecutive same-color observations and then at least one observation of the other color.
Hence, when $\chi_{z} \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ satisfies the mentioned condition, we can be sure that $S_{z}(t) \leq m^{3} L$ $\forall t \leq\left(i_{1}-1\right) m^{2}$, i.e. $E_{z, i_{1}}$ holds. If the condition is not met, then the word $\chi_{z} \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ is not considered for $g$-reproduction, it will be filtered out.
Suppose now $z=0$. In this case we want that $\chi_{z} \mid\left[0,\left(i_{1}-1\right) m^{2}\right]=\left(i_{1}-1\right) m$. This gives a big chance
for $g$-reproduction of the $i_{1}$-th cell $\hat{g}\left(\chi \mid\left[\left(i_{1}-1\right) m^{2}, i_{1} m^{2}\right]\right) \sqsubseteq g\left(\psi \mid\left[\left(i_{1}-1\right) m, i_{1} m\right]\right)$. But in this case the observation-word $\chi \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ definitely contains $m^{3}$ consecutive same colors followed by the different color and such a word will be filtered out. Therefore, we must adjust the described condition to make sure that (with certain probability) the word $\chi \mid\left[0,\left(i_{1}-1\right) \mathrm{m}^{2}\right]$ will be not filtered out. For this note: in order to reach from $z<0$ to $z^{\prime}>m^{3} L$, the random walk must generate (in the observations) at least $m^{3}$ consecutive same colors, having the different color at the beginning and at the end. On the other hand, to reach from 0 to $z^{\prime}>m^{3} L$, the random walk can follow the path that begins with $m^{3}$ same colors, hence the word $\chi \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ will not necessarily contain least $m^{3}$ consecutive same colors with the different color in the beginning (although this event has probability bigger than $\frac{1}{2}$ ).
A word $(w(0), w(1), \ldots, w(u-1), w(u))$ is called block with length $u$, if $w(0) \neq w(1)=\cdots=w(u-1) \neq$ $w(u)$. Hence the filtering rule is: the word $\chi \mid\left[0,\left(i_{1}-1\right) \mathrm{m}^{2}\right]$ will be filtered out, if it contains a block with length at least $m^{3}$. Such blocks are called big.
For each block $B$ in $\psi$, we define the reading length of $B$ as the length of the smallest block that the random walk generates in observations by crossing it. If the length of $B$ is $L m^{3}$, then the reading length of $B$ is roughly $m^{3}$ (see Section 2.3 for the formal definition and examples). Suppose now that $i_{1}>2 L m^{2}$ and there is one block with $B$ the reading length at least $m^{3}$ between $m^{3} L+L$ and $\left(i_{1}-1\right) m-L m^{2}$. Then, to reach $\left(i_{1}-1\right) m$ from $y$, the random walk necessarily generates at least one big block in observation. To reach $\left(i_{1}-1\right) m$ from $z<0$, the random walk necessarily generates at least two big block in observations. Hence, the filtering rule in this case is: $\chi \mid\left[0,\left(i_{1}-1\right) m^{2}\right]$ will be filtered out, if it contains more than one big block.

Generally, we proceed as follows: we define $\mathcal{I}(\psi)$ to be indexes if cells that are not only OK, but have the additional property: if $i \in \mathcal{I}(\psi)$ then $\psi \mid\left[(i-1) m-L m^{2}, i m+l m^{2}\right]$ cannot be a part of any block with reading length at least $m^{3}$ (see Section 2.1). This means that any block $B$ with the reading length at least $m^{3}$ must end before $(i-1) m-L m^{2}$. This makes our $q$-reproduction procedure to work. We call a group of blocks with reading length at least $m^{3}$ a big cluster if the random walk can cross the group by generating only one big block in observations. Note that all big clusters of $\psi \mid[0, l m]$ are located in the pieces of $\psi$ corresponding to the cells $\{1, \ldots, l\} \backslash \mathcal{I}(\psi)=: \mathcal{I}^{c}(\psi)$.
For each $i$ we count all big clusters in $\psi \mid[0, i m]$, for each $i=1,2, \ldots, l$ and we compare them with the big clusters in $\chi_{z} \mid\left[0, i m^{2}\right]$ for each $i$. Formally, e define the functions

$$
q:\{0,1\}^{l m+1} \mapsto \mathbb{N}^{l}, \quad \text { and } \quad \hat{q}:\{0,1\}^{l m^{2}+1} \mapsto \mathbb{N}^{l}
$$

as follows: $q(w)=\left(q_{1}(w), \ldots, q_{l}(w)\right), \hat{q}(v)=\left(\hat{q}_{1}(v), \ldots, \hat{q}_{l}(v)\right)$ where

$$
\begin{aligned}
q_{i}(w) & :=\text { number of big clusters contained in sub-vector }(w(0), \ldots, w(i m)) \\
\hat{q}_{i}(v) & :=\text { number of big blocks contained in sub-vector }\left(v(0), \ldots, v\left(i m^{2}\right)\right) .
\end{aligned}
$$

As usually we define $q^{*}(w):=q\left(w^{-}\right)$and $\hat{q}^{*}(v)=\hat{q}^{*}\left(v^{-}\right)$.
We denote

$$
\hat{q}(v) \leq q(w) \quad\left(\hat{q}^{*}(v) \leq q^{*}(w)\right) \quad \text { if and only if } \quad \hat{q}_{i}(v) \leq q_{i}(w) \quad\left(\hat{q}_{i}^{*}(v) \leq q_{i}^{*}(w)\right) \text { for all } i
$$

Hence, if $\hat{q}\left(\chi_{z} \mid\left[0, l m^{2}\right]\right) \leq q(\psi \mid[0, m l])=: q(\psi)$, then for each $i$, the number of big blocks in $\left.\chi_{z} \mid\left[0, i m^{2}\right]\right)$ is not bigger than the number of big clusters in $\psi \mid[0, m i]$. The foregoing argument shows that in case $z<y$, this implies that $S_{z}$ is always "one cluster-end behind" implying $E_{z, \mathcal{I}(\psi)}$.
If $z=0$, then the observation word $\chi \mid\left[0, l m^{2}\right]$ will be not filtered out if, for each $i \in \mathcal{I}(\psi)$, the $S$ moves from 0 to $(i-1) m$ generating as few big blocks in observations as possible. In Proposition 2.1 we show that this event has the probability bigger than

$$
\left(p_{\min }\right)^{\left|\mathcal{I}^{c}(\psi)\right| m^{2}}
$$

This follows from the observation that this particular event restricts the behavior if $S_{y}$ during its stay on the cells in $\mathcal{I}^{c}(\psi)$, only. The bound on the previous display is big enough to still have the bound like (1.13) (Theorem 2.3).

### 1.4.6 Final selection

We are now ready to define the final version of the selection rule.

* Note, for every $u \in\{0,1\}^{l m+1}, q(u)=\left(q_{1}, \ldots, q_{l}\right)$ is vector, such that $q_{i}=\{0,1, \ldots, l\}, q_{1}=0$ and $q_{i} \leq q_{i+1} \leq q_{i}+1$. Any such vector is called a $q$-vector. Hence, for every $u, q(u)$ and $q^{*}(u)$ are $q$-vectors.

Recall that, for any $u \in\{0,1\}^{l m+1}, g(u)=\left(g_{1}, \ldots, q_{l}\right)$, where $g_{i} \in\{0,1\}^{n^{2}+1}$. Any such word is called a $g$-word. Hence, for each $u, g(u)$ and $g^{*}(u)$ are $g$-words.

In section 2.1 we shall give the formal definition of $\mathcal{I}_{y}(\xi)$ and $\mathcal{I}_{x}^{*}(\xi)$. When $B_{\text {intervals ok }}^{1}$ holds, then $\left|\mathcal{I}_{y}(\xi)\right|,\left|\mathcal{I}_{x}^{*}(\xi)\right| \geq(1-3 \epsilon) l$ for each pair $x, y \in I_{1}$.

* We call $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ a set of attributes, if $I^{*}, I \subset\{1, \ldots, l\},\left|I^{*}\right|,|I| \geq l(1-3 \epsilon(n)), q, q^{*}$ are $q$-vectors and $g^{*}, g$ are $g$-words.

Recall the definition if observation words $w^{1}(j), w^{2}(j), w^{3}(j), j=1, \ldots, \exp \left(\alpha l_{1}\right)$. For each set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ we define the set $J\left(I^{*}, I, q^{*}, q, g^{*}, g\right) \subset\left[1, \exp \left(\alpha l_{1}\right)\right]$ as follows:
$j \in J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ if and only if $j$ satisfies

$$
\begin{equation*}
\hat{q}^{*}\left(w^{1}(j)\right) \leq q^{*}, \quad \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{I^{*}} g^{*}, \quad \hat{q}\left(w^{3}(j)\right) \leq q, \quad \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{I} g \tag{1.16}
\end{equation*}
$$

As described, the selection rule is based on $g$ - and $q$-reproduction, and it consists of two parts - getting selected and avoiding non-ladder words. The principle of the final selection is exactly the same as the one of simplified selection described in Subsection 1.4.1.

With $g$ - and $q$-reproduction, the getting selected part (a)) means that (with high probability) for each $x, y \in I_{1}, y-x=L c_{1} l_{1}$ there exists a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ and at least $\exp \left(\gamma l_{1}\right)$ stopping times $\tau(j)$ with corresponding index set $J(x, y)$ such that for each $j \in J(x, y),(1.16)$ hold and the word $w^{2}(j)$ is the same, say $w$. Hence the first requirement of selection rule is to check whether there exists a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that $\exists J^{\prime} \subset J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that $\left|J^{\prime}\right| \geq \exp \left(\gamma l_{1}\right)$ and $j \mapsto w^{2}(j)$ is constant on $J^{\prime}$. The existence of such set of attributes and index-set $J^{\prime}$ can be easily checked.
The second requirement of the selection rule (b)) is avoiding the non-ladder words. We already know that if $(x, y)$ form a barrier then (with high probability) the vectors $q_{x}^{*}(\xi), q_{y}(\xi)$ and words $g_{x}^{*}(\xi)$ and $g_{y}(\xi)$ cannot be read somewhere else. Hence, if $I^{*}, I, q^{*}, q, g^{*}, g$ found in the first step are indeed $\mathcal{I}_{x}^{*}(\xi), \mathcal{I}_{y}(\xi)$ $q_{x}^{*}(\xi), q_{y}(\xi), g_{x}^{*}(\xi), g_{y}(\xi)$ as we want them to be, and if $w$ is the word to be selected, then the following must hold: whenever there is a stopping time index $j$ satisfying (1.16), then $w^{2}(j)=w$. Thus, the set $J^{\prime}$ must actually be $J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$.

We now give the formal definition of the selection rule.
Definition 1.3 We define the set $\mathcal{W}=\mathcal{W}\left(\chi^{12 \alpha l_{1}}, \tau\right)$ as follows. A word $w \in\{0,1\}^{c_{1} l_{1}+1}$ belongs to $\mathcal{W}$ if and only if there exists a complect of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that the following conditions are satisfied:
a) $\left|J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)\right| \geq \exp \left(\gamma l_{1}\right)$
b) if $j \in J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$, then $w^{2}(j)=w$.

## 2 Iteration

In this Section, we formalize $g$ - and $q$-reproduction, described in Subsection 1.4.3. We begin with the definition of the OK-pieces of scenery, and we prove that a long piece of random scenery is typically OK (Theorem 2.1). In Subsection 2.2, we prove the inequality (1.15) (Theorem 2.2). In Subsection 2.3, we formalize $q$-reproduction and we found a suitable lower bound for (1.12) (Theorem 2.3). This is the main ingredient for obtaining the lower bound (1.13). Finally, in Subsection 2.4 we show how the barriers make the whole name-reading procedure to work.

Throughout the section, $n, m(n)$ and $l>2 L m^{2}$ are fixed integer.

### 2.1 OK cells

In Theorem 1.2 we defined the set $B_{\text {cell_OK }}(n) \in \sigma(\xi(z) \mid z \in[-c m, c m])$ that contains all typical pieces of sceneries in interval $[-c m, c m]$. In this definition, $c>1$ is a fixed integer not depending on $m$. Thus, any word $w \in\{0,1\}^{2 c m+1}$, regarded as a piece of scenery restricted to $[-c m, c m]$ either belongs to $B_{\text {cell_OK }}(n)$ or not. We say that such a word $w$ is completely OK, if $w \in B_{\text {cell_OK }}(n)$.

* Let $w_{1}^{N}:=(w(1), \ldots, w(N)), w(j) \in\{0,1\}$ be a binary word. Consider a sub-word $w_{a}^{a+m}$ of $w$. We say that $w_{a}^{a+m}$ is weak-OK, if $a-c m \geq 1, a+c m \leq N$ and the extension of $w, w_{a-c m}^{a+c m}$ is completely OK. Thus, any word of length $m$ is weak-OK, if it is a certain sub-word of a larger word of length 2 cm that is completely OK.
* Define integer intervals

$$
D_{i}:=:\left[d_{i-1}, d_{i}\right]:=\left(d_{i-1}, \ldots, d_{i}\right), \quad \text { where } \quad d_{i}:=i m, \quad i=1,2 \ldots
$$

Clealy $D_{i}$-s are not disjoint, $D_{i} \cap D_{i+1}=\left\{d_{i}\right\}$. It is also clear that $D_{1} \cup \cdots \cup D_{l}=[0, l m]$.

* Consider the words $w \in\{0,1\}^{l m+1}$. For each such a word we define $l$ sub-words, called cells $w_{1}, \ldots, w_{l}$ as follows:

$$
\begin{equation*}
w_{i} \in\{0,1\}^{m+1}, w_{i}:=w_{d_{i-1}}^{d_{i}}=\left(w\left(d_{i-1}\right), \ldots, w\left(d_{i}\right)\right), \quad i=1, \ldots, l \tag{2.1}
\end{equation*}
$$

Hence, when speaking about a cell $w_{i}$, we always consider it as a sub-word of a longer word $w$ with the length $l m$. Regarding $w$ as a mapping, we equivalently define $w_{i}=w \mid D_{i}$.

* Using the representation (2.1) we define the sets of indexes

$$
\mathcal{I}_{I}(w):=\left\{i \in\left[2 L m^{2}, l\right]: w_{i} \text { is weak-OK }\right\} .
$$

Hence $\mathcal{I}_{I}(w)$ is a set of all indexes bigger than $2 L m^{2}$ such that $w_{i}$ is weak-OK.

* We say that binary word $w=(w(1), \ldots, w(N))$ of length at least $N \geq m^{1.1}$ is empty, if there is no index $j$ such that $w(j)=w(j+1)=\cdots=w\left(j+m^{0.9}\right)$. We say that a cell $w_{i}$ has empty neighborhood if $d_{i}+L m^{2} \leq l m, d_{i-1}-L m^{2} \geq 0$ and $\left(w\left(d_{i-1}-L m^{2}\right), \ldots, w\left(d_{i}+L m^{2}\right)\right)$ is empty.
* We say that a word $(w(1), \ldots, w(N))$ contains a fence if $\exists 1 \leq i \leq N-2 L+1$ such that

$$
w(i)=\cdots=w(i+L-1) \neq w(i+L)=\cdots=w(i+2 L-1)
$$

We say that a cell $w_{i}$ in representation (2.1) is isolated, if $L m+2 \leq i \leq l-L m-1$ and both (sub-)words, $w_{i+L m+1}=\left(w\left(d_{i}+L m^{2}\right), \ldots, w\left(d_{i}+L m^{2}+m\right)\right)$ and $w_{i-(L m+1)}=\left(w\left(d_{i-1}-L m^{2}-m\right), \ldots, w\left(d_{i-1}-L m^{2}\right)\right)$ contain a fence.

* Let $w$ be as in (2.1). Define

$$
\begin{aligned}
& \mathcal{I}_{I I}^{1}(w):=\left\{i \in\left[2 L m^{2}, l\right]: w_{i} \text { is isolated }\right\} \\
& \mathcal{I}_{I I}^{2}(w):=\left\{i \in\left[2 L m^{2}, l\right]: w_{i} \text { has empty neighborhood }\right\} \\
& \mathcal{I}_{I I}(w):=\mathcal{I}_{I I}^{1}(w) \cap \mathcal{I}_{I I}^{2}(w), \quad \mathcal{I}(w):=\mathcal{I}_{I}(w) \cap \mathcal{I}_{I I}(w)
\end{aligned}
$$

* Let $\epsilon(n)=: P\left(B_{\text {cell_OK }}(n)^{c}\right) \vee \exp \left(-m^{0.7}\right)$. We know, that $\epsilon(n) \rightarrow 0$. Consider a word $w \in\{0,1\}^{l m+1}$. We say that $w$ is $\mathbf{O K}$ if

$$
\left|\mathcal{I}_{I}(w)\right| \geq l(1-2 \epsilon(n)) \quad \text { and } \quad\left|\mathcal{I}_{I I}(w)\right| \geq l\left(1-\exp \left(-m^{0.7}\right)\right)
$$

Recall the definition $\xi^{m l}:=\xi \mid[0, l m]$ and let us define the events

$$
\begin{aligned}
E_{\mathrm{OK}} & :=\left\{\xi^{m l} \text { is } \mathrm{OK}\right\} \\
E_{\mathrm{OK} a} & :=\left\{\left|\mathcal{I}_{I}\left(\xi^{m l}\right)\right| \geq l(1-2 \epsilon(n))\right\} \\
E_{\mathrm{OK} b} & :=\left\{\left|\mathcal{I}_{I I}\left(\xi^{m l}\right)\right| \geq l\left(1-\exp \left(-m^{0.7}\right)\right)\right\}
\end{aligned}
$$

Clearly,

$$
\begin{equation*}
E_{\mathrm{OK}}=E_{\mathrm{OK} a} \cap E_{\mathrm{OK} b} \tag{2.2}
\end{equation*}
$$

and on $E_{\mathrm{OK}}$

$$
\begin{equation*}
\left|\mathcal{I}\left(\xi^{m l}\right)\right| \geq l(1-3 \epsilon(n)) \tag{2.3}
\end{equation*}
$$

provided $n$ is big enough.
The following theorem states that for $n$ big enough, the probability of $E_{\mathrm{OK}}^{c}$ is exponentially decreasing in $l$. Hence, $E_{\mathrm{OK}}$ represents the typical behavior of $\xi^{m l}$. The proof is based on Höffding's inequalities and we leave it to Appendix.
Theorem 2.1 There exists $N<\infty$ such that for each $n>N$ there exists $a(n)>0$ not depending on $l$ such that for all $l$ big enough the event $E_{\mathrm{OK}}$ is independent on $\xi^{L m^{3}}$ and

$$
P\left(E_{\mathrm{OK}}\right) \geq 1-e^{-a l}
$$

### 2.2 Iterated $g$-functions

Recall the function $g:\{0,1\}^{m+1} \mapsto\{0,1\}^{n^{2}+1}$ and $\hat{g}:\{0,1\}^{m^{2}+1} \mapsto\{0,1\}^{n^{2}}$ from Theorem 1.2. In the present section we extend these definitions to the sets $\{0,1\}^{l m}$ and $\{0,1\}^{l m^{2}+1}$.

* Let $w \in\{0,1\}^{l m+1}$. Using the cell-representation (2.1) we extend the definition of $g$ as follows

$$
\begin{equation*}
g:\{0,1\}^{l m+1} \mapsto\{0,1\}^{l\left(n^{2}+1\right)}, \quad g(w):=\left(g\left(w_{1}\right), g\left(w_{2}\right), \ldots, g\left(w_{l}\right)\right) \tag{2.4}
\end{equation*}
$$

Note: by definition $w_{i}$ and $w_{i+1}$ are not disjoint - they have a common bit. However, by the definition, $g$ does not depend on the first bit. Hence, applied on the scenery $\xi^{m l}$, the components $g_{i}\left(\xi^{m l}\right)$ and $g_{j}\left(\xi^{m l}\right)$ are independent.

* Define intervals

$$
T_{i}:=:\left[t_{i-1}, t_{i}\right]:=\left(t_{i-1}, \ldots, t_{i}\right), \quad \text { where } \quad t_{i}:=i m^{2}, \quad, i=1,2 \ldots
$$

So, $T_{i}$-s are defined as $D_{i}$-s with $m^{2}$ instead of $m$.
Clearly $T_{i}$-s are not disjoint, $T_{i} \cap T_{i+1}=\left\{t_{i}\right\}$. It is also clear that $T_{1} \cup \cdots \cup T_{l}=\left[0, l m^{2}\right]$.

* Consider binary words $v=(v(1), \ldots, v(l m)) \in\{0,1\}^{l m^{2}+1}$. For each such a word we define $l$ sub-words, $v_{1}, \ldots, v_{l}$ as follows:

$$
\begin{equation*}
v_{i} \in\{0,1\}^{m+1}, v_{i}:=v_{t_{i-1}}^{t_{i}}=\left(v\left(t_{i-1}\right), \ldots, v\left(t_{i}\right)\right), \quad i=1, \ldots, l \tag{2.5}
\end{equation*}
$$

Regarding $v$ as a mapping, we equivalently define $v_{i}=v \mid T_{i}$.
Using the sub-words (2.5) we define

$$
\hat{g}:\{0,1\}^{l m+1} \mapsto\{0,1\}^{l n^{2}}, \quad \hat{g}(v):=\left(\hat{g}\left(v_{1}\right), \hat{g}\left(v_{2}\right), \ldots, \hat{g}\left(v_{l}\right)\right)
$$

* Let $A=\left(a_{1}^{\prime}, \ldots, a_{l}^{\prime}\right), B=\left(b_{1}^{\prime}, \ldots, b_{l}^{\prime}\right)$ be $l p$ and $l r$ dimensional vectors, respectively. Let $I \subseteq\{1,2, \ldots, l\}$. We define the following notation:

$$
A \sqsubseteq_{I} B \quad \text { iff for each } i \in I \text { we have that } \quad a_{i}^{\prime} \sqsubseteq b_{i}^{\prime} .
$$

Recall the definition of $E_{z, I}$ in (1.14). The event $E_{z, I}$ says that for each $i \in I$ we have that at time $t_{i-1}$ the random walk $S_{z}$ is further away than $L\left(m^{2}\right)$ from the point $d_{i-1}$. In that case, during the time interval $T_{i}$ the random walk $S_{z}$ can not visit the (location) set $D_{i}$. This, in turn, implies that the observation $\chi_{z} \mid T_{i}$ are independent of $\xi \mid D_{i}$. Then, obviously, $\hat{g}\left(\chi_{z} \mid T_{i}\right)$ is independent of $g\left(\xi \mid D_{i}\right)$.

The following theorem yields the bound (1.15).
Theorem 2.2 There exists $\alpha_{I}(n)>0$ not depending on $l$, such that for all $z<0$ the following holds:

$$
\begin{equation*}
P\binom{\exists I \subset\{1,2, \ldots, l\} \text { with }|I|=l(1-3 \varepsilon(n)) \text { such that }}{E_{z, I} \text { holds and } \hat{g}\left(\chi_{z}^{l m^{2}}\right) \sqsubseteq_{I} g\left(\xi^{m l}\right)} \leq e^{-\alpha_{I} l} \tag{2.6}
\end{equation*}
$$

provided $l$ and $n$ are both big enough.
Proof. Let $z>0$. Denote $\xi_{i}=\xi\left|D_{i}, \chi_{z, i}:=\chi_{z}\right| T_{i}$. Let $Y_{i}, X_{i} i=1, \ldots, l$ be Bernoulli random variables, where

$$
\begin{array}{rll}
X_{i}=1 & \text { iff } & \hat{g}\left(\chi_{z, i}\right) \sqsubseteq g\left(\xi_{i}\right) \\
Y_{i}=1 & \text { iff } & S_{z}\left(t_{i-1}\right)<d_{i}-L m^{2}
\end{array}
$$

By definition, $g\left(\xi_{i}\right)$ is a $n^{2}+1$ dimensional random vector, with elements being Bernoulli iid with parameter $\frac{1}{2}$. For each fixed $n^{2}$-dimensional binary vector $w$ we, therefore, get:

$$
\begin{equation*}
P\left(w \sqsubseteq g\left(\xi_{i}\right)\right)=(0.5)^{n^{2}-1} \tag{2.7}
\end{equation*}
$$

Note, when $\left\{Y_{i}=1\right\}$ holds, then $g\left(\xi_{i}\right)$ is independent of $\hat{g}\left(\chi_{z, i}\right)$. By (2.7) then

$$
P\left(X_{i}=1 \mid Y_{i}=1\right)=P\left(\hat{g}\left(\chi_{z, i}\right) \sqsubseteq g\left(\xi_{i}\right) \mid Y_{i}=1\right)=(0.5)^{n^{2}-1}
$$

Let $I \subset\{1, \ldots, l\}$. Consider the probability $P\left(X_{i}=1, i \in I \mid Y_{i}=1, i \in I\right)$. If $\left\{Y_{i}=1, i \in I\right\}$ holds, then, $\left\{X_{i}, i \in I\right\}$ are iid random variables, with parameter (0.5) ${ }^{n^{2}-1}$. Hence

$$
P\left(X_{i}=1, i \in I \mid Y_{i}=1, i \in I\right)=(0.5)^{\left(n^{2}-1\right)|I|}
$$

Thus, for each $I \subseteq\{1, \ldots, l\}$ we have

$$
\begin{align*}
& P\left(E_{z, I} \cap\left\{\hat{g}\left(\chi_{z}^{l m^{2}}\right) \sqsubseteq_{I} g\left(\xi^{m l}\right)\right\}\right)=E\left(\prod_{i \in I} X_{i} Y_{i}\right)=P\left(\prod_{i \in I} X_{i} Y_{i}=1\right)= \\
& P\left(X_{i}=1, i \in I \mid Y_{i}=1, i \in I\right) P\left(Y_{i}=1, i \in I\right) \leq(0.5)^{\left(n^{2}-1\right)|I|} \tag{2.8}
\end{align*}
$$

Using(2.8), the probability in (2.6) can bound by

$$
\begin{equation*}
\sum_{\substack{I \subset\{1,2, \ldots, l\},|I|=l(1-3 \epsilon(n))}} P\left(E_{z, I} \cap\left\{\hat{g}\left(\chi_{z}^{l m^{2}}\right) \sqsubseteq_{I} g\left(\xi^{m l}\right)\right\}\right) \leq\binom{ l}{3 l \varepsilon(n)}\left(\frac{1}{2}\right)^{\left(n^{2}-1\right) l(1-3 \epsilon(n))} . \tag{2.9}
\end{equation*}
$$

Using Stirling's approximation, one can show that for $l$ big enough

$$
\binom{l}{3 l \varepsilon(n)} \leq \exp \left[-l((3 \varepsilon(n) \ln (3 \varepsilon(n))+(1-3 \varepsilon(n)) \ln (1-3 \varepsilon(n)))]=\exp \left(-l \epsilon_{2}(n)\right)\right.
$$

where $\epsilon_{2}(n):=3 \varepsilon(n) \ln (3 \epsilon(n))+(1-3 \epsilon(n)) \ln (1-3 \epsilon(n)) \rightarrow 0$, as $n$ grows. Hence, if $n$ is big enough, then the sum in (2.9) can bounded by

$$
\exp \left(-l \epsilon_{2}(n)\right)\left((0.5)^{\left(n^{2}-1\right) l(1-3 \epsilon(n))}\right) \leq \exp \left(-\ln ^{2} \frac{\ln 2}{2}\right)=\exp \left(-l \alpha_{I}(n)\right)
$$

where $\alpha_{I}(n)=n^{2} \frac{\ln 2}{2}$.

### 2.3 Counting blocks

We now give the formal definition of block.

* Let $w=(w(u), \ldots, w(v))$ be a word. We say that $w$ is a block, if

$$
w(u) \neq w(u+1)=w(u+2)=\cdots=w(v-1) \neq w(v) .
$$

The length of block is defined as $v-u$. The $w(u)$ and $w(v)$ (or $u$ and $v$ ) are the beginning of the block and the end of the block, respectively. The color $w(u+1)$ is called the color of block. For two blocks, $A=\left[a_{1}, a_{2}\right], B=\left[b_{1}, b_{2}\right]$ we denote $A<B$ if $a_{1}<b_{1}$.
Let $\phi: D \longrightarrow\{0,1\}$ be a piece of scenery. Let $T=\left[t_{1}, t_{2}\right] \subset D$ be an integer interval of length at least 3. Since $\phi \mid T$ can be considered as a word, the definition of block applies to $\phi \mid T$ as well.

For given $\phi$, we also call a location interval $T=\left[t_{1}, t_{2}\right]$ a block of $\phi$, if $\phi \mid T$ is a block (as word). So, in the following, a block can be a certain pattern (word) or a certain location $(T)$, where $\phi$ has a block. In the latter case, speaking about blocks, the piece $\phi$ must be specified. We call a block big if its length is at least $m^{3}$.

Note: although the block basically means many consecutive bits of the same color, by definition the first and last bit of a block must be different. For example, 01110 is a block with length 4 , but 00001 is not a block.

* Let $\left[t_{1}, t_{2}\right] \in \mathbb{N}$ be a (time) interval. We call $R \in \mathbb{Z}^{\left[t_{1}, t_{2}\right]}$ an admissible path of length $t_{2}-t_{1}$, if for all $t \in\left[t_{1}, t_{2}-1\right]$

$$
P(S(1)-S(0)=R(t+1)-R(t))>0
$$

So an admissible path is just a possible trajectory of $S$ in time interval $\left[t_{1}, t_{2}\right]$, starting at $R\left(t_{1}\right)$ and ending at $R\left(t_{2}\right)$. The word "possible" means that the probability of such a trajectory is positive.
Let $\mathcal{R}(n)$ be the set of all admissible paths of length $n$. Thus

$$
\mathcal{R}(n):=\left\{R \in \mathbb{Z}^{[0, n]}: P(S(1)-S(0)=R(i+1)-R(i))>0, i=0, \ldots, n-1\right\}
$$

* Let $B=\left[b_{1}, b_{2}\right] \subset \mathbb{Z}$ be a block of scenery $\psi$. Define

$$
l(B):=\min \left\{n>\left.1\right|^{\exists R \in \mathcal{R}(n) \quad \text { such that } \psi \circ R=\psi(R(0)), \ldots, \psi(R(n))} \begin{array}{c}
\text { is a block, } R(0) \leq b_{1}, R(n) \geq b_{2} \tag{2.10}
\end{array}\right\}
$$

The number $l(B)$ will be called as the reading-length of $B$.
Suppose $l(B)=n$ and $R(0), \ldots, R(n)$ is the admissible path that attains the minimum in (2.10). Then the points $R(0)$ and $R(n)$ are called the reading-beginning and the reading-end of $B$, respectively. The reading length of a block is the length of the smallest block in observations, generated under conditions that $S$ crosses $B$. Clearly, $l(B)$ is approximately $\frac{b_{2}-b_{1}}{L}$, but it depends also on $\psi$ outside the block $B$. Let us consider some examples.
Examples: 1. If $S$ is a simple random walk (i.e. $L=1$ ), then $l(B)=b_{2}-b_{1}$ and reading beginning (reading end) and the beginning (the end) of the block coincide.
2. Let $L=3$. Consider the word $(w(1), w(2), \ldots, w(11))=00111111000$. This word contains a block with the length 7 . The reading length of this block is, obviously, 3 . The beginning of the block is $w(2)$,
the end of the block is $w(9)$. The reading beginning is $w(2)$ or $w(1)$ with the reading ends $w(11)$ or $w(10)$, respectively.
3. Let $L=3$. Consider the word $(w(1), w(2), \ldots, w(11))=0011111111000$. It contains a block of length 9 , the reading length of the block is 3 , the reading beginning of the block is $w(2)$, the reading end of the block is $w(11)$.
4. Suppose $L=4$ and $P(S(1)-S(0)=2)=P(S(1)-S(0)=3)=0$. Consider the word $w(1), \ldots, w(18)=$ 011101111111110111 . This word contains a block of length $10 B=(w(5), \ldots, w(15))$. The reading length of this block is 5 .
5. Change the word without changing the block and consider the word 1110111111111000 . The reading length of $B$ is now 3 , the reading-beginning is $w(4)$, the reading-end is $w(16)$.
6. Consider now the words as in the last 2 examples. Suppose $P(S(1)-S(0)=i)>0, i=-4,-3, \ldots, 3,4$. Then the block has reading length 3 no matter what the neighborhood of the block is.

* Let $A=\left[a_{1}, a_{2}\right], B=\left[b_{1}, b_{2}\right]$ be two blocks of $\psi, A<B$. We say that $A$ and $B$ are connected if they are of same color, say 1 , and there is an admissible path from $A$ to $B$ such that moving along this path, only the color 1 is observed. Formally, $A$ and $B$ is connected, if there exists an $n$ and $R \in \mathcal{R}(n)$ such that $R(0) \in\left(a_{1}, a_{2}\right), R(n) \in\left(b_{1}, b_{2}\right)$ and $\psi \circ R(0)=\psi \circ R(1)=\cdots=\psi \circ R(n)$.
In other words, the blocks of the same color are connected, if it is possible to read them as one block.
Let $B_{1}<B_{2}<\ldots<B_{h}$ be blocks of $\psi$. We say that $B_{1} \cup \ldots \cup B_{h}$ is a big cluster, if
- $B_{i}$ has the reading length at least $m^{3}, i=1, \ldots, h$;
- $B_{1}, \ldots, B_{h}$ are connected;
- there is no more blocks with the reading length at least $m^{3}$ connected to $B_{1}$.

We define the reading-path of a big cluster in the same way as the reading path of a block (which can be a big cluster consisting of one block) - this is the shortest admissible path to cross the big bluster and producing exactly one block. Formally, for a big cluster $C:=B_{1} \cup \ldots \cup B_{h}$ we define the reading length of the big cluster as

$$
l(C):=\min \{n>1: \exists R \in \mathcal{R}(n) \text { such that } \quad \psi(R(0)), \ldots, \psi(R(n)) \text { is a block, } R(0) \leq c, R(n) \geq d\}
$$

where $c$ is the beginning of $B_{1}$ and $d$ is the end of $B_{h}$. These points are referred to as the beginning and the end of $C$, respectively. Clearly, $l(C) \geq m^{3}$. The reading-path of $C$ is any path that attains the minimum above.
${ }^{*}$ Let us fix $\psi \in E_{\mathrm{OK}}$. Denote $\mathcal{I}=\mathcal{I}\left(\psi^{m^{2} l}\right), \mathcal{I}_{I}=\mathcal{I}_{I}\left(\psi^{m^{2} l}\right), \mathcal{I}_{I I}=\mathcal{I}_{I I}\left(\psi^{m^{2} l}\right)$.
Consider the set $\mathcal{I}_{I I}^{c}:=[1, l]-\mathcal{I}_{I I}$. Clearly $\mathcal{I}_{I I}^{c}$ is an union of disjoint intervals, i.e.

$$
\begin{equation*}
\mathcal{I}_{I I}^{c}=\left[l_{1}, l_{2}\right] \cup\left[l_{3}, l_{4}\right] \cup \cdots \cup\left[l_{2 k-1}, l_{2 k}\right], \tag{2.11}
\end{equation*}
$$

where $l_{1}=1, l_{2}, l_{3}, \ldots \in\left[2 \mathrm{Lm}^{2}, l\right], l_{i} \leq l_{i+1}$.
The set of cell-indexes $\left[l_{2 i-1}, l_{2 i}\right]$ corresponds to the location-interval (cells) $\left[\left(l_{2 i-1}-1\right) m, l_{2 i} m\right]$ or $\left[d_{l_{2 i-1}-1}, d_{l_{2 i}}\right]$. Let us denote

$$
\begin{equation*}
r_{i}:=\left(l_{2 i-1}-1\right) m, \quad s_{i}=l_{2 i} m, \quad i=1, \ldots, k \tag{2.12}
\end{equation*}
$$

By definition, $S$ visits every point in $\mathbb{Z}$ i.o.. This means, there exists an integer $k \geq 1$ such that $P(S(k)-S(0)=1)$. Let $\bar{v}:=\inf \{k: P(S(k)-S(0)=1)>0\}$. Thus there is an admissible path $R(0), \ldots, R(\bar{v})$ such that $R(0)=0$ and $R(\bar{v})=1$. Similarly, between points $a<b$ there exists an admissible path $R(0), \ldots, R((b-a) \bar{v})$ such that $R(0)=a, R(\bar{v})=a+1, R(2 \bar{v})=a+2, \ldots R((b-a) \bar{v})=b$. We
say that $S$ moves stepwise from $a$ to $b$, if it moves along the path just described. Obviously, $\bar{v} \ll m$.
In Subsection 1.4.5 we defined big cluster counter $q:\{0,1\}^{l m+1} \mapsto \mathbb{N}^{l}$ and block counter $\hat{q}:\{0,1\}^{l m^{2}+1} \mapsto$ $\mathbb{N}^{l}$.
Define the events

$$
\begin{aligned}
F_{\min }(1):= & \left\{\hat{q}\left(\chi \mid\left[0, s_{1} m\right]\right) \leq q\left(\psi \mid\left[0, s_{1}\right]\right), \chi \mid\left[s_{1}-m \bar{v}, s_{1}\right] \text { contains both colors, } \quad S\left(s_{1} m\right)=s_{1}\right\} . \\
F_{\min }(j):= & \left\{\hat{q}\left(\chi_{r_{j}} \mid\left[0,\left(s_{j}-r_{j}\right) m\right]\right) \leq q\left(\psi \mid\left[r_{j}, s_{j}\right]\right), \chi_{r_{j}} \mid[0, m \bar{v}] \text { and } \chi_{r_{j}} \mid\left[\left(s_{j}-r_{j}\right) m-m \bar{v},\left(s_{j}-r_{j}\right) m\right]\right. \\
& \text { contain both colors } \left.S_{r_{j}}\left(\left(s_{j}-r_{j}\right) m\right)=s_{j}\right\}, j=2, \ldots, k-1
\end{aligned}
$$

For the last interval in (2.12) we define $F(k)$ as $F(j), j>1$, if $s_{k}<l$. If $r_{k}=l$, we define

$$
F_{\min }(k):=\left\{\hat{q}\left(\chi_{r_{k}} \mid\left[0,\left(l-r_{k}\right) m\right]\right) \leq q\left(\psi \mid\left[r_{k}, l\right]\right), \chi_{r_{k}} \mid[0, m \bar{v}] \text { contains both colors, } S_{r_{k}}\left(\left(l-r_{k}\right) m\right)=l\right\}
$$

Obviously, the events $F_{\min }(j)$ depend on the random walk, $S$, only. Moreover, by definition, the event $F_{\min }(j)$ depends on the behavior of the random walk during the time interval $\left[0,\left(s_{j}-r_{j}\right) m\right]$. This means, if for a $j$, there exists at least one admissible path $R_{j} \subset \mathcal{R}\left(\left(s_{j}-r_{j}\right) m\right)$ such that

$$
\begin{align*}
& R_{j}(0)=r_{j}, \quad R_{j}\left(\left(s_{j}-r_{j}\right) m\right)=s_{j}  \tag{2.13}\\
& \hat{q}\left(\psi \circ R_{j}\right) \leq q\left(\psi\left[s_{j}, r_{j}\right]\right)  \tag{2.14}\\
& \text { if } r_{j} \neq 0 \text { and } s_{j} \neq l \\
& \text { then }\left(\psi \circ R_{j}\right) \mid[0, m \bar{v}], \quad \text { and }\left(\psi \circ R_{i}\right) \mid\left[\left(s_{j}-r_{j}\right)-m \bar{v},\left(s_{j}-r_{j}\right)\right] \text { have both colors, } \tag{2.15}
\end{align*}
$$

then $F_{\min }(i) \neq \emptyset$ and $P_{\psi}\left(F_{\min }(i)\right) \geq\left(p_{\min }\right)^{\left(s_{i}-r_{i}\right) m}$. The following proposition, proved in Appendix, shows that for each $i$, at least one such admissible path exists.

Proposition 2.1 For each $i=1, \ldots k$ the following holds:

$$
\begin{equation*}
P_{\psi}\left(F_{\text {min }}(i)\right) \geq\left(p_{\text {min }}\right)^{\left(s_{i}-r_{i}\right) m}=\left(p_{\text {min }}\right)^{\left(l_{2 i}-l_{2 i-1}+1\right) m^{2}} \tag{2.16}
\end{equation*}
$$

The next theorem is the main ingredient of the "getting selected" part of the reconstruction. It gives a lower bound for the probability that $g$ - and $q$-reproduction to work.

Theorem 2.3 There exist constant $\alpha_{I I}(n)>0$ not depending on $l$, such that for all $\psi \in E_{\mathrm{OK}}$ the following holds:

$$
\begin{equation*}
P_{\psi}\left(\hat{g}\left(\chi^{l m^{2}}\right) \sqsubseteq_{\mathcal{I}} g\left(\xi^{m l}\right), \quad \hat{q}\left(\chi^{m^{2} l}\right) \leq q\left(\xi^{m l}\right), \quad S\left(m^{2} l\right)=m l\right) \geq e^{-l \alpha_{I I}} \tag{2.17}
\end{equation*}
$$

Proof. For each subset $i \in[1, l]$ and $I \subset[1, l]$ we define

$$
\begin{aligned}
E_{S}(i) & :=\left\{S\left(t_{i-1}\right)-S\left(t_{i}\right)=m\right\}, \quad E_{S}(I):=\cap_{i \in I} E_{S}(i) \\
E_{\sqsubseteq}(i) & :=\left\{\hat{g}\left(\chi_{z} \mid T_{i}\right) \sqsubseteq g\left(\psi \mid D_{i}\right)\right\}, \quad E_{\sqsubseteq}(I):=\cap_{i \in I} E_{\sqsubseteq}(i) ; \\
E_{\text {no-block }}(i) & :=\left\{\text { the sequence } \chi_{z} \mid T_{i} \text { contains both colors }\right\}, \quad E_{n o-b l o c k}(I):=\cap_{i \in I} E_{n o-b l o c k}(i) .
\end{aligned}
$$

Use $\left[r_{j}, s_{j}\right], j=1, \ldots, k$ as in (2.12) to define

$$
\begin{aligned}
E_{\min }(1):= & \left\{\hat{q}\left(\chi \mid\left[0, s_{1} m\right]\right) \leq q\left(\psi \mid\left[0, s_{1}\right]\right), \quad S\left(s_{1} m\right)=s_{1}\right. \\
& \left.\chi \mid\left[s_{1}-m \bar{v}, s_{1}\right] \text { contain both colors }\right\} \\
E_{\min }(j):= & \left\{\hat{q}\left(\chi \mid\left[r_{j} m, s_{j} m\right]\right) \leq q\left(\psi \mid\left[r_{j}, s_{j}\right]\right), \quad S\left(s_{j} m\right)=s_{j},\right. \\
& \left.\chi \mid\left[r_{j}, r_{j}+m \bar{v}\right] \text { contain both colors, } \chi \mid\left[s_{j}-m \bar{v}, s_{j}\right] \text { contain both colors }\right\}, \\
& j=2, \ldots, k-1 \text { and } \\
E_{\min }:= & \cap_{j=1}^{k} E_{\min }(j) .
\end{aligned}
$$

If $s_{k}=l$, then the requirement $\left\{\chi \mid\left[s_{k}-m \bar{v}, s_{k}\right]\right.$ contain both colors $\}$ is dropped for the definition of $E_{\text {min }}(k)$.

Consider the event $E_{\text {min }} \cap E_{S}\left(\mathcal{I}_{I I}\right)$. Use the relation $E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right) \subset E_{S}([1, l])$ to deduce that

$$
\begin{align*}
P_{\psi}\left(E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right) & =\prod_{j=1}^{k} P_{\psi}\left(E_{\min }(j) \mid S\left(m r_{j}\right)=r_{j}\right) P_{\psi}\left(E_{S}\left(\mathcal{I}_{I I}\right)\right) \\
& =\prod_{j=1}^{k} P_{\psi}\left(F_{\min }(j)\right) P_{\psi}\left(E_{S}\left(\mathcal{I}_{I I}\right)\right) \geq\left(p_{\min }\right)^{\left|\mathcal{I}_{I I}^{c}\right| m^{2}} P_{\psi}\left(E_{S}\left(\mathcal{I}_{I I}\right)\right) \tag{2.18}
\end{align*}
$$

Let $i \in \mathcal{I}_{I I}$. If $i \neq l_{2 j-1}$ for each $j=1, \ldots, k$, then $E_{n o-b l o c k}^{c}(i)$ does not depend on $F_{\min }(j)$. If $i=l_{2 j-1}$ for a $j=1, \ldots, k$, then $P_{\psi}\left(E_{n o-b l o c k}^{c}(i) \mid F_{\min }(j)\right)=P_{\psi}\left(E_{n o-b l o c k}^{c}(i) \mid S\left(t_{i-1}\right)=d_{i-1}\right)$ and $E_{n o-b l o c k}^{c}(i)$ is independent of $F_{\min }\left(j^{\prime}\right), j^{\prime} \neq j$. Hence

$$
P\left(E_{\text {no-block }}^{c}(i) \cap E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)=\prod_{i=1}^{k} P_{\psi}\left(F_{\min }(i)\right) P\left(E_{n o-b l o c k}^{c}(i) \mid S\left(t_{i-1}\right)=d_{i-1}, E_{S}(i)\right) P_{\psi}\left(E_{S}\left(\mathcal{I}_{I I}\right)\right)
$$

implying that, for each $i \in \mathcal{I}_{I I}$

$$
\begin{align*}
& P_{\psi}\left(E_{\text {no-block }}^{c}(i) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)=  \tag{2.19}\\
& P_{\psi}\left(E_{n o-b l o c k}^{c}(i) \mid S\left(t_{i-1}\right)=d_{i-1}, S\left(t_{i}\right)=d_{i}\right)=P_{\psi}\left(E_{n o-b l o c k}^{c}(i) \mid E_{S}([1, l])\right)
\end{align*}
$$

Let us estimate (2.19). If $S\left(t_{i-1}\right)=d_{i-1}$ and $S\left(t_{i}\right)=d_{i}$, then during $T_{i}$ random walk stays in the $L m^{2}$ neighborhood of $D_{i}$. But $\psi \mid D_{i}$ is isolated and has empty neighborhood. Thus, during $T_{i}$ the random walk stays on the area where is no $m^{0.9}$ consecutive colors. In this case, the probability of generating a block of length at least $m^{2}$ is, for big $m$, bounded above by $\exp \left(-\frac{a m^{2}}{m^{1.8}}\right)=\exp \left(-a m^{0.2}\right)$, where $a>0$ is a constant that does not depend on $m$ (see, e.g. Lemma 2.1 in [15]).
Denote

$$
p_{m}:=P\left(S\left(m^{2}\right)=m\right)
$$

Then

$$
\begin{equation*}
P\left(E_{S}([1, l])\right)=\left(p_{m}\right)^{l} \tag{2.20}
\end{equation*}
$$

So, for each $i \in \mathcal{I}_{I I}$, it holds

$$
P_{\psi}\left(E_{\text {no-block }}^{c}(i) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)=P_{\psi}\left(E_{\text {no-block }}^{c}(i) \mid E_{S}([1, l])\right) \leq \frac{\exp \left(-a m^{0.2}\right)}{\left(p_{m}\right)^{l}}
$$

Now, by local central limit theorem, $p_{m}$ is of order $\frac{1}{m}$. Thus, when $m$ is big enough

$$
\begin{equation*}
P_{\psi}\left(E_{\text {no-block }}(i) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)>0.75, \quad P_{\psi}\left(E_{\text {no-block }}\left(\mathcal{I}_{I I} \backslash \mathcal{I}\right) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)>(0.75)^{\left|\mathcal{I}_{I I}\right|-|\mathcal{I}|} \tag{2.21}
\end{equation*}
$$

The second inequality follows because conditional on $E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)$, everything that happens during the time-interval $T_{i}$, is independent of events happening during the time-interval $T_{j}, j, i \in \mathcal{I}_{I I}$. Hence, for $j, i \in \mathcal{I}_{I I}$ the events $E_{\text {no-block }}(i)$ and $E_{\text {no-block }}(j)$ are conditionally independent.

Suppose now $i \in \mathcal{I} \subset \mathcal{I}_{I I}$. Then $\psi \mid D_{i}$ is weak-OK. By 2) of Theorem 1.2 we now get that

$$
P_{\psi}\left(E_{\sqsubseteq}^{c}(i) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)=P_{\psi}\left(E_{\sqsubseteq}^{c}(i) \mid S_{d_{i-1}}\left(m^{2}\right)=d_{i}\right) \leq 0.25
$$

This also means that, with $i \in \mathcal{I}$

$$
\begin{aligned}
& P_{\psi}\left(\left(E_{\text {no-block }}(i) \cap E_{\sqsubseteq}(i)\right)^{c} \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right) \leq \\
& P_{\psi}\left(E_{\text {no-block }}^{c}(i) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)+P_{\psi}\left(E_{\sqsubseteq}^{c}(i) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)<0.5 .
\end{aligned}
$$

And, by independence, again

$$
\begin{equation*}
P_{\psi}\left(E_{n o-b l o c k}(\mathcal{I}) \cap E_{\sqsubseteq}(\mathcal{I}) \mid E_{\min } \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)>(0.5)^{|\mathcal{I}|} \tag{2.22}
\end{equation*}
$$

Finally, by the same independence-argument, (2.22) and (2.21),

$$
\begin{align*}
& P_{\psi}\left(E_{\text {no-block }}\left(\mathcal{I}_{I I}\right) \cap E_{\sqsubseteq}(\mathcal{I}) \mid E_{\text {min }} \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)= \\
& P_{\psi}\left(\left(E_{\text {no-block }}(\mathcal{I}) \cap E_{\sqsubseteq}(\mathcal{I})\right) \cap E_{\text {no-block }}\left(\mathcal{I}_{I I} \backslash \mathcal{I}\right) \mid E_{\text {min }} \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)>(0.5)^{l} \tag{2.23}
\end{align*}
$$

Consider $\left[r_{j}, s_{j}\right], j=1, \ldots, k$ as in (2.12). By the definition of $\mathcal{I}_{I I},\left[s_{i}-L m^{2}, r_{i+1}+L m^{2}\right]$ is empty, for each $i=1, \ldots, k-1$ as well as for $\left[s_{k}-L m^{2}, l\right]$, if $s_{k}<l$. This implies that these intervals do not contain any small block (and, therefore, no big clusters). Also [ $\left.s_{i}-L m^{2}-m, s_{i}-L m^{2}\right]$ as well as $\left[r_{i+1}+L m^{2}, r_{i+1}+L m^{2}+m\right](i=1, \ldots, k-1)$ and $\left[s_{k}-L m^{2}-m, s_{i}-L m^{2}-m\right]$, if $s_{k}<l$, contain a fence. This means that a interval $\left[s_{i}-L m^{2}, r_{i+1}+L m^{2}\right](i=1, \ldots, k-1)$ as well as $\left[s_{k}-L m^{2}-m, s_{i}-L m^{2}-m\right]$ (if $s_{k}<l$ ) is not inside a big cluster (without fences this could be a case even if the interval is empty). The emptiness and the isolation of $\left[s_{i}, r_{j}\right]$ imply that the cluster-counting vector $q\left(\psi^{m^{2} l}\right)$ is constant on $\mathcal{I}_{I I}$.
The event $E_{n o-b l o c k}\left(\mathcal{I}_{I I}\right) \cap E_{\text {min }}$ ensures that the word $\chi \mid\left[s_{i}-m \bar{v}, r_{i+1}+m \bar{v}\right], i=1, \ldots, k-1$ does not contain more than $m \bar{v}+m^{2}$ consecutive colors. The same is true for the word $\chi \mid\left[s_{k}-m \bar{v}, l\right]$. The event $E_{\min }$ also guarantees that all big blocks in observations end before time interval $T_{i}, i \in \mathcal{I}_{I I}$. Hence, the block-counting vector $\hat{q}\left(\chi^{m^{2} l}\right)$ is constant on $\mathcal{I}_{I I}$. Thus, $\hat{q}\left(\chi^{t_{l}}\right) \leq q\left(\psi^{t_{l}}\right)$ if $\hat{q}_{i}\left(\chi^{t_{l}}\right) \leq q_{i}\left(\psi^{t_{l}}\right)$ for each $i \in \mathcal{I}_{I I}^{c}$. The latter holds if and only if $\hat{q}\left(\chi \mid\left[r_{j} m, s_{j} m\right]\right) \leq q\left(\psi \mid\left[r_{j}, s_{j}\right]\right)$ for each $j=1, \ldots, k$. Hence

$$
E_{\min } \cap E_{\text {no-block }}\left(\mathcal{I}_{I I}\right) \subset\left\{\hat{q}\left(\chi^{m^{2} l}\right) \leq q\left(\xi^{m l}\right)\right\}
$$

This means
$E_{\text {min }} \cap E_{\text {no-block }}\left(\mathcal{I}_{I I}\right) \cap E_{\sqsubseteq}(\mathcal{I}) \cap E_{S}\left(\mathcal{I}_{I I}\right) \subset\left\{\hat{g}\left(\chi^{l m^{2}}\right) \sqsubseteq \mathcal{I} g\left(\xi^{m l}\right), \quad \hat{q}\left(\chi^{m^{2} l}\right) \leq q\left(\xi^{m l}\right), \quad S\left(m^{2} l\right)=m l\right\}$.
From (2.23), (2.20) and (2.18) it follows

$$
\begin{align*}
& P_{\psi}\left(E_{\text {min }} \cap E_{\text {no-block }}\left(\mathcal{I}_{I I}\right) \cap E_{\sqsubseteq}(\mathcal{I}) \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)= \\
& P_{\psi}\left(E_{\sqsubseteq}(\mathcal{I}) \cap E_{\text {no-block }}\left(\mathcal{I}_{I I}\right) \cap \mid E_{\text {min }} \cap E_{S}\left(\mathcal{I}_{I I}\right)\right) P_{\psi}\left(E_{\text {min }} \cap E_{S}\left(\mathcal{I}_{I I}\right)\right)> \\
& (0.5)^{l} P_{\psi}\left(E_{\text {min }} \cap E_{S}\left(\mathcal{I}_{I I}\right)\right) \geq(0.5)^{l}\left(p_{\text {min }}\right)^{\left|\mathcal{I}_{I I}^{c}\right| m^{2}} P_{\psi}\left(E_{S}\left(\mathcal{I}_{I I}\right)\right) \geq(0.5)^{l}\left(p_{\text {min }}\right)^{\left|\mathcal{I}_{I I}^{c}\right| m^{2}}\left(p_{m}\right)^{l} \tag{2.25}
\end{align*}
$$

Hence (2.24), (2.25) and the inequality $\left|\mathcal{I}_{I I}^{c}\right| \leq l \exp \left(-m^{0.7}\right)$ imply

$$
\begin{aligned}
& P_{\psi}\left(\hat{g}\left(\chi^{l m^{2}}\right) \sqsubseteq_{\mathcal{I}} g\left(\xi^{m l}\right), \quad \hat{q}\left(\chi^{m^{2} l}\right) \leq q\left(\xi^{m l}\right), \quad S\left(m^{2} l\right)=m l\right) \geq(0.5)^{l}\left(p_{m i n}\right)^{\left|\mathcal{I}_{I I}^{c}\right| m^{2}}\left(p_{s}\right)^{l} \geq \\
& {\left[0.5 p_{m}\left(p_{\text {min }}\right)^{m^{2} \exp \left(-m^{0.7}\right)}\right]^{l}=\exp \left[l\left(\ln \left(0.5 p_{m}\right)+m^{2} \exp \left(-m^{0.7}\right) \ln \left(p_{\text {min }}\right)\right)\right]=\exp \left[-l \alpha_{I I}(m)\right]}
\end{aligned}
$$

Let us show that, for $n$ big enough,

$$
\begin{equation*}
8 \alpha_{I I}(n)=-8 \ln \left(0.5 p_{s}\right)-m(n)^{2} \exp \left(-m(n)^{0.7}\right) \ln \left(p_{\min }\right)<n^{2} \frac{\ln 2}{2}=\alpha_{I}(n) \tag{2.26}
\end{equation*}
$$

By the LCLT, $p_{m}$ is of order $\frac{1}{m}$, meaning that $-\ln \left(0.5 p_{m}\right)$ is of order $\ln 2 m$. On the other hand, $m(n)<\exp (2 n)([15],(3.10))$, implying that $-\ln \left(0.5 p_{m}\right)$ is of order $n$. The expression

$$
-m(n)^{2} \exp \left(-m(n)^{0.7}\right) \ln \left(p_{\min }\right)
$$

is negligible in comparison with $-\ln \left(0.5 p_{s}\right)$. So, if $n$ is big enough, it holds $\alpha_{I I}(n)<K n$, for some $K<\infty$. Since $\alpha_{I I}(n)$ is of order $n^{2}$, for big $n$, the inequality (2.26) clearly holds.

### 2.4 Block at origin

Define the event

$$
E_{\text {origin }}:=\left\{\xi(-L)=\cdots=\xi(-1) \neq \xi(0)=\cdots=\xi\left(m^{3} L\right) \neq \xi\left(m^{3} L+1\right)=\cdots=\xi\left(m^{3} L+L\right)\right\}
$$

The reason of block-counting is the following observation. Recall the definition of $E_{z, I}$ given in (1.14). The next theorem formalizes the argument explained in Subsection 1.4.5.

Theorem 2.4 If $z<0$ then

$$
\begin{equation*}
E_{\text {origin }} \cap\left\{\hat{q}\left(\chi_{z}^{t_{l}}\right) \leq q\left(\xi^{m l}\right)\right\} \subset E_{z, \mathcal{I}\left(\xi^{m l}\right)} \tag{2.27}
\end{equation*}
$$

Proof. Let $\xi=\psi, \mathcal{I}=\mathcal{I}(\psi)$. Let $i \in \mathcal{I}$. The interval $D_{i}$ is isolated and, hence, $D_{i}$ is not included into any big cluster of $\psi$, i.e. $q_{i}\left(\psi^{d_{l}}\right)=q_{i}\left(\psi^{d_{i}}\right)$. The interval $D_{i}$ has empty neighborhood, which together with the isolation implies that the number of big clusters in $\left[0, d_{i}\right]$ is the same as the number of big clusters in $\left[0, d_{i}-L m^{2}-m\right)=\left[0, d_{i-1-L m}\right)$ or

$$
\begin{equation*}
q_{i}\left(\psi^{m^{2} l}\right)=q_{i-1-L m}\left(\psi^{m^{2} l}\right) \tag{2.28}
\end{equation*}
$$

Let $z<0$. By crossing an interval, the random walk cannot produce less big blocks than the number of big clusters in this interval. Hence, the number of big blocks in observations generated by $S_{z}$ by crossing the interval $\left[z, d_{i-1-L m}\right]$ is at least the number of big clusters in $\left[z, d_{i-1-L m}\right]$. Suppose now that $E_{\text {origin }}$ holds. Then the interval $\left[z, d_{i-1-L m}\right]$ contains strictly more big clusters than the interval $\left[0, d_{i-1-L m}\right]$. Therefore, the number of big blocks in observations generated by $S_{z}$ by crossing the interval $\left[z, d_{i-1-L m}\right]$ is strictly bigger than the number of big clusters in $\psi \mid\left[0, d_{i-1-L m}\right]$. By (2.28), this number equals $q_{i}\left(\psi^{m^{2} l}\right)$. Hence, if $S_{z}\left(t_{i}\right) \geq d_{i-1-L m}$, then $\hat{q}_{i}\left(\chi_{z}^{m^{2} l}\right)>q_{i}\left(\psi^{m^{2} l}\right)$. Consequently, $E_{\text {origin }} \cap E_{z, \mathcal{I}}^{c} \subset E_{\text {origin }} \cap\left\{\hat{q}\left(\chi_{z}^{t_{l}}\right) \leq q\left(\xi^{m l}\right)\right\}^{c}$. This proves the statement.

Define

$$
E_{\text {mistake }}(z):=\left\{\hat{q}\left(\chi_{z}^{m^{2} l}\right) \leq q\left(\xi^{m l}\right)\right\} \cap\left\{\hat{g}\left(\chi_{z}^{m^{2} l}\right) \sqsubseteq_{\mathcal{I}\left(\xi^{m l}\right)} g\left(\xi^{m l}\right)\right\} \cap E_{\text {origin }}
$$

Corollary 2.1 If $z<0$, then for $n$ and $l$ big enough

$$
\begin{equation*}
P\left(E_{\text {mistake }}(z) \cap E_{O K}\right) \leq \exp \left(-\alpha_{I} l\right) \tag{2.29}
\end{equation*}
$$

Proof. By (2.27) we have

$$
E_{\text {mistake }}(z) \subset E_{z, \mathcal{I}\left(\xi^{m l}\right)} \cap\left\{\hat{g}\left(\chi_{z}^{m^{2} l}\right) \sqsubseteq_{\mathcal{I}\left(\xi^{m l}\right)} g\left(\xi^{m l}\right)\right\} .
$$

Thus

$$
\begin{equation*}
E_{\text {mistake }}(z) \cap E_{O K} \subset E_{z, \mathcal{I}\left(\xi^{m l}\right)} \cap\left\{\hat{g}\left(\chi_{z}^{m^{2} l}\right) \sqsubseteq_{\mathcal{I}\left(\xi^{m l}\right)} g\left(\xi^{m l}\right)\right\} \cap E_{O K} \tag{2.30}
\end{equation*}
$$

Consider the right side of (2.30). By $E_{O K}^{*}$ and (2.3), $\left|\mathcal{I}\left(\xi^{m l}\right)\right| \geq l(1-3 \epsilon(n)) \mid$. Thus, if the right side of (2.30) holds, then there exists a subset $I \subset \mathcal{I}\left(\xi^{m l}\right)$ such that $|I|=|l(1-3 \epsilon(n))|,\left\{\hat{g}\left(\chi_{z}^{m^{2} l}\right) \sqsubseteq_{I} g\left(\xi^{m l}\right)\right\}$ and $E_{z, I}$ holds. By Theorem 2.2, this event has probability not bigger than $\exp \left(-l \alpha_{I}\right)$.

## 3 Reconstruction at level $l_{1}$

### 3.1 Some definitions

* A vector $I \in \mathbb{Z}^{[0, n]}$ is ladder interval of length $n$, if $I=(a, a+L, a+2 L, \ldots, a+n L)$ for some $a \in \mathbb{Z}$. Let $\mathcal{L}(n)$ be the set of all ladder intervals of length $n$.

Let $I$ be a ladder interval and $\phi \in\{0,1\}^{I}$. The mapping $\phi$ is called a ladder piece. If $\phi \in\{0,1\}^{D}$, $I \subset D$ is a ladder interval, we sometimes say that $\phi \mid I$ is a ladder piece of $\phi($ or $\phi \mid D)$.

Hence, a ladder piece of non-random scenery $\psi$ is any vector $(\psi(a), \psi(a+L), \ldots, \psi(a+n L)), a \in \mathbb{Z}, n \in \mathbb{N}$.
Recall: two pieces of scenery $\phi: \in\{0,1\}^{D}$ and $\phi^{\prime} \in\{0,1\}^{D^{\prime}}$ are equivalent, $\phi \approx \phi^{\prime}$, if there is a mapping $T: \mathbb{Z} \mapsto \mathbb{Z}, T(z)=a z+b, a \in\{+1,-1\}, b \in \mathbb{Z}$ such that $T: D \mapsto D^{\prime}$ is a bijection and $\phi(z)=\phi^{\prime}(T(z))$ $\forall z \in D$. We also write $T \phi=\phi^{\prime}$.
Every word $w \in\{0,1\}^{n+1}$ is a mapping $w \in\{0,1\}^{[0, n]}$. Hence, if $I=[a, a+n]$ for some $n$, and $\phi \in\{0,1\}^{I}$, then the equivalence $\phi \approx w$ is well defined. By definition, it means that $\phi(a)=w(1), \ldots, \phi(a+n)=$ $w(n+1)$ or $\phi(a)=w(n+1), \ldots, \phi(a+n)=w(1)$.
Let $I=(a, a+L, \ldots, a+n L)$ be a ladder-interval and let $\phi \in\{0,1\}^{I}$ be a ladder piece. We write $\phi \approx_{l} w$, if $\phi(a)=w(1), \ldots, \phi(a+L n)=w(n+1)$ or $\phi(a)=w(n+1), \ldots, \phi(a+L n)=w(1)$. Hence, if $L=1$, then the relation $" \approx_{l} "$ is the same as the equivalence $" \approx "$.

Given a ladder piece $\phi \in\{0,1\}^{I}, I \in \mathcal{L}(n)$, we say that $w \in\{0,1\}^{n+1}$ is a ladder word of $\phi$, if $\phi \approx_{l} w$. Hence, any ladder piece has at most two ladder words that are equivalent. Also note that two ladder pieces are equivalent, if and only if their ladder words coincide. (In the notation of $[\mathrm{L}, \mathrm{M}, \mathrm{M}], w$ is a ladder word of $\phi$, if $w \in\left\{(\phi)_{\rightarrow},(\phi)_{\leftarrow}\right\}$.)

* In this chapter, $l_{1}, c_{1}$ stand for positive integers, they will be specified later. We denote

$$
I_{1}:=\left[-\exp \left(3 l_{1}\right), \exp \left(3 l_{1}\right)\right]
$$

* The following event, $B_{\text {unique fit }}^{1}$, states that any ladder piece of $\xi \mid I_{1}$ of length $\frac{l_{1} c_{1}}{4}$ has unique ladder word up to equivalence. Formally,

$$
B_{\text {unique fit }}^{1}:=\left\{\text { if } I, J \in \mathcal{L}\left(l_{1} c_{1} / 4\right), I, J \subset I_{1} \text { and } I \neq J \text { then } \xi|I \not \approx \xi| J\right\}
$$

* Suppose $x, y \in \mathbb{Z}, y=x+\left(l_{1} c_{1}\right) L$. In this case there is only one admissible path of length $c_{1} l_{1}$ from $x$ to $y$, i.e. there exist unique $R \in \mathcal{R}\left(l_{1} c_{1}\right)$ such that $R\left(c_{1} l_{1}\right)-R(0)=\left(l_{1} c_{1}\right) L$. Obviously, this path consists of maximum jumps, only, i.e. $R(i+1)-R(i)=L, i=0,1, \ldots, l_{1} c_{1}-1$.
Suppose now that $x, y \in \mathbb{Z}, x<y$ are such that $y<x+\left(l_{1} c_{1}\right) L$. In this case, this might happen that there is no admissible path going from $x$ to $y$ with exactly $l_{1} c_{1}$ steps. However, if there is one such admissible path, then it is clearly not unique. The following event, $B_{\text {recon straight }}^{1}$, states that if $x, y \in I_{1}$, then among these admissible paths, there are at least two that generate different words in the observations. More precisely,

$$
B_{\text {recon straight }}^{1}:=\left\{\begin{array}{c}
\text { if } R \in \mathcal{R}\left(l_{1} c_{1}\right) \text { such that } R(0), R\left(l_{1} c_{1}\right) \in I_{1} \text { and } R\left(l_{1} c_{1}\right)-R(0)<\left(l_{1} c_{1}\right) L, \text { then } \\
\exists R^{\prime} \in \mathcal{R}\left(l_{1} c_{1}\right) \text { such that } R(0)=R^{\prime}(0), R\left(c_{1} l_{1}\right)=R^{\prime}\left(c_{1} l_{1}\right) \text { and } \xi \circ R \neq \xi \circ R^{\prime}
\end{array}\right\}
$$

* Let $\psi$ be a non-random scenery. We say that $x$ is a left-barrier point of $\psi$, if

$$
\psi(x+L)=\cdots=\psi(x+1) \neq \psi(x)=\cdots=\psi\left(x-m^{3} L\right) \neq \psi\left(x-m^{3} L-1\right)=\cdots=\psi\left(x-m^{3} L-L\right)
$$

We say $y$ is a right-barrier point of $\psi$, if

$$
\psi(y-L)=\cdots=\psi(y-1) \neq \psi(y)=\cdots=\psi\left(y+m^{3} L\right) \neq \psi\left(y+m^{3} L+1\right)=\cdots=\psi\left(y+m^{3} L+L\right)
$$

The pair $(x, y)$ is called a barrier of $\psi$, if $x$ is a left- and $y$ is a right-barrier point. Recall the event $E_{\text {origin }}$. The point $y$ is a right-barrier point of $\psi$, if the translated scenery $\psi_{y}:=(\psi(i+y))_{i \in \mathbb{Z}}$ belongs to the event $E_{\text {origin }}$. Similarly, $x$ is a left-barrier point, if the translated and reflected scenery $\psi_{x}^{-}:=(\psi(x-i))_{i \in \mathbb{Z}}$ belongs to the event $E_{\text {origin }}$.

We consider the barriers of $\xi,(x, y)$ such that $y-x=\left(c_{1} l_{1}\right) L$. In order to carry on the reconstruction in level $l_{1}$, every interval $\left[z, z+\left(c_{1} l_{1} / 4\right) L\right], z \in I_{1}$ should contain enough left-barrier points of such barriers. This is the meaning of the event $B_{\text {enough barriers }}^{1}$. More precisely,

$$
B_{\text {enough barriers }}^{1}:=\left\{\begin{array}{c}
\text { for any } j=0, \ldots, L-1 \text { and for any } z \in I_{1} \\
\text { there exists } x \in\left[z, z+\left(c_{1} l_{1} / 4\right) L\right] \text { such that: } \\
x \bmod L=j \text { and }\left(x, x+\left(c_{1} l_{1}\right) L\right) \text { is a barrier of } \xi
\end{array}\right\}
$$

* We now define the left-side counterparts of $g, \hat{g}, q$ and $\hat{q}$. For a word $u=\left(u_{1}, \ldots, u_{n}\right)$ denote by $u^{-}$its reflection, i.e. $u^{-}:=\left(u_{n}, \ldots, u_{1}\right)$. Now let

$$
q^{*}:\{0,1\}^{l m+1} \mapsto \mathbb{N}^{l}, \quad \hat{q}^{*}:\{0,1\}^{l m^{2}+1} \mapsto \mathbb{N}^{l}, \quad g^{*}:\{0,1\}^{l m+1} \mapsto\{0,1\}^{l n^{2}+1}, \quad \hat{g}^{*}:\{0,1\}^{l m^{2}+1} \mapsto\{0,1\}^{l n^{2}}
$$

be as follows

$$
\begin{align*}
q^{*}(w)=q\left(w^{-}\right), \quad g^{*}(w)=g\left(w^{-}\right), \quad w \in\{0,1\}^{l m+1}  \tag{3.1}\\
\hat{q}^{*}(v)=\hat{q}\left(v^{-}\right), \quad \hat{g}^{*}(v)=\hat{g}\left(v^{-}\right), \quad v \in\{0,1\}^{l m^{2}+1} \tag{3.2}
\end{align*}
$$

* Finally, we put

$$
l=l_{1} \cdot l_{2}
$$

Hence, the requirement " $l$ big enough" in all statement of previous chapter is equivalent to the requirement " $l_{1}$ big enough".

### 3.2 Stopping-time events

### 3.2.1 Right side

${ }^{*}$ Let $\tau(1), \tau(2), \ldots$ be a sequence of $\mathcal{F}$-adapted stopping times satisfying

$$
\begin{equation*}
\tau(k)-\tau(k-1) \geq 2 \exp \left(2 l_{1}\right), \quad k=2,3, \ldots \tag{3.3}
\end{equation*}
$$

Let $z \in \mathbb{Z}$. Define

$$
\kappa^{3}(z, 1):=\min \left\{j: S\left(\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}+c_{1} l_{1}\right)=z\right\}
$$

and, inductively,

$$
\kappa^{3}(z, k):=\min \left\{j>\kappa^{3}(z, k-1): S\left(\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}+c_{1} l_{1}\right)=z\right\}
$$

Thus, $\kappa^{3}(z, k)$ is the index of $k$-th stopping time $\tau(j)$, for which $S\left(\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}+c_{1} l_{1}\right)=z$. Hence, at time

$$
T_{z}^{3}(k):=\tau\left(\kappa^{3}(z, k)\right)+\exp \left(2 l_{1}\right)+l m^{2}+c_{1} l_{1}
$$

the random walk is at position $z$. Let $w_{z}^{3}(k)$ denote the observation-word of length $l m^{2}$ starting at time $T_{z}^{3}(k)$, i.e.

$$
w_{z}^{3}(k):=\chi \mid\left[T_{z}^{3}(k), T_{z}^{3}(k)+l m^{2}\right] .
$$

[^0]in the values of $\hat{q}\left(w_{k}^{3}\right)$ and $\hat{g}\left(w_{k}^{3}\right)$ (note that the length of $w_{k}^{3}$ is $\left.l m^{2}\right)$. Having sufficiently many $k$ s , the values $\hat{q}\left(w_{k}^{3}\right)$ and $\hat{g}\left(w_{k}^{3}\right)$ give us some information about $g(\xi \mid[z, z+m l])$ and $q(\xi \mid[z, z+m l])$. Indeed, by Theorem 2.3, there is a proportion of words $w_{k}^{3}$ such that $\hat{q}\left(w_{k}^{3}\right) \leq q(\xi \mid[z, z+m l])$ and $\hat{g}\left(w_{k}^{3}\right) \sqsubseteq_{\mathcal{I}(\xi \mid[z, z+m l])} g(\xi \mid[z, z+m l])$. On the other hand, if $y \in \mathbb{Z}$ is a right-barrier point bigger than $z$, then, by Corollary 2.1, the probability for the relations above to hold is rather small. Hence we expect that for such $y$ the relations
\[

$$
\begin{equation*}
\hat{q}\left(w_{k}^{3}\right) \leq q(\xi \mid[y, y+m l]), \quad \hat{g}\left(w_{k}^{3}\right) \sqsubseteq_{\mathcal{I}(\xi \mid[y, y+m l])} g(\xi \mid[y, y+m l]) \tag{3.4}
\end{equation*}
$$

\]

do not occur.
To make these ideas precise, for each $y \in \mathbb{Z}$ we define

$$
g_{y}(\xi):=g(\xi \mid[y, y+m l]), \quad q_{y}(\xi):=q(\xi[y, y+m l])
$$

The following event is a counterpart of $E_{\text {mistake }}(z)$. It states that although $y$ is a right-barrier point and $z<y$, the mistake (3.4) still holds.
$E_{\text {mistake-r }}^{1}(z, y, k):=\left\{\hat{q}\left(w_{z}^{3}(k)\right) \leq q_{y}(\xi)\right\} \cap\left\{\hat{g}\left(w_{z}^{3}(k)\right) \sqsubseteq_{\mathcal{I}(\xi \mid[y, y+m l])} g_{y}(\xi)\right\} \cap\{y$ is a right-barrier point $\}$.
Finally, let

$$
E_{\text {mistake-r }}^{1}:=\bigcup E_{\text {mistake-r }}^{1}(z, y, k)
$$

where the union is taken over all $z, y, k$ such that $z<y, z, y \in I_{1}$ and $k \leq \exp \left(\alpha l_{1}\right)$.

### 3.2.2 Left side

We now introduce the left-side counterparts of defined notions. At first, let

$$
\kappa^{1}(z, 1):=\min \left\{j: S\left(\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}\right)=z\right\}
$$

and, inductively

$$
\kappa^{1}(z, k):=\min \left\{j>\kappa^{1}(z, k-1): S\left(\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}\right)=z\right\}
$$

Thus, $\kappa^{l}(z, k)$ is the index of $k$-th stopping time $\tau(j)$, for which $S\left(\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}\right)=z$. Hence, at time

$$
T_{z}^{1}(k):=\tau\left(\kappa^{1}(z, k)\right)+\exp \left(2 l_{1}\right)+l m^{2}
$$

the random walk is at position $z$. Let $w_{z}^{1}(k)$ denote the observation-word of length $l m^{2}$ ending at time $T_{z}^{1}(k)$, i.e.

$$
w_{z}^{1}(k):=\chi \mid\left[T_{z}^{1}(k)-l m^{2}, T_{z}^{1}(k)\right]
$$

As previously, we consider the characteristics $\hat{q}^{*}\left(w_{z}^{1}(k)\right), \hat{g}^{*}\left(w_{z}^{1}(k)\right)$ and we compare them with the corresponding functions $q^{*}(\xi \mid[x-m l, x])$ and $g^{*}(\xi \mid[x-m l, x])$, where $x<z$ is a left-barrier point. For this define

$$
\left.q_{x}^{*}(\xi):=q^{*}(\xi \mid[x-m l, x]), \quad g_{x}^{*}(\xi):=g^{*}(\xi \mid[x-m l, x]), \quad \mathcal{I}^{*}(\xi \mid[x-m l, x]):=\mathcal{I}(\xi \mid[x-m l, x])^{-}\right)
$$

The counterpart of $E_{\text {mistake-r }}^{1}$ is as follows
$E_{\text {mistake-1 }}^{1}(z, x, k):=\left\{\hat{q}^{*}\left(w_{z}^{1}(k)\right) \leq q_{x}^{*}(\xi)\right\} \cap\left\{\hat{g}^{*}\left(w_{z}^{1}(k)\right) \sqsubseteq_{\mathcal{I}^{*}(\xi \mid[x-m l, x])} g_{x}^{*}(\xi)\right\} \cap\{x$ is a left-barrier point $\}$.
Finally, let

$$
E_{\text {mistake-1 }}^{1}:=\bigcup E_{\text {mistake-l }}^{1}(z, x, k)
$$

where the union is taken over all $z, x, k$ such that $x<z, z, x \in I_{1}$ and $k \leq \exp \left(\alpha l_{1}\right)$.
Finally, let

$$
E_{\text {no mistake }}^{1}:=\left(E_{\text {mistake-1 }}^{1} \cap E_{\text {mistake-r }}^{1}\right)^{c}
$$

### 3.2.3 Attributes

* Define

$$
T^{1}(j):=\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}, \quad T^{3}(j):=\tau(j)+\exp \left(2 l_{1}\right)+l m^{2}+c_{1} l_{1}=T^{1}(j)+c_{1} l_{1}
$$

Hence $T^{1}(j)$ (or $T^{3}(j)$ ) is defined as $T_{z}^{1}(k)$ (or $T_{z}^{1}(k)$ ) by dropping the requirement that the random walk is at position $z$. We now define the counterparts of $w_{z}^{1}(k)$ and $w_{z}^{3}(k)$.

Let, for each $j=1,2, \ldots$

$$
\begin{aligned}
w^{1}(j) & :=\chi \mid\left[T^{1}(j)-l m^{2}, T^{1}(j)\right] \\
w^{2}(j) & :=\chi \mid\left[T^{1}(j), T^{3}(j)\right] \\
w^{3}(j) & :=\chi \mid\left[T^{3}(j), T^{3}(j)+l m^{2}\right]
\end{aligned}
$$

Let $x, y \in I$ be such that $y-x=L c_{1} l_{1}$. We consider stopping times $\tau(1), \tau(2), \ldots, \tau\left(\exp \left(\alpha l_{1}\right)\right)$. The following event states that among these stopping times there is at least $\exp \left(\gamma l_{1}\right)$ stopping times, $\tau(j)$ such that: $S\left(T^{1}(j)\right)=x, S\left(T^{3}(j)\right)=y$ and

$$
\begin{align*}
& \hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \quad \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{\mathcal{I}^{*}} g_{x}^{*}(\xi)  \tag{3.5}\\
& \hat{q}\left(w^{3}(j)\right) \leq q_{y}(\xi), \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}} g_{y}(\xi), \tag{3.6}
\end{align*}
$$

where

$$
\mathcal{I}^{*}:=\mathcal{I}^{*}(\xi \mid[x-l m, x]), \quad \mathcal{I}:=\mathcal{I}(\xi \mid[y, y+l m])
$$

More precisely,

$$
E_{\text {enough times }}^{1}(x, y):=\left\{\begin{array}{c}
\text { there exists a set } \quad J(x, y) \subset\left[1, \exp \left(\alpha l_{1}\right)\right] \quad \text { such that } \\
|J(x, y)| \geq \exp \left(\gamma l_{1}\right) \text { and for every } j \in|J(x, y)| \\
S\left(T^{1}(j)\right)=x, S\left(T^{3}(j)\right)=y, \\
\hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{\mathcal{I}^{*}} g_{x}^{*}(\xi), \\
\hat{q}\left(w^{3}(j)\right) \leq q_{y}(\xi), \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{\mathcal{I}} g_{y}(\xi)
\end{array}\right\} .
$$

Finally, let

$$
E_{\text {enough times }}^{1}:=\bigcap_{x, y \in I_{1}, x-y=L c_{1} l_{1}} E_{\text {enough times }}^{1}(x, y) .
$$

* Note, for every $u \in\{0,1\}^{l m+1}, q(u)=\left(q_{1}, \ldots, q_{l}\right)$ is vector, such that $q_{i}=\{0,1, \ldots, l\}, q_{1}=0$ and $q_{i} \leq q_{i+1} \leq q_{i}+1$. Any such vector is called a $q$-vector. Hence, for every $u, q(u)$ and $q^{*}(u)$ are $q$-vectors.

Recall that, for any $u \in\{0,1\}^{l m+1}, g(u)=\left(g_{1}, \ldots, q_{l}\right)$, where $g_{i} \in\{0,1\}^{n^{2}+1}$. Any such word is called a $g$-word. Hence, for each $u, g(u)$ and $g^{*}(u)$ are $g$-words.

We say that a word $u \in\{0,1\}^{l m+1}$ is $\mathrm{OK}^{*}$, if $u^{-}$is OK. Clearly, if $\xi \mid[x-l m, x]$ is $\mathrm{OK}^{*}$, then $\left|\mathcal{I}^{*}\right| \geq$ $l(1-3 \epsilon(n))$. Similarly, if $\xi \mid[y, y+l m]$ is OK, then $|\mathcal{I}| \geq l(1-3 \epsilon(n))$.

Note, if (3.5) and (3.6) hold and $\xi \mid[x-l m, x]$ together with $\xi \mid[y, y+l m]$ are OK* and OK, respectively, then there exists subsets $I^{*}, I \subset\{1, \ldots, l\}$, such that $\left|I^{*}\right|,|I| \geq l(1-3 \epsilon(n))$, the $q$-vectors $q$, $q^{*}$ and $g$-words $g, g^{*}$ such that

$$
\begin{equation*}
\hat{q}^{*}\left(w^{1}(j)\right) \leq q^{*}, \quad \hat{g}^{*}\left(w^{1}(j)\right) \sqsubseteq_{I^{*}} g^{*}, \quad \hat{q}\left(w^{3}(j)\right) \leq q, \quad \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{I} g . \tag{3.7}
\end{equation*}
$$

We call $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ a set of attributes, if $I^{*}, I \subset\{1, \ldots, l\},\left|I^{*}\right|,|I| \geq l(1-3 \epsilon(n)), q, q^{*}$ are $q$-vectors and $g^{*}, g$ are $g$-words.

For every set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ we define the index-set $J\left(I^{*}, I, q^{*}, q, g^{*}, g\right) \subset\left[1, \exp \left(\alpha l_{1}\right)\right]$ as follows: $j \in J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ if and only if $j$ satisfies (3.7).

Hence, if $E_{\text {enough times }}^{1}(x, y)$ holds and $\xi \mid[x-l m, x]$ together with $\xi \mid[y, y+l m]$ are OK, then there exists a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that

$$
\begin{equation*}
\left|J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)\right| \geq \exp \left(\gamma l_{1}\right) \tag{3.8}
\end{equation*}
$$

Then also $S\left(T^{1}(j)\right)=x$ and $S\left(T^{3}(j)\right)=y=x+L c_{1} l_{1}$, i.e. $S\left(T^{3}(j)\right)-S\left(T^{1}(j)\right)=L c_{1} l_{1}$.

* The following event implies: if $\exists J^{\prime} \subset J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that $\left|J^{\prime}\right| \geq \exp \left(\gamma l_{1}\right)$ and $\forall j \in J^{\prime}$ it holds $S\left(T^{3}(j)\right)-S\left(T^{1}(j)\right)<L c_{1} l_{1}$, then there exists at least two indexes $j^{\prime}, j^{\prime \prime} \in J^{\prime}$ such that $w^{2}\left(j^{\prime}\right) \neq w^{2}\left(j^{\prime \prime}\right)$.

Formally, we fix a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$, we consider the stopping times $\tau(0), \tau(1), \ldots$ and we define the indexes

$$
\begin{align*}
& j_{1}:=\min \{j \geq 0:  \tag{3.9}\\
& \left.\quad \hat{q}^{*}\left(w^{1}(j)\right) \leq q^{*}, \hat{g}\left(w^{1}(j)\right) \sqsubseteq_{I^{*}} g^{*}, \hat{q}\left(w^{3}(j)\right) \leq q, \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{I} g,\left|S\left(T^{3}(j)\right)-S\left(T^{1}(j)\right)\right|<L c_{1} l_{1}\right\}  \tag{3.10}\\
& j_{k}:=\min \left\{j>j_{k-1}:\right.  \tag{3.11}\\
& \left.\quad \hat{q}^{*}\left(w^{1}(j)\right) \leq q^{*}, \hat{g}\left(w^{1}(j)\right) \sqsubseteq_{I^{*}} g^{*}, \hat{q}\left(w^{3}(j)\right) \leq q, \hat{g}\left(w^{3}(j)\right) \sqsubseteq_{I} g,\left|S\left(T^{3}(j)\right)-S\left(T^{1}(j)\right)\right|<L c_{1} l_{1}\right\} . \tag{3.12}
\end{align*}
$$

Here the minimum over empty set is defined to be $\infty$. Let $\kappa:=\max \left\{k: j_{k}<\infty\right\}$.
Clearly the subindexes $j_{1}, j_{2}, \ldots$ depend on chosen attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$.
Recall $B_{\text {recon straight }}^{1}$. The following events are of similar nature. Let

$$
\begin{aligned}
E_{\text {recon straight }}^{1}\left(I^{*}, I, q^{*}, q, g^{*}, g\right) & :=\left\{\kappa \geq \exp \left(\gamma l_{1}\right), \quad \exists k \leq \exp \left(\gamma l_{1}\right) \text { such that } w^{2}\left(j_{1}\right) \neq w^{2}\left(j_{k}\right)\right\} \cup\left\{\kappa \leq \exp \left(\gamma l_{1}\right)\right\} \\
E_{\text {recon straight }}^{1} & :=\bigcap_{I^{*}, I, q^{*}, q, g^{*}, g} E_{\text {recon straight }}^{1}\left(I^{*}, I, q^{*}, q, g^{*}, g\right)
\end{aligned}
$$

where the intersection is taken over all sets of attributes.

### 3.3 Algorithm

We are ready to give the precise definition of the algorithm $\mathcal{A}^{1}$. The input of $\mathcal{A}^{1}$ consists of three ingredients

- $\exp \left(12 \alpha l_{1}\right)+1$ observations, $\chi \mid\left[0, \exp \left(12 \alpha l_{1}\right)\right]$;
- $\mathcal{F}$-adapted stopping times $\tau=\left(\tau(1), \ldots, \tau\left(\exp \left(\alpha l_{1}\right)\right) \subset\left[0, \exp \left(12 \alpha l_{1}\right)\right]\right.$ satisfying (3.3);
- a piece of original scenery $\psi^{o}=\xi \mid I^{o}$, where $\left|I^{o}\right| \geq 2 c_{1} L l_{1}$ and $I^{o} \subset\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]$.

The output of $\mathcal{A}^{1}$ is a piece of scenery of length $4 \exp \left(l_{1}\right)$. Thus, formally,

$$
\mathcal{A}^{1}:\{0,1\}^{\left[0, \exp \left(12 \alpha l_{1}\right)\right]} \times\left[0, \exp \left(12 \alpha l_{1}\right)\right]^{\left[1, \exp \left(\alpha l_{1}\right)\right]} \times\left(\cup_{k>2 c_{1} l_{1} L}\{0,1\}^{k}\right) \mapsto\{0,1\}^{\left[-2 \exp \left(l_{1}\right), 2 \exp \left(l_{1}\right)\right]}
$$

The aim of $\mathcal{A}^{1}$ is to produce a piece of original scenery that lies between $\xi \mid\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]$ and $\xi \mid\left[-4 \exp \left(l_{1}\right), 4 \exp \left(l_{1}\right)\right]$. Thus, $\mathcal{A}^{1}$ does well, if the following event holds

$$
\begin{equation*}
E_{\text {alg works }}^{1}\left(\tau, I^{o}\right):=\left\{\xi\left|\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right] \sqsubseteq \mathcal{A}^{1}\left(\chi^{\exp \left(12 \alpha l_{1}\right)}, \tau, \xi \mid I^{o}\right) \sqsubseteq \xi\right|\left[-3 \exp \left(l_{1}\right), 3 \exp \left(l_{1}\right)\right]\right\} . \tag{3.13}
\end{equation*}
$$

Obviously the event (3.13) depends on $\tau$ as well as on the chosen interval $I^{o}$. In the following we do not know exactly the interval $I^{o}$. Hence, we want that $\mathcal{A}^{1}$ works with any given interval $I^{o}$. The corresponding event is

$$
E_{\text {alg works }}^{1}(\tau):=\bigcap_{I^{o} \subset\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]} E_{\text {alg works }}^{1}\left(\tau, I^{o}\right)
$$

The construction of $\mathcal{A}^{1}$ consists of two phases.
Phase I Collect the ladder words of $\xi \mid I^{1}$. For this, the observation-words triples $\left(w^{1}(j), w^{2}(j), w^{3}(j)\right)$, defined by $\tau(j)$, are used. The word $w^{2}(j)$ will be collected as a ladder word, if it passes certain selection procedure. We shall specify the selection rule below, this is the core of $\mathcal{A}^{1}$. The set of collected works, i.e. the set of all words, that pass the selection rule, will be denoted by $\mathcal{W}^{1}$.

Phase II We assemble the words from $\mathcal{W}^{1}$ to get a big word of length $4 \exp \left(l_{1}\right)$ as the output. This means the construction of a big word (of length $4 \exp \left(l_{1}\right)$ ) by attaching, one by one, suitable words from $\mathcal{W}^{1}$. We start from $\psi^{o}$, and we attach to it a word from $\mathcal{W}^{1}$, which has an overlap with $\psi^{o}$ at least $\frac{c_{1} l_{1}}{4}$. We then attach a word from $\mathcal{W}^{1}$ to the enlarged $\psi^{o}$ using the same overlapping-criterion. We proceed so, until the desired length has been achieved.

We now give the description and the definition of the selection rule for Phase $\mathbf{I}$ and the precise definition of assembling rule for Phase II. These definitions complete the definition of $\mathcal{A}^{1}$.

The selection rule is the most crucial part of the whole scenery reconstruction. The selection rule must be restrictive enough to ensure that only ladder words of ladder pieces of original scenery $\xi$ can pass it (with high probability). Formally, the following event should hold

$$
E_{\text {only ladders }}^{1}:=\left\{\forall w \in \mathcal{W}^{1} \text { there exists } I \in \mathcal{L}\left(c_{1} l_{1}\right) \text { such that } I \subset I_{1} \text { and } \xi \mid I \approx_{l} w\right\}
$$

On the other hand, the selection rule must be flexible enough to ensure that enough ladder words pass it (otherwise the set $\mathcal{W}^{1}$ is too small). More precisely, the following event should hold

$$
E_{\text {enough ladders }}^{1}:=\left\{\begin{array}{c}
\text { for any } j=0, \ldots, L-1 \text { and for any } z \in I_{1} \\
\text { there exists } x \in\left[z, z+\left(c_{1} l_{1} / 4\right) L\right] \text { such that: } \\
x \bmod L=j \text { and }\left(\xi(x), \xi(x+L), \ldots, \xi\left(x+\left(c_{1} l_{1}\right) L\right)\right) \in \mathcal{W}^{1}
\end{array}\right\}
$$

Let us briefly introduce the main ideas behind the selection rule. The construction of the selection rule used for $\mathcal{A}^{1}$ starts from the fact that, with high probability, the events $E_{\text {enough times }}^{1}$ and

$$
B_{\text {intervals OK }}^{1}:=\left\{\xi \mid[z, z+m l] \text { is OK } \forall z \in I_{1}\right\} \cap\left\{\xi \mid[z, z-m l] \text { is } \mathrm{OK}^{*} \forall z \in I_{1}\right\}
$$

both hold. This means that for each $x, y \in I_{1}, y-x=L c_{1} l_{1}$ there exists a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ and at least $\exp \left(\gamma l_{1}\right)$ stopping times $\tau(j)$ with corresponding index set $J(x, y)$ such that for each $j \in J(x, y)$, (3.7) hold and the word $w^{2}(j)$ is the same, say $w$. This yields the first requirement of selection rule - check whether there exists $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ that satisfies the following condition: $\exists J^{\prime} \subset J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that $\left|J^{\prime}\right| \geq \exp \left(\gamma l_{1}\right)$ and $j \mapsto w^{2}(j)$ is constant on $J^{\prime}$. The existence of such set of attributes and index-set $J^{\prime}$ can be easily checked.

The second requirement of the selection rule is based on the fact that, with high probability the events $B_{\text {enough barriers }}^{1}$ and $E_{\text {no mistake }}^{1}$ hold. This means that if $(x, y)$ form a barrier then the vectors $q_{x}^{*}(\xi), q_{y}(\xi)$ and words $g_{x}(\xi)$ and $g_{y}(\xi)$ cannot be read somewhere else. Hence, if $I^{*}, I, q^{*}, q, g^{*}, g$ found in the first step are indeed $\mathcal{I}^{*}(\xi \mid[x-l m, x]), \mathcal{I}(\xi[y, y+l m]) q_{x}^{*}(\xi), q_{y}(\xi), g_{x}(\xi), g_{y}(\xi)$ as we want them to be, and if $w$ is the word to be selected, then the following must hold: whenever there is a stopping time index $j$ satisfying (3.7), then $w^{2}(j)=w$. Thus, the set $J^{\prime}$ must actually be $J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$. This is the second requirement of the selection rule.
From the argument above, it is clear that if $E_{\text {enough times }}^{1}, B_{\text {intervals OK }}^{1}$, and $E_{\text {no mistake }}^{1}$ hold, then the selection rule will select all ladder words $\left(\xi(x), \xi(x+L), \ldots, \xi\left(x+L\left(c_{1}-1\right) L\right), \xi(y)\right)$, where $(x, y)$ is barrier of $\xi$ and $y-x=L c_{1} l_{1}$. With $B_{\text {enough barriers }}^{1}$ the latter yields $E_{\text {enough ladders }}^{1}$ (see Lemma 3.1 for formal proof). If, in addition $E_{\text {recon straight }}^{1}$ holds, then, as it is not hard to see, the selection procedure will select only ladder words (Lemma 3.1). Hence, the selection rule consisting of two requirements described above is sufficient for our purposes. We now give the formal definition of the selection rule.

Definition 3.1 We define the set $\mathcal{W}^{1}=\mathcal{W}^{1}\left(\chi^{12 \alpha l_{1}}, \tau\right)$ as follows. A word $w \in\{0,1\}^{c_{1} l_{1}+1}$ belongs to $\mathcal{W}^{1}$ if and only if there exists a complect of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that the following conditions are satisfied:

1. $\left|J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)\right| \geq \exp \left(\gamma l_{1}\right)$
2. if $j \in J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$, then $w^{2}(j)=w$.

Let us now formalize the Phase II.
For a ladder interval $I$ and a set $D \subset \mathbb{Z}$ we write $|I \cap D| \geq r$ if there exists a ladder interval $J \in \mathcal{L}(r)$ such that $J \subset D \cap I$. Recall that two pieces of scenery $\phi$ and $\phi^{\prime}$ are strongly equivalent, $\phi \equiv \phi^{\prime}$, if $\phi$ is obtained by some translation of $\phi^{\prime}$. Let $\psi^{o} \in\{0,1\}^{k+1}$ be the given piece of original scenery. Thus, $\psi^{o} \equiv \xi \mid I^{o}$ for some interval $I^{o} \subset\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]$.
Definition 3.2 We say that the piece of scenery $\phi \in\{0,1\}^{\left[-2 \exp \left(l_{1}\right), 2 \exp \left(l_{1}\right)\right]}$ is a solution, formally $\phi \in$ $\mathcal{S}\left(\chi^{12 \alpha l_{1}}, \tau, \psi^{o}\right)$, if and only if there exist $\phi_{i} \in\{0,1\}^{D_{i}}, i=1,2, \ldots, n$ such that $D_{i} \subset\left[-3 \exp \left(l_{1}\right), 3 \exp \left(l_{1}\right)\right]$ and the following conditions are satisfied:

1. $D_{1}=[0, k], \phi_{1} \equiv \psi^{o}$;
2. for each $i=2, \ldots, n$ it holds $\phi_{i} \mid D_{i-1}=\phi_{i-1}$;
3. for each $i=2, \ldots, n$ there exists $I_{i} \in \mathcal{L}\left(c_{1} l_{1}\right)$ such that

3a) $D_{i}=D_{i-1} \cup V_{i}$;
3b) $\left|D_{i} \cap V_{i}\right| \geq \frac{c_{1} l_{1}}{4}$;
3c) $\exists w_{i} \in \mathcal{W}^{1}\left(\chi^{12 \alpha l_{1}}, \tau\right)$ such that $\phi_{i} \mid V_{i} \approx_{l} w_{i}$;
4. $\left[-2 \exp \left(l_{1}\right), 2 \exp \left(l_{1}\right)\right] \subset D_{n}, \quad \phi=\phi_{n} \mid\left[-2 \exp \left(l_{1}\right), 2 \exp \left(l_{1}\right)\right]$.

Finally, the formal definition of $\mathcal{A}^{1}$. The output is any element of $\mathcal{S}$; we choose one of them, if $\mathcal{S}$ is not empty.

Definition 3.3 We define $\mathcal{A}^{1}\left(\chi^{12 \alpha l_{1}}, \tau, \psi^{o}\right)$ as follows:

- If $\mathcal{S}\left(\chi^{12 \alpha l_{1}}, \tau, \psi^{o}\right)$ is nonempty, then we define $\mathcal{A}^{1}\left(\chi^{12 \alpha l_{1}}, \tau, \psi^{o}\right)$ to be its lexicographically smallest element;
- otherwise, $\mathcal{A}^{1}\left(\chi^{12 \alpha l_{1}}, \tau, \psi^{o}\right):=(1)_{\left[-2 \exp \left(l_{1}\right), 2 \exp \left(l_{1}\right)\right]}$.


### 3.4 Combinatorics for main theorem

Recall the stopping times $\tau=\left(\tau(0), \ldots, \tau\left(\exp \left(\alpha l_{1}\right)\right)\right)$. The aim of the stopping times is to stop the random walk $S$ near the origin. It is enough, if $S(\tau(k)) \in\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right]$. Thus, the stopping times do well, if the following event holds

$$
E_{\text {stop }}^{1}(\tau):=\left\{|S(\tau(j))| \leq \exp \left(l_{1}\right), \quad j=0,1,2, \ldots, \exp \left(\alpha l_{1}\right)\right\}
$$

Roughly speaking, the main theorem of the paper states that the algorithm $\mathcal{A}^{1}$ reconstructs correctly with high probability, provided the stopping times $\tau$ indeed stop the random walk close to the origin.

Theorem 3.4 There exists $a\left(l_{2}, n, c_{1}\right)>0$ such that, for $l_{1}$ big enough

$$
\begin{equation*}
P\left(E_{\text {stop }}^{1}(\tau) \cap\left(E_{\text {alg works }}^{1}(\tau)\right)^{c}\right) \leq e^{-a l_{1}} \tag{3.14}
\end{equation*}
$$

The rest of the paper is the proof of Theorem 1.1. At first we prove some inclusions.
Lemma 3.1 The following inclusions hold

$$
\begin{align*}
& E_{\text {recon straight }}^{1} \cap E_{\text {stop }}^{1}(\tau) \subset E_{\text {only ladders }}^{1} ;  \tag{3.15}\\
& B_{\text {intervals } O K}^{1} \cap E_{\text {stop }}^{1}(\tau) \cap E_{\text {no mistake }}^{1} \cap E_{\text {enough times }}^{1} \cap E_{\text {enough barriers }}^{1} \subset E_{\text {enough ladders }}^{1} ;  \tag{3.16}\\
& E_{\text {only ladders }}^{1} \cap E_{\text {enough ladders }}^{1} \cap B_{\text {unique fit }}^{1} \subset E_{\text {alg works }}^{1}(\tau) \tag{3.17}
\end{align*}
$$

provided $l_{1}$ is big enough.
Proof. At first note: if $E_{\text {stop }}^{1}(\tau)$ holds, then, for each $j=1,2, \ldots, \exp \left(\alpha l_{1}\right)$, it holds

$$
\begin{equation*}
\left|S\left(T^{3}(j)\right)\right| \leq\left|S(\tau(j))+L\left(\exp \left(2 l_{1}\right)+l m^{2}+c_{1} l_{1}\right)\right| \leq \exp \left(3 l_{1}\right) \tag{3.18}
\end{equation*}
$$

provided $l_{1}$ is big enough. Thus, in this case, during the time interval $\left[T^{1}(j), T^{3}(j)\right], S$ stays on $I_{1}$, $j=1,2, \ldots, \exp \left(\alpha l_{1}\right)$.

Proof of (3.15):
We prove

$$
\begin{equation*}
\left(E_{\text {only ladders }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau) \subset\left(E_{\text {recon straight }}^{1}\right)^{c} \tag{3.19}
\end{equation*}
$$

Suppose $\left(E_{\text {only ladders }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)$ holds. Then there exists a $w \in \mathcal{W}^{1}$ that is not a ladder word of any ladder piece $\xi \mid I$ of length $l_{1} c_{1}$ such that $I \subset I_{1}$. However, the word $w$ has passed the selection rule. This means that for a complect of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ the conditions 1. and 2. of Definition 3.1 hold. This means, that

$$
\begin{equation*}
\left|S\left(T^{3}(j)\right)-S\left(T^{1}(j)\right)\right|<c_{1} l_{1}, \quad \forall j \in J\left(I^{*}, I, q^{*}, q, g^{*}, g\right) \tag{3.20}
\end{equation*}
$$

Indeed, if there were an index $j^{*} \in J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that (3.20) fails, then there would be a ladder interval $I$ of length $c_{1} l_{1}$ such that $\xi \mid I \approx_{l} w$. Clearly, during the time interval $\left[T^{1}\left(j^{*}\right), T^{3}\left(j^{*}\right)\right]$, the random walk $S$ is on $I$. Since then $S$ is also on $I_{1}$, we get $I \subset I_{1}$. This contradicts our assumption on $w$.
Recall the definition of $\kappa$. Since $\left|J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)\right| \geq \exp \left[\gamma l_{1}\right]$, we have $\kappa \geq \exp \left[\gamma l_{1}\right]$. On the other hand, by 2 of Definition 3.1, for each $j_{k}, k=1,2, \ldots, \exp \left[\gamma l_{1}\right]$, it holds $w\left(j_{k}\right)=w\left(j_{1}\right)=w$. Thus, $E_{\text {recon straight }}^{1}\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ fails. This completes the proof of (3.19).

Proof of (3.16):
Let $x, y \in I_{1}$ and $y-x=c_{1} l_{1}$. Since $B_{\text {intervals OK }}^{1}$ holds then, by (2.3), $I^{*}=\mathcal{I}(\xi \mid[x-l m])$ and $I=\mathcal{I}(\xi \mid[y, y+l m])$ satisfy $\left|I^{*}\right|,|I| \geq l(1-3 \epsilon(n))$. Since $E_{\text {enough times }}^{1}(x, y)$ holds, there exists $q$-vectors $q^{*}=q_{x}^{*}(\xi), q=q_{y}(\xi)$ and $g$-words $g^{*}=g_{x}(\xi), g=g_{y}(\xi)$ such that for each $j \in J(x, y),(3.7)$ holds. Moreover, $|J(x, y)| \geq \exp \left(\gamma l_{1}\right)$ and for each $j \in J(x, y)$ it holds $S\left(T^{1}(j)\right)=x$ and $S\left(T^{3}(j)\right)=y$. Then, obviously, $w^{2}(j)=(\xi(x), \xi(x+L), \ldots, \xi(y))$. Hence, we have a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ and an
index set $J^{\prime}=J(x, y) \subset J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ such that $\left|J^{\prime}\right| \geq \exp \left(\gamma l_{1}\right)$ and $w^{2}(j)$ is constant on $J^{\prime}$. Assume $(x, y)$ is a barrier. Then $J^{\prime}=J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$. Suppose not. Then there exists $j^{*} \in$ $J\left(I^{*}, I, q^{*}, q, g^{*}, g\right) \backslash J^{\prime}$. This means that $j^{*}$ satisfies (3.7), but $w^{2}\left(j^{*}\right) \neq(\xi(x), \xi(x+L), \ldots, \xi(y))$. The latter is possible only, if $S\left(T^{1}\left(j^{*}\right)\right)>x$ or $S\left(T^{1}\left(j^{*}\right)\right)<y$. Let $S\left(T^{1}\left(j^{*}\right)=z>x\right.$. The event $E_{\text {stop }}^{1}(\tau)$ implies (3.18) and then $z \in I_{1}$. Hence, there is $z \in I_{1}$ and $k^{*} \leq j^{*}$ such that $E_{\text {mistake-1 }}^{1}\left(z, x, k^{*}\right)$ holds. This is a contradiction with $E_{\text {no mistake }}^{1}$. Hence $J^{\prime}=J\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ and $(\xi(x), \xi(x+L), \ldots, \xi(y)) \in \mathcal{W}^{1}$.
Now it remains to show that there are enough barriers in $I_{1}$. This follows immediately from $B_{\text {enough barriers }}^{1}$.
Proof of (3.17):
It suffices to show that $E_{\text {only ladders }}^{1} \cap E_{\text {enough ladders }}^{1} \cap B_{\text {unique fit }}^{1}$ ensures that $\mathcal{S}\left(\chi^{12 \alpha l_{1}}, \tau, \xi \mid I^{o}\right)$ consists of one element that satisfies (3.13).
Consider the "puzzle-playing" algorithm formalized in Definition 3.2. We show that there is an unique way to combine the words from $\mathcal{W}^{1}$, i.e. the solution set $\mathcal{S}$ is unique. Let $\exists \phi \in \mathcal{S}$ and let $D_{1} \subset D_{2} \subset \cdots \subset D_{n}$ be the sequence of sets ensured by the definition of $\phi$. By $\mathbf{1}, \phi \mid D_{1}$ is translated from a piece of $\xi \mid I_{1}$ by some $b$ satisfying $|b| \leq \exp \left(l_{1}\right)$, i.e. $\xi \mid I^{o}=T\left[\phi \mid D_{1}\right]$, where $T z=z+b$ is the translation and $I_{o} \subset\left[-e^{l_{1}}, e^{l_{1}}\right] \subset I_{1}$. We show: if $\phi \mid D_{i}$ is translated from a piece of $\xi \mid I_{1}$ by $b$, i.e. $\xi \mid J_{i}=T\left[\phi \mid D_{i}\right]$, for some $J_{i} \subset I_{1}$, then the same applies for $\phi \mid D_{i+1}$. Recall that $\phi \mid D_{i+1}$ and $\phi \mid D_{i}$ differ on $V_{i+1}$, only. By $\mathbf{3 c}$ ) and $E_{\text {only ladders }}^{1}$, $\phi\left|V_{i+1} \approx \xi\right| J(w)$ for some $J(w) \subset I_{1}$. Thus, there is an affine $T^{\prime}$ such that $\xi \mid J(w)=T^{\prime}\left[\phi \mid V_{i+1}\right]$ and, hence, there is a ladder interval $J^{\prime} \subset J(w)$ such that $\xi \mid J^{\prime}=T^{\prime}\left[\phi \mid\left(V_{i+1} \cap D_{i}\right)\right]$. So, $\phi \mid\left(V_{i+1} \cap D_{i}\right)$ is equivalent with some ladder word of $\xi \mid I_{1}$ by $T^{\prime}$. On the other hand, $\phi \mid\left(V_{i+1} \cap D_{i}\right)$ is translated by $b$, hence it is equivalent with some ladder word of $\xi \mid I_{1}$ by $T$. Let this word be $\xi \mid J$. Clearly $\xi|J \approx \xi| J^{\prime}$. By $\mathbf{3 b}$ ), the length of the ladder interval $V_{i+1} \cap D_{i}$ as well as $J^{\prime}$ and $J$ is at least $\frac{c_{1} l_{1}}{4}$. If $T \neq T^{\prime}$, then $J \neq J^{\prime}$, which contradicts $B_{\text {unique fit }}^{1}$. Hence, $T^{\prime}=T$ and $\phi \mid V_{i+1}$ is translated from a piece of $\xi \mid I_{1}$ by $b$ and $\phi \mid D_{i+1}$ is translated from a piece of $\xi \mid I_{1}$ by $b$ as well. The same holds for $\phi$, i.e. $\phi \equiv \xi \mid I(\phi)$ for some interval $I(\phi)$. By 4, $I(\phi)=\left[a_{o}-2 \exp \left(l_{1}\right), a_{o}+2 \exp \left(l_{1}\right)\right]$, where $I_{o}:=\left[a_{o}, b_{o}\right]$. So, $\phi$ is obtained from a fixed piece of scenery $\xi \mid I(\phi)$ by a fixed translation, $T$. Clearly such a $\phi$ is unique.
Let us show that $\phi$ satisfies (3.13). Since $\left|a_{o}\right| \leq \exp \left(l_{1}\right)$, we have that

$$
\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right] \subset I(\phi) \subset\left[-3 \exp \left(l_{1}\right), 3 \exp \left(l_{1}\right)\right]
$$

This means

$$
\xi\left|\left[-\exp \left(l_{1}\right), \exp \left(l_{1}\right)\right] \sqsubseteq \phi \sqsubseteq \xi\right|\left[-3 \exp \left(l_{1}\right), 3 \exp \left(l_{1}\right)\right],
$$

i.e. (3.13) holds.

It remains to show that $\mathcal{S}$ is not empty. Consider again the "puzzle playing" algorithm. Let $D_{i}, i \geq 1$ be the domain of $\phi_{i}$ at state $i$. It suffices to show that there exists $V_{i+1} \in \mathcal{L}\left(c_{1} l_{1}\right)$ satisfying all requirements of 3 and such that $\left|V_{i+1} \backslash D_{i}\right| \geq \frac{c_{1} l_{1}}{2}$. Note that $D_{i}=\cup_{j=0}^{L-1} I(j)$, where $I(j) \subset I_{1}$ is a ladder interval with length at least $c_{1} l_{1}$. Fix $j$ and let $a_{j}<b_{j}$ be the endpoints of $I(j)$. Consider the ladder interval $I^{b}(j):=\left(b_{j}-2\left(c_{1} l_{1} / 4\right) L, \ldots, b_{j}-1\left(c_{1} l_{1} / 4\right) L\right) \subset I(j)$. By $E_{\text {enough ladders }}^{1}$ there exists $z \in I^{b}(j)$ such that a ladder word of $\xi \mid V(z)$, with $V(z)=\left(z, z+L, \ldots, z+\left(c_{1} l_{1}\right) L\right) \in \mathcal{L}\left(c_{1} l_{1}\right)$ belongs to $\mathcal{W}^{1}$. Let this word be $w(z)$. Clearly, $\left|V(z) \cap D_{i}\right| \geq \frac{c_{1} l_{1}}{4}$. By $B_{\text {unique fit }}^{1}, w(z)$ is not a ladder word of any ladder piece $\phi_{i} \mid V_{j}, j=1, \ldots, i$. Hence $w(z)$ and $V(z)$ can be taken as $w_{i+1}$ and $V_{i+1}$. The same argument applies for $I^{a}(j):=\left(a_{j}+1\left(c_{1} l_{1} / 4\right) L, \ldots, a_{j}+2\left(c_{1} l_{1} / 4\right) L\right)$, implying that $D_{i}$ can be efficiently enlarged in other direction as well.

### 3.5 Probabilities for main theorem

### 3.5.1 Scenery-dependent events

At first, estimate the probabilities of $B$-events. These events depend on $\xi$, only. Note that all exponential bounds are valid for $l_{1}$ being big enough.

Estimate $P\left(\left(B_{\text {intervals OK }}^{1}\right)^{c}\right)$

Let

$$
E:=\left\{\xi \mid[z, z+m l] \text { is } \mathrm{OK} \forall z \in I_{1}\right\}, \quad E^{*}:=\left\{\xi \mid[z, z-m l] \text { is } \mathrm{OK}^{*} \forall z \in I_{1}\right\} .
$$

Now, by translation invariancy of $\xi$ and Theorem 2.1, it holds that for $l_{1}$ big enough

$$
P\left(E^{c}\right) \leq \sum_{z \in I_{1}} P(\xi \mid[z, z+m l] \text { is not } \mathrm{OK}) \leq 2 e^{3 l_{1}} P\left(E_{\mathrm{OK}}^{c}\right) \leq 2 \exp \left[3 l_{1}-a l\right]
$$

Similarly,

$$
P\left(E^{* c}\right) \leq \sum_{z \in I_{1}} P\left(\xi \mid[z-m l, z] \text { is not } \mathrm{OK}^{*}\right) \leq 2 e^{3 l_{1}} P\left(E_{\mathrm{OK}}^{c}\right) \leq 2 \exp \left[3 l_{1}-a l\right]
$$

Hence, if $l_{1}$ is sufficiently big, then

$$
\begin{equation*}
P\left(\left(B_{\text {intervals OK }}^{1}\right)^{c}\right) \leq 4 \exp \left[\left(3-a l_{2}\right) l_{1}\right] \tag{3.21}
\end{equation*}
$$

The following proposition also specifies the choice of $c_{1}$.

Proposition 3.1 There exists constants $C_{1}(n)$ and $k_{1}, k_{2}, k_{3}>0$ not depending on $l_{1}$ such that for $c_{1}>C_{1}(n)$ it holds:

$$
\begin{align*}
& P\left(\left(B_{\text {unique fit }}^{1}\right)^{c}\right) \leq \exp \left[-k_{1} l_{1}\right]  \tag{3.22}\\
& P\left(\left(B_{\text {recon straight }}^{1}\right)^{c}\right) \leq \exp \left[-k_{2} l_{1}\right]  \tag{3.23}\\
& P\left(\left(B_{\text {enough barriers }}^{1}\right)^{c}\right) \leq \exp \left[-k_{3} l_{1}\right] \tag{3.24}
\end{align*}
$$

provided $l_{1}$ is big enough.
Proof. It follows from Lemma 6.33 in [LMM] that for some constants $a_{1}, a_{2}$ depending on $L$, only, the bound $P\left(\left(B_{\text {unique fit }}^{1}\right)^{c}\right) a_{1} \leq \exp \left[-a_{2} l_{1}\right]$ is valid. Also, there is a fixed constant $C_{r}$ such that $a_{2}>0$ if $c_{1}>C_{r}$. This implies (3.22) for $l_{1}$ sufficiently big.
$\underline{\text { Estimate } P\left(\left(B_{\text {recon straight }}^{1}\right)^{c}\right)}$
Let $\mathcal{R}\left(l_{1} c_{1}\right)(x, y):=\left\{\mathcal{R}\left(l_{1} c_{1}\right)(x, y): R(0)=x, R\left(l_{1} c_{1} L\right)=y\right\}$. Thus $\mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ is (possibly empty) the set of admissible path from $x$ to $y$ with $l_{1} c_{1}$ steps. Fix $x, y$ such that $|y-x|<\left(l_{1} c_{1}\right) L$. At first note: if $l_{1}$ is big enough, then (for any value of $\left.C_{1} \geq 1\right) \mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ is either empty or has cardinality at least 2 . Any admissible path $R \in \mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ is a sequence $R=\left(t_{1}, \ldots, t_{c_{1} l_{1}}\right)$ of steps, where $\left|t_{i}\right| \leq L$. Hence, there exists a $R=\left(t_{1}, \ldots, t_{c_{1} l_{1}}\right) \in \mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ such that $t_{i} \neq t_{1}$ for a $i=2, \ldots c_{1} l_{1}$ (if no, then $\mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ would consists of one path, only). Let $R$ be one of such paths. Let $c_{1} \geq\left\ulcorner\frac{100}{2 L+1}\right\urcorner$. The number of possible steps is bounded by $2 L+1$. Hence, there is a step $t$ that occurs in $R$ at least $k:=100 l_{1}$ times. Formally, $\exists t \in\{-L, \ldots, L\}$ such that $\left|\left\{i=1, \ldots, c_{1} l_{1}: t_{i}=t\right\}\right| \geq k$. Any rearrangement of the order of steps in $R$ corresponds to another path in $\mathcal{R}\left(l_{1} c_{1}\right)(x, y)$. We consider two rearrangements of $R$. The first, $R^{1}$, starts with $k l_{1}$ steps of size $t$. Thus $R^{1}=\left\{t_{1}^{1}, \ldots, t_{c_{1} l_{1}}^{1}\right\} \in \mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ is such that $t_{1}^{1}=\cdots=t_{k}^{1}=t$. Let $u$ be another step if $R$ such that $u \neq t$. The second path, $R^{2}$, starts with $u$, and then is followed by $k$-steps of size $t$. Formally, $R^{2}=\left\{t_{1}^{2}, \ldots, t_{c_{1} l_{1}}^{2}\right\} \in \mathcal{R}\left(l_{1} c_{1}\right)(x, y)$ is such that $t_{1}^{2}=u, t_{2}^{2}=\cdots=t_{k+1}^{2}=t$. We now
estimate the probability that the paths $R^{1}$ and $R^{2}$ generate the same word in observation; we estimate

$$
\begin{aligned}
& P\left(\xi \circ R^{1}=\xi \circ R^{2}\right) \leq P((\xi(x+t), \ldots, \xi(x+k t))=(\xi(x+u), \xi(x+u+t), \ldots, \xi(x+u+k t))) \\
& \leq P(\xi(x+t)=\xi(x+u)) P(\xi(x+2 t)=\xi(x+u+t) \mid \xi(x+t)=\xi(x+u)) \times \\
& \times P(\xi(x+3 t)=\xi(x+u+2 t) \mid \xi(x+t)=\xi(x+u)), \xi(x+2 t)=\xi(x+u+t)) \cdots \\
& \cdots P(\xi(x+k t)=\xi(x+u+(k-1) t) \mid \xi(x+t)=\xi(x+u) \cdots \xi(x+(k-1) t)=\xi(x+u+(k-2) t)) \\
& \leq 2^{-k}=\exp \left[-100 \ln 2 l_{1}\right]
\end{aligned}
$$

Now,

$$
\begin{gathered}
E_{\text {recon straight }}=\bigcup_{x, y \in I_{1},|x-y|<l_{1} c_{1}} E_{\text {recon straight }}(x, y) \\
P\left(\left(E_{\text {recon straight }}\right)^{c}\right) \leq \sum_{x, y \in I_{1}} P\left(E_{\text {recon straight }}(x, y)\right) \leq 4 \exp \left(6 l_{1}\right) \exp \left[-100 \ln 2 l_{1}\right] \leq \exp \left[-50 l_{1}\right]
\end{gathered}
$$

$\underline{\text { Estimate } P\left(\left(B_{\text {enough barriers }}^{1}\right)^{c}\right)}$
For each $z, j$ define

$$
B_{\text {enough barriers }}^{1}(z, j):=\left\{\begin{array}{c}
\text { there exists } x \in\left[z, z+\left(\frac{c_{1} l_{1}}{4}\right) L\right] \text { such that } x \bmod L=j \\
\text { and } \left.\left(x, x+\left(c_{1} l_{1}\right)\right) L\right\} \text { is a barrier of } \xi
\end{array}\right\}
$$

Define

$$
B(x):=\left\{\left(x, x+\left(c_{1} l_{1}\right) L\right) \text { is a barrier of } \xi\right\}, \quad Y_{x}:=I_{B(x)}
$$

Note, if $x^{\prime}-x \geq 3 m^{3} L=: r$, then, by the definition, the events $B(x)$ and $B\left(x^{\prime}\right)$ are independent. Clearly the probability of $B(x)$ does not depend on $x$, let us denote $p=P(B(x))$. By definition, $p>2^{-3 m^{3} L}$. Denote $w=\left\llcorner\frac{c_{1} l_{1}}{4 r}-\frac{L}{r}\right\lrcorner>\frac{c_{1}-4 L}{4 r} l_{1}$ and use Höffding's inequality again

$$
\begin{aligned}
P\left(\left(B_{\text {enough barriers }}^{1}(z, j)\right)^{c}\right) & =P\left(\sum_{x=z+j}^{\frac{z+c_{1} l_{1}}{4}} Y_{x}=0\right) \leq P\left(\sum_{k=0}^{\frac{c_{1} l_{1}}{4 r}} Y_{r k+z+j}=0\right) \\
& \leq P\left(\sum_{k=0}^{w}\left(Y_{r k+z+j}-p\right) \leq w p\right) \leq 2 \exp \left[-2 w p^{2}\right] \\
& \leq 2 \exp \left[-2 \frac{c_{1}-4 L}{4 r} 2^{-6 m^{3} L} l_{1}\right]=2 \exp \left[-k_{2}^{\prime} l_{1}\right]
\end{aligned}
$$

for $k_{2}^{\prime}:=\frac{c_{1}-4 L}{4 r} 2^{-\left(6 m^{3} L+1\right)}$. Obviously, $k_{2}^{\prime}>0$, if $c_{1}>4 L$. Thus

$$
P\left(\left(B_{\text {enough barriers }}^{1}\right)^{c}\right) \leq \sum_{z \in I_{1}, j=\{0, \ldots, L\}} P\left(\left(B_{\text {enough barriers }}^{1}(z, j)\right)^{c}\right) \leq 8 \exp \left[\left(6-k_{2}^{\prime}\right) l_{1}\right] \leq \exp \left[-l_{1}\right]
$$

if $k_{2}^{\prime} \geq 8$. The latter implies $c_{1}-4 L \geq 4 \cdot 2^{6 m^{3}+6}$ or $c_{1} \geq r 2^{6 m^{3}+8}+4 L=3 m^{3} L 2^{6 m^{3}+8}+4 L$.
Hence, Proposition 3.1 holds with $C_{1}(n):=\max \left\{C_{r},\left\ulcorner\frac{100}{2 L+1}\right\urcorner, 3 m^{3} L 2^{6 m^{3}+8}+4 L\right\}$.

### 3.5.2 Random-walk depending events

$\underline{\text { Estimate } P\left(E_{\text {mistake-r }}^{1} \cap B_{\text {intervals OK }}^{1}\right) .}$
Fix $y, z \in I_{1}, z<y$ and note

$$
\begin{equation*}
E_{\text {mistake-r }}^{1}(z, y, k) \cap B_{\text {intervals OK }}^{1} \subset E_{\text {mistake-r }}^{1}(z, y, k) \cap\{\xi \mid[y, y+l m] \text { is OK }\}, \quad k=1,2, \ldots \tag{3.25}
\end{equation*}
$$

We now estimate the right side of (3.25). Recall the definitions of $T_{z}^{3}(k), w_{z}^{3}(k)$ and $g_{y}(\xi)$. Consider the events

$$
\begin{align*}
& E_{\text {mistake-r }}^{1}(y, z, k) \cap\{\xi \mid[y, y+l m] \text { is } \mathrm{OK}\}= \\
& \left\{\hat{q}\left(w_{z}^{3}(k)\right) \leq q_{y}(\xi), \hat{g}\left(w_{z}^{3}(k)\right) \sqsubseteq_{\mathcal{I}(\xi \mid[y, y+m l])} g_{y}(\xi), y \text { is a right barrier point, } \xi \mid[y, y+l m] \text { is } \mathrm{OK}\right\}, \tag{3.26}
\end{align*}
$$

$k=1,2, \ldots$. Because of (3.3), conditionally on $\xi$ the events (3.26) are independent and identically distributed. Hence, the events (3.26) all have the probability equal to

$$
\begin{equation*}
P\left(\hat{q}\left(\chi_{z}^{m^{2} l}\right) \leq q_{y}(\xi), \quad \hat{g}\left(\chi_{z}^{m^{2} l}\right) \preccurlyeq_{\mathcal{I}(\xi \mid[y, y+m l])} g_{y}(\xi), \quad y \text { is a right barrier point, } \quad \xi \mid[y, y+l m] \text { is } \mathrm{OK}\right) . \tag{3.27}
\end{equation*}
$$

The event in (3.27) depends on $\xi$, only. The distribution of $\xi$ is obviously translation invariant. Therefore, by Corollary $2.1,(3.27)$ can be estimated

$$
\begin{aligned}
& \left.P\left(\hat{q}\left(\chi_{z-y}^{m^{2} l}\right) \leq q_{0}(\xi)\right), \quad \hat{g}\left(\chi_{z-y}^{m^{2} l}\right) \preccurlyeq{\mathcal{I}\left(\xi^{m l}\right)} g_{0}(\xi), \quad 0 \text { is a right barrier point }, \quad \xi^{m l} \text { is OK }\right)= \\
& P\left(\left\{\hat{q}\left(\chi_{z-y}^{m^{2} l}\right) \leq q_{0}(\xi), \quad \hat{g}\left(\chi_{z-y}^{m^{2} l}\right) \preccurlyeq \mathcal{I}\left(\xi^{m l}\right) g_{0}(\xi)\right\} \cap E_{\text {origin }} \cap E_{O K}^{*}\right)= \\
& P\left(E_{\text {mistake }}(z-y) \cap E_{O K}^{*}\right) \leq \exp \left(-l \alpha_{I}\right),
\end{aligned}
$$

provided $l_{1}$ is big enough. Therefore,

$$
\begin{align*}
P\left(E_{\text {mistake-r }}^{1} \cap B_{\text {intervals OK }}^{1}\right) & \leq \sum_{y, z, k} P\left(E_{\text {mistake-r }}^{1}(y, z, k) \cap B_{\text {intervals OK }}^{1}\right) \\
& \leq \sum_{y, z, k} \exp \left(-l \alpha_{I}\right)<4 \exp \left[(6+\alpha) l_{1}-\alpha_{I} l\right] \tag{3.28}
\end{align*}
$$

The sum here is taken over all $z, y \in I_{1}, z<y$ and $k=1, \ldots, \exp \left(\alpha l_{1}\right)$.
Estimate $P\left(E_{\text {mistake-1 }}^{1} \cap B_{\text {intervals OK }}^{1}\right)$.
We need some additional notations. Recall $T_{z}^{1}(k)$. Now fix $x^{\prime} \in I_{1}$ and define $T_{z}^{1}\left(k_{i}\right), i=1,2, \ldots, N\left(x^{\prime}\right)$. as the $i$-th stopping time $T_{z}^{1}(k)$, for which $S\left(T_{z}^{1}(k)+\exp \left(2 l_{1}\right)\right)=x^{\prime}$. The indexes $k_{i}$ depend on chosen $x^{\prime}$. Define now

$$
E_{\text {mistake-1 }}^{1}\left(z, x, i, x^{\prime}\right):=\left\{\hat{q}^{*}\left(w_{z}^{1}\left(k_{i}\right)\right) \leq q_{x}^{*}(\xi)\right\} \cap\left\{\hat{g}^{*}\left(w_{z}^{1}\left(k_{i}\right)\right) \preccurlyeq \mathcal{I}^{*}(\xi \mid[x-m l, x]) g_{x}^{*}(\xi)\right\} \cap\{x \text { is a left b. p. }\}
$$

$i=1,2, \ldots, N\left(x^{\prime}\right)$.
Clearly, for each $k$ there exist $i, x^{\prime}$ such that $E_{\text {mistake-1 }}^{1}(z, x, k)=E_{\text {mistake-1 }}^{1}\left(z, x, i, x^{\prime}\right)$. The counterpart of (3.25) is

$$
E_{\text {mistake-1 }}^{1}\left(z, x, i, x^{\prime}\right) \cap B_{\text {intervals OK }}^{1} \subset E_{\text {mistake-r }}^{1}\left(z, x, i, x^{\prime}\right) \cap\left\{\xi \mid[x-l m, x] \text { is } \mathrm{OK}^{*}\right\}=: E\left(i, x^{\prime}\right)
$$

$i=1,2, \ldots, N\left(x^{\prime}\right)$.
As previously, we observe that $P\left(E\left(i, x^{\prime}\right)\right)$ is equal to
$P\left(\hat{q}^{*}\left(\chi_{x^{\prime}}^{m^{2} l}\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(\chi_{x^{\prime}}^{m^{2} l}\right) \preccurlyeq_{\mathcal{I}^{*}(\xi \mid[x-m l, x])} g_{x}^{*}(\xi), S_{x^{\prime}}\left(m^{2} l\right)=z, x\right.$ is a left b. p., $\xi \mid[x-l m, x]$ is $\left.\mathrm{OK}^{*}\right)$.
To calculate (3.29), at first note the following. Let $R(i), i=0,1, \ldots, k$ be an admissible path such that $R(0)=x^{\prime}, R(k)=z$. Thus, for any scenery $\psi$, the observation $\chi \mid[0, k]$ equals $\psi(R(i)), i=0, \ldots, k$. This means, $(\chi \mid[0, k])^{-}=\psi\left(R^{-}(i)\right)$, where $R^{-}(i)=-R(k-i), i=0, \ldots, k$. By symmetry of $S$, any admissible
path $R[0, k]$ has the same probability as its reverse $R^{-}[0, k]$. This means that for any $u \in\{0,1\}^{k+1}$ and for any fixed scenery $\psi$ we have with $P_{\psi}(\cdot):=P(\cdot \mid \xi=\psi)$,

$$
P_{\psi}\left((\chi \mid[0, k])^{-}=u, S_{x^{\prime}}(k)=z\right)=P_{\psi}\left(\chi \mid[0, k]=u, S_{z}(k)=x^{\prime},\right)
$$

or

$$
P_{\psi}\left(\left(\chi_{x^{\prime}}^{k}\right)^{-}=u, S_{x^{\prime}}(k)=z\right)=P_{\psi}\left(\chi_{z}^{k}=u, S_{z}(k)=x^{\prime}\right)
$$

By symmetry, again, the left side of last equality equals

$$
P_{-\psi}\left(\chi_{-z}^{k}=u, S_{-z}(k)=-x^{\prime}\right)
$$

In particular, since $(\psi \mid[x-l m, x])^{-}=-\psi \mid[-x,-x+l m]$
$P_{\psi}\left(\hat{q}\left(\left(\chi_{x^{\prime}}^{k}\right)^{-}\right) \leq q\left((\psi \mid[x-l m, x])^{-}\right), \hat{g}\left(\left(\chi_{x^{\prime}}^{k}\right)^{-}\right){\preccurlyeq \mathcal{I}\left((\psi \mid[x-l m, x])^{-}\right)} g\left((\psi \mid[x-l m, x])^{-}\right), S_{x^{\prime}}(k)=z\right)=$ $P_{-\psi}\left(\hat{q}\left(\chi_{-z}^{k}\right) \leq q(-\psi \mid[-x,-x+l m]), \hat{g}\left(\left(\chi_{-z}^{k}\right)\right) \preccurlyeq \mathcal{I}(-\psi \mid[-x+l m,-x]) g(-\psi \mid[-x+l m,-x]), S_{-z}(k)=-x^{\prime}\right)$.

Recall the definitions of $\hat{q}^{*}, q^{*}, \hat{g}^{*}, g^{*}$. Clearly $x$ is a left barrier point for $\psi$ if and only if $-x$ is a right barrier point for $-\psi$ and, by definition, $\psi \mid[x-l m, x]$ is $\mathrm{OK}^{*}$ if and only $(\psi \mid[x-l m, x])^{-}=-\psi \mid[-x,-x+l m]$ is OK. Let

$$
\begin{aligned}
A^{*}(x) & :=\left\{x \text { is a left b. p. of } \psi, \psi \mid[x-l m, x] \text { is } \mathrm{OK}^{*}\right\} \\
A(x) & :=\{x \text { is a right b. p. of } \psi, \psi \mid[x, x+l m] \text { is OK }\} .
\end{aligned}
$$

Thus, for each $\psi$,

$$
\begin{aligned}
& P_{\psi}\left(\hat{q}^{*}\left(\chi_{x^{\prime}}^{k}\right) \leq q^{*}(\psi \mid[x-l m, x]), \hat{g}^{*}\left(\chi_{x^{\prime}}^{k}\right) \preccurlyeq \mathfrak{I}^{*}(\psi \mid[x-l m, x]) g^{*}(\psi \mid[x-l m, x]), S_{x^{\prime}}(k)=z\right) I_{A^{*}(x)}(\psi)= \\
& P_{-\psi}\left(\hat{q}\left(\chi_{-z}^{k}\right) \leq q(-\psi \mid[-x,-x+l m]), \hat{g}\left(\chi_{-z}^{k}\right) \preccurlyeq \mathcal{I}(-\psi \mid[-x+l m,-x])\right. \\
& \times I_{A(-x)}(-\psi) .
\end{aligned}
$$

Finally, integrate over $\xi$ and use the fact that $\xi$ and $-\xi$ have the same distribution to get

$$
\begin{align*}
& P\left(\hat{q}^{*}\left(\chi_{x^{\prime}}^{k}\right) \leq q^{*}(\xi \mid[x-l m, x]), \hat{g}^{*}\left(\chi_{x^{\prime}}^{k}\right) \preccurlyeq \mathfrak{I}^{*}(\xi \mid[x-l m, x]) g^{*}(\xi \mid[x-l m, x]), S_{x^{\prime}}(k)=z, \xi \in A^{*}(x)\right)= \\
& P\left(\hat{q}\left(\chi_{-z}^{k}\right) \leq q(\xi \mid[-x,-x+l m]), \hat{g}\left(\chi_{-z}^{k}\right) \preccurlyeq \mathcal{I}(\xi \mid[-x+l m,-x]) g(\xi \mid[-x+l m,-x]), S_{-z}(k)=-x^{\prime}, \xi \in A(-x)\right) \tag{3.30}
\end{align*}
$$

Now take $k=m^{2} l, y:=-x, z:=-z$ and sum over $x^{\prime}$ to obtain that $P\left(E\left(i, x^{\prime}\right)\right)$ equals

$$
P\left(\hat{q}\left(\chi_{z}^{m^{2} l}\right) \leq q_{y}(\xi), \hat{g}\left(\chi_{z}^{m^{2} l}\right) \preccurlyeq \mathcal{I}(\xi \mid[y+l m, y]) g_{y}(\xi), y \text { is a right b. p., } \xi \mid[y, y+l m] \text { is } \mathrm{OK}\right) .
$$

Hence, $P\left(E\left(i, x^{\prime}\right)\right.$ equals (3.27) and, hence, it is bounded by $\exp \left(-l \alpha_{I}\right)$. This means

$$
\begin{align*}
P\left(E_{\text {mistake-1 }}^{1} \cap B_{\text {intervals OK }}^{1}\right) & \leq \sum_{y, z, i, x^{\prime}} P\left(E_{\text {mistake-r }}^{1}\left(y, z, i, x^{\prime}\right) \cap B_{\text {intervals OK }}^{1}\right) \\
& \leq \sum_{y, z, i, x^{\prime}} \exp \left(-l \alpha_{I}\right)<8 \exp \left[(9+\alpha) l_{1}-\alpha_{I} l\right] \tag{3.31}
\end{align*}
$$

where the sum is taken over all $z, y, x^{\prime} \in I_{1}, z<y$ and $i=1, \ldots, \exp \left(\alpha l_{1}\right)$.
Estimate $P\left(E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1} \cap\left(E_{\text {recon straight }}^{1}\right)^{c}\right)$

Fix a set of attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$ and consider random indexes $j_{1}, \ldots, j_{\kappa}$ as in (3.9) (3.12). They depend on chosen attributes. We consider the set $E^{c}$, where $E:=E_{\text {reconstraight }}^{1}\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$. On $E^{c}$, the following hold: $\kappa>\exp \left(\gamma l_{1}\right)$ and for every $k=1, \ldots, \exp \left(\gamma l_{1}\right)+1$, it holds $w^{2}\left(j_{k}\right)=w^{2}\left(j_{1}\right)$. Define

$$
Y_{k}:=1-I_{w^{2}\left(j_{1}\right)}\left(w^{2}\left(j_{k}\right)\right), \quad k=2, \ldots, \kappa
$$

Hence $Y_{k}=1$ if and only if $w^{2}\left(j_{k}\right) \neq w^{2}\left(j_{1}\right)$. Therefore, $E^{c}=\left\{\sum_{k=1}^{\exp \left(\gamma l_{1}\right)+1} Y_{k}=0\right\}$. We now consider the following $\sigma$-algebra

$$
\mathcal{A}:=\sigma\left(\xi(z), S(\tau(j)), S\left(T^{1}\left(j_{k}\right)\right), S\left(T^{3}\left(j_{k}\right)\right), z \in \mathbb{Z}, j=1, \ldots, \exp \left(\alpha l_{1}\right), k=1, \ldots, \kappa\right)
$$

Given $\mathcal{A}$, the values of $\kappa$ as well as $S\left(T^{1}\left(j_{k}\right)\right)=x_{k}$ and $S\left(T^{3}\left(j_{k}\right)\right)=y_{k}, k=1, \ldots, \kappa$ are known. This means that the random variables $Y_{1}, \ldots, Y_{\kappa}$ depend on the behavior of $S$ from $x_{k}$ to $y_{k}$ during $c_{1} l_{1}$ steps. Hence, given $\mathcal{A}$ the random variables $Y_{1}, \ldots, Y_{\kappa}$ are independent.
Consider now the events $E_{\text {stop }}^{1}(\tau)$ and $B_{\text {reconstraight }}^{1}$. Obviously they both belong to $\mathcal{A}$. Note that on $E_{\text {stop }}^{1}(\tau)$, it holds $x_{k}, y_{k} \in I_{1}$, for every $k=1, \ldots, \kappa$. Hence, if in addition also $B_{\text {recon straight }}^{1}$ holds, then for each $k=2, \ldots, \kappa$ there exists at least one admissible path from $x_{k}$ to $y_{k}$ that generates different words in observations. Recall the definition of $p_{\text {min }}$ and deduce that on $E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1}$ it holds $P\left(Y_{k}=1 \mid \mathcal{A}\right) \geq\left(p_{\min }\right)^{c_{1} l_{1}}, k=2, \ldots, \kappa$. Hence, by Höffding's inequality on $E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1}$ we get for a constant $d>0$

$$
\begin{equation*}
P\left(E^{c} \mid \mathcal{A}\right)=P\left(\sum_{k=2}^{\exp \left(\gamma l_{1}\right)+1} Y_{k}=0 \mid \mathcal{A}\right) \leq \exp \left[-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right] \tag{3.32}
\end{equation*}
$$

Indeed, for $Y_{1}, \ldots, Y_{e^{b}}$ independent Bernoulli random variables with $E\left(X_{i}\right) \geq e^{a}$, the Höffding's inequality states

$$
P\left(\sum_{i=1}^{e^{b}} Y_{i}=0\right)=P\left(\sum_{i=1}^{e^{b}}\left(Y_{i}-E Y_{i}\right) \leq-\sum_{i=1}^{e^{b}} E Y_{i}\right) \leq \exp \left[-d e^{-b}\left(\sum_{i=1}^{e^{b}} E Y_{i}\right)^{2}\right] \leq \exp \left[-d e^{b+2 a}\right]
$$

Now take $b=\gamma l_{1}, a=c_{1} l_{1} \ln \left(p_{\text {min }}\right)$ to obtain (3.32).
Integrate (3.32) over $E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1}$ to obtain

$$
\begin{equation*}
P\left(E^{c} \cap E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1}\right) \leq \exp \left[-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right] \tag{3.33}
\end{equation*}
$$

Finally, estimate

$$
P\left(\left(E_{\text {recon straight }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1}\right) \leq \sum_{\left(I^{*}, I, q^{*}, q, g^{*}, g\right)} E_{\text {recon straight }}^{1}\left(I^{*}, I, q^{*}, q, g^{*}, g\right)
$$

where the sum is taken over all attributes $\left(I^{*}, I, q^{*}, q, g^{*}, g\right)$. There are less than $2^{2\left(n^{2} l+l\right)} l^{4 l}$ attributes. Thus, the right side of the previous display is bounded by

$$
\begin{aligned}
& 2^{2\left(n^{2} l+l\right)} l^{4 l} \exp \left[-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right]= \\
& \exp \left[2\left(n^{2} l+l\right) \ln 2+(4 l) \ln l-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right]= \\
& \exp \left[l_{1}\left(2\left(n^{2} l_{2}+l_{2}\right) \ln 2+\left(4 l_{2}\right)\left(\ln l_{1}+\ln l_{2}\right)\right)-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right]
\end{aligned}
$$

So,

$$
\begin{align*}
& \left(E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1} \cap\left(E_{\text {recon straight }}^{1}\right)^{c}\right) \leq \\
& \leq \exp \left[l_{1}\left(2\left(n^{2} l_{2}+l_{2}\right) \ln 2+\left(4 l_{2}\right)\left(\ln l_{1}+\ln l_{2}\right)\right)-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\text {min }}\right) l_{1}\right)\right] \tag{3.34}
\end{align*}
$$

$\underline{\text { Estimate } P\left(E_{\text {stop }}^{1}(\tau) \cap\left(E_{\text {enough times }}^{1}\right)^{c} \cap B_{\text {intervals } O K}^{1}\right)}$
Let $p_{L}:=P(S(1)-S(0)=L)$ and define

$$
p^{*}:=\exp \left[-\left(1.5+2 \alpha_{I I} l_{2}+c_{1} \ln p_{L}\right) l_{1}\right]
$$

Proposition 3.2 If

$$
\begin{equation*}
\exp \left(\alpha l_{1}\right) p^{*} \geq 2 \exp \left(\gamma l_{1}\right) \tag{3.35}
\end{equation*}
$$

then

$$
\begin{equation*}
P\left(E_{\text {stop }}^{1}(\tau) \cap\left(E_{\text {enough times }}^{1}\right)^{c} \cap B_{\text {intervals } O K}^{1}\right) \leq 36 \exp \left[(2-\exp (2 \gamma-\alpha)) l_{1}\right] \tag{3.36}
\end{equation*}
$$

provided $l_{1}$ is big enough.
Proof. Recall the definitions of $T^{1}(j), T^{3}(j), j=1, \ldots \exp \left(\alpha l_{1}\right)$. Let $x, y \in I_{1}$ be such that $y=x+c_{1} l_{1} L$ and define

$$
E_{j}(x, y):=\left\{\begin{array}{c}
S\left(T^{1}(j)-l m^{2}\right)=x-l m \\
S\left(T^{1}(j)\right)=x, S\left(T^{3}(j)\right)=y \\
\hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \preccurlyeq \mathcal{I}^{*} g_{x}^{*}(\xi), \\
\hat{q}\left(w^{3}(j)\right) \leq q_{y}(\xi), \hat{g}\left(w^{3}(j)\right) \preccurlyeq \mathcal{I} g_{y}(\xi)
\end{array}\right\}, \quad Y_{j}:=I_{E_{j}}, \quad j=1, \ldots, e^{\alpha l_{1}} .
$$

Obviously,

$$
\begin{equation*}
\left\{\sum_{j=1}^{e^{\alpha l_{1}}} Y_{j} \geq e^{\gamma l_{1}}\right\} \subset E_{\text {enough times }}^{1}(x, y) \tag{3.37}
\end{equation*}
$$

For each $j$ and for every scenery $\psi$, it holds

$$
\begin{aligned}
& P_{\psi}\left(Y_{j}=1\right)=P_{\psi}\left(S\left(T^{1}(j)-l m^{2}\right)=x-l m\right) \times \\
& P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \preccurlyeq \mathcal{I}^{*} g_{x}^{*}(\xi) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m^{2}\right) \times \\
& P_{\psi}\left(S\left(T^{3}(j)\right)=y \mid S\left(T^{1}(j)-l m^{2}\right)=x-\operatorname{lm}, S\left(T^{1}(j)\right)=x, \hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \preccurlyeq \mathcal{I}^{*} g_{x}^{*}(\xi)\right) \times \\
& P_{\psi}\left(\hat{q}\left(w^{3}(j)\right) \leq q_{y}(\xi), \hat{g}\left(w^{3}(j)\right) \preccurlyeq \mathcal{I} g_{y}(\xi) \mid\right. \\
& \left.S\left(T^{1}(j)-l m^{2}\right)=x-l m^{2}, S\left(T^{1}(j)\right)=x, S\left(T^{3}(j)\right)=y, \hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \preccurlyeq \mathcal{I}^{*} g_{x}^{*}(\xi)\right) .
\end{aligned}
$$

Now, by the Markov property of $S$

$$
\begin{aligned}
P_{\psi}\left(Y_{j}=1 \mid E_{\text {stop }}(\tau)\right) & \left.=P_{\psi}\left(S\left(T^{1}(j)-l m^{2}\right)=x-\operatorname{lm} \mid E_{\text {stop }}(\tau)\right)\right) \\
& \times P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \preccurlyeq \mathcal{I}^{*} g_{x}^{*}(\xi) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m^{2}\right) \\
& \times P_{\psi}\left(S\left(T^{3}(j)\right)=y \mid S\left(T^{1}(j)\right)=x\right) \\
& \times P_{\psi}\left(\hat{q}\left(w^{3}(j)\right) \leq q_{y}(\xi), \hat{g}\left(w^{3}(j)\right) \preccurlyeq \mathcal{I} g_{y}(\xi) \mid S\left(T^{3}(j)\right)=y\right)
\end{aligned}
$$

Use local CLT to estimate

$$
\begin{aligned}
P_{\psi}\left(S\left(T^{1}(j)-l m^{2}\right)=x-l m \mid E_{\text {stop }}(\tau)\right) & =P_{\psi}\left(S\left(T^{1}(j)-l m^{2}\right)=x-l m \mid S\left(T^{1}(j)-l m^{2}-e^{2 l_{1}}\right)\right. \\
& \geq \exp \left(-1.5 l_{1}\right)
\end{aligned}
$$

By Theorem 2.3 and symmetry of $S$, it holds

$$
\begin{aligned}
& P_{\psi}\left(S\left(T^{1}(j)\right)=x, \hat{q}^{*}\left(w^{1}(j)\right) \leq q_{x}^{*}(\xi), \hat{g}^{*}\left(w^{1}(j)\right) \preccurlyeq \mathcal{I}^{*} g_{x}^{*}(\xi) \mid S\left(T^{1}(j)-l m^{2}\right)=x-l m^{2}\right) \geq \exp \left[-l \alpha_{I I}\right] \\
& P_{\psi}\left(\hat{q}\left(w^{3}(j)\right) \leq q_{y}(\xi), \hat{g}\left(w^{3}(j)\right) \preccurlyeq \mathcal{I} g_{y}(\xi) \mid S\left(T^{3}(j)\right)=y\right) \geq \exp \left[-l \alpha_{I I}\right]
\end{aligned}
$$

provided $\psi \in B_{\text {intervals OK }}^{1}$.
Finally,

$$
P_{\psi}\left(S\left(T^{3}(j)\right)=y \mid S\left(T^{1}(j)\right)=x\right)=\left(p_{L}\right)^{c_{1} l_{1} L} .
$$

This means, for $\psi \in B_{\text {intervals }}^{1}$ OK

$$
\begin{equation*}
P_{\psi}\left(Y_{j}=1 \mid E_{\text {stop }}(\tau)\right) \geq \exp \left[-1.5 l_{1}\right] \exp \left[-2 l \alpha_{I I}\right]\left(p_{L}\right)^{c_{1} l_{1} L}=p^{*} \tag{3.38}
\end{equation*}
$$

Conditional on $\left.E_{\text {stop }}(\tau)\right)$ and $\psi$, the random variables $Y_{i}$ are independent. That follows from the definition of $\left.E_{\text {stop }}(\tau)\right)$. Hence

$$
\begin{equation*}
P_{\psi}\left(\sum_{j=1}^{e^{\alpha l_{1}}} Y_{j}<e^{\gamma l_{1}} \mid E_{\text {stop }}(\tau)\right) \leq P\left(\sum_{j=1}^{e^{\alpha l_{1}}} Z_{j}<e^{\gamma l_{1}}\right)=P\left(\sum_{j=1}^{e^{\alpha l_{1}}}\left(Z_{j}-p^{*}\right)<e^{\gamma l_{1}}-e^{\alpha l_{1}} p^{*}\right), \tag{3.39}
\end{equation*}
$$

where $Z_{i}$ are independent Bernoulli random variables with parameter $p^{*}$. By (3.35), the right side of (3.39) is bounded by

$$
P\left(\sum_{j=1}^{e^{\alpha l_{1}}}\left(Z_{j}-p^{*}\right)<e^{\gamma l_{1}}-e^{\alpha l_{1}} p^{*}\right) \leq P\left(\sum_{j=1}^{e^{\alpha l_{1}}}\left(Z_{j}-p^{*}\right)<-e^{\gamma l_{1}} p^{*}\right) .
$$

Use Höffding's inequality to get

$$
P\left(\sum_{j=1}^{e^{\alpha l_{1}}}\left(Z_{j}-p^{*}\right)<-e^{\gamma l_{1}}\right) \leq \exp \left[-d e^{(2 \gamma-\alpha) l_{1}}\right] .
$$

Finally, integrate over $E_{\text {stop }}^{1}(\tau) \cap B_{\text {intervals }}^{1} O K$ and use (3.37) to deduce

$$
\left.P\left(E_{\text {stop }}^{1}(\tau) \cap\left(E_{\text {enough times }}^{1}(x, y)\right)^{c} \cap B_{\text {intervals Ок }}^{1}\right) \leq \exp [-\exp (2 \gamma-\alpha)) l_{1}\right] .
$$

Sum over all pairs $(x, y) \in I_{1}$ to get (3.36).

### 3.6 Tuning parameters

Recall that for big $n, \alpha_{I}>8 \alpha_{I I}$.

- Choose $n$ so big that statements of Theorem 2.1, Theorem 2.2, relation (2.26) and the statement of Corollary 2.1 hold.
- Then choose $c_{1}(n)>C_{1}(n)$, where $C_{1}(n)$ is specified in Proposition 3.1.
- Then choose $l_{2}\left(c_{1}, n\right)$ so big that simultaneously

$$
\begin{align*}
\alpha_{I I} l_{2} & >1.5+\ln 2-c_{1} \ln p_{L}  \tag{3.40}\\
\left(\alpha_{I}-7 \alpha_{I I}\right) l_{2} & >9  \tag{3.41}\\
4 \alpha_{I I} l_{2} & >-2 c_{1} \ln p_{\min }  \tag{3.42}\\
\alpha_{I I} l_{2} & >\ln 2  \tag{3.43}\\
a l_{2} & >3 \tag{3.44}
\end{align*}
$$

- Then take $\gamma\left(n, c_{1}, l_{2}\right)=4 \alpha_{I I} l_{2}$
- Then take $\alpha\left(n, c_{1}, l_{2}\right)=7 \alpha_{I I} l_{2}$


### 3.7 Proof of the main theorem

Recall Lemma 3.1. By (3.15), (3.16) and (3.17), for $l_{1}$ big enough, it holds

$$
\begin{align*}
& P\left(\left(E_{\text {alg works }}^{1}(\tau)\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right) \leq \\
& P\left(\left(E_{\text {only ladders }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right)+P\left(\left(E_{\text {all ladders }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right)+P\left(\left(B_{\text {unique fit }}^{1}\right)^{c}\right) ;  \tag{3.45}\\
& P\left(\left(E_{\text {only ladders }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right) \leq P\left(\left(E_{\text {recon straight }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right) \leq  \tag{3.46}\\
& P\left(\left(E_{\text {recon straight }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau) \cap B_{\text {recon straight }}^{1}\right)+P\left(\left(B_{\text {recon straight }}^{1}\right)^{c}\right) ; \\
& P\left(\left(E_{\text {all ladders }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right) \leq P\left(\left(B_{\text {enough barriers }}^{1}\right)^{c}\right) \\
& +P\left(\left(E_{\text {no mistake }}^{1}\right)^{c}\right)+P\left(\left(B_{\text {enough paths }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right) ;  \tag{3.47}\\
& P\left(\left(E_{\text {no mistake }}^{1}\right)^{c}\right) \leq P\left(\left(E_{\text {no mistake }}^{1}\right)^{c} \cap B_{\text {intervals OK }}^{1}\right)+P\left(\left(B_{\text {intervals OK }}^{1}\right)^{c}\right) ;  \tag{3.48}\\
& P\left(\left(B_{\text {enough times }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau)\right) \leq P\left(\left(B_{\text {enough times }}^{1}\right)^{c} \cap E_{\text {stop }}^{1}(\tau) \cap B_{\text {intervals OK }}^{1}\right)+P\left(\left(B_{\text {intervals OK }}^{1}\right)^{c}\right) . \tag{3.49}
\end{align*}
$$

Recall the definitions of $l_{2}$. The condition (3.45) states $7 \alpha_{I I} l_{2}>4 \alpha_{I I} l_{2}+1.5+\ln 2-c_{1} \ln p_{L}+2 \alpha_{I I} l_{2}$ or, equivalently,

$$
\alpha l_{1}>\left(\gamma+1.5+\ln 2-c_{1} \ln p_{L}\right) l_{1}+2 \alpha_{I I} l
$$

Taking exponentials,

$$
\exp \left(\alpha l_{1}\right) \exp \left(-1.5 l_{1}-2 \alpha_{I I} l\right)\left(p_{L}\right)^{c_{1} l_{1}}>2 \exp \left(\gamma l_{1}\right) .
$$

Recall the definition of $p^{*}$ and note that the inequality in the previous display is (3.35). Hence, by Proposition 3.2, we have the bound (3.36). By (3.43), $k_{4}:=\exp (2 \gamma-\alpha)=\exp \left(\alpha_{I I} l_{2}\right)>2$, implying that (3.36) is exponentially small in $l_{1}$. By (3.44), there exist $k_{5}>0$ such that (3.21) is bounded by $4 \exp \left[-k_{5} l_{1}\right]$. With (3.36), we obtain that (3.49) is bounded by $40 \exp \left[-\left(k_{4} \wedge k_{5}\right) l_{1}\right]$.

Use (3.31) and (3.28) with (3.41) to obtain that $P\left(\left(E_{\text {no mistake }}^{1}\right)^{c} \leq 12 \exp \left[(9+\alpha) l_{1}-\alpha_{I} l\right]=12 \exp \left[-k_{6} l_{1}\right]\right.$ for a $k_{6}>0$. Hence, (3.48) is bounded by $12 \exp \left[-k_{6} l_{1}\right]+4 \exp \left[-k_{5} l_{1}\right] \leq 16 \exp \left[-\left(k_{6} \wedge k_{5}\right) l_{1}\right]$.

By (3.24), we now get that (3.47) is bounded by $40 \exp \left[-\left(k_{4} \wedge k_{5}\right) l_{1}\right]+16 \exp \left[-\left(k_{6} \wedge k_{5}\right) l_{1}\right]+\exp \left[-k_{3} l_{1}\right] \leq$ $56 \exp \left[-k_{7} l_{1}\right]$ for a $k_{7}>0$.

The requirement (3.42) states that $\gamma+2 c_{1} \ln p_{\text {min }}>0$ implying that

$$
\exp \left[l_{1}\left(2\left(n^{2} l_{2}+l_{2}\right) \ln 2+\left(4 l_{2}\right)\left(\ln l_{1}+\ln l_{2}\right)\right)-d \exp \left(\left(\gamma+2 c_{1} \ln p_{\min }\right) l_{1}\right)\right] \leq \exp \left[k_{8} l_{1}\right]
$$

for $l_{1}$ big enough. This means (3.46) is bounded by $\exp \left[-k_{9} l_{1}\right]$ for $l_{1}$ big enough.
Finally, we get that (3.45) is bounded by $\exp \left[-k l_{1}\right]$, if $l_{1}$ is big enough. This proves Theorem 1.1.

## 4 Appendix

### 4.1 Proof of Theorem 2.1

Recall $m(n)>n$.
For each $i=1, \ldots, l$ random cells $\xi_{i}=\xi \mid D_{i}=\left(\xi\left(d_{i-1}\right), \ldots, \xi\left(d_{i}\right)\right)$.
Consider the event $E_{\text {OKa }}$. We can rewrite

$$
E_{\mathrm{OK} a}=\left\{\sum_{i=2 L m^{2}}^{l} X_{i} \leq l 2 \epsilon(n)\right\}
$$

where $X_{i}$ is Bernoulli random variable that is one iff $\xi_{i}$ is not weak-OK. Let

$$
l_{*}:=L m^{2}+c+2, \quad l^{*}=l-c+1
$$

Then $\left(l_{*}-1\right) m-c m=L m^{3}+m$ and $\left(l^{*}-1\right) m+c m=l m$. Clearly $P\left(X_{i}=1\right) \leq \epsilon(n)$, if $l_{*} \leq i \leq l^{*}$. If $i>l^{*}$, then, by definition, $\xi_{i}$ cannot be weak-OK and, hence, $X_{i}=1$. Now, let $n$ be so big that $l_{*} \leq 2 L m^{2}$ i.e. $c+2 \leq L m^{2}$. This means, $E_{\mathrm{OK} a}$ is independent on $\xi^{L m^{3}}$. Then also $l-l^{*}=c-1 \leq 2 L m^{2}$. Let us estimate

$$
\begin{aligned}
& E_{\mathrm{OK}_{a}^{c}}=\left\{\sum_{i=2 L m^{2}}^{l} X_{i}>l 2 \epsilon(n)\right\} \subset\left\{\sum_{i=2 L m^{2}}^{l-2 L m^{2}} X_{i}>l 2 \epsilon(n)-2 L m^{2}\right\} \\
& \subset \bigcup_{j=-c+1}^{c}\left\{\sum_{k=k_{*}}^{k^{*}} X_{i k 2 c-j}>\frac{l 2 \epsilon(n)-2 L m^{2}}{2 c}\right\} \\
& \subset \bigcup_{j=-c+1}^{c}\left\{\sum_{k=k_{*}}^{k^{*}} X_{k 2 c-j}-\left(k^{*}-k_{*}+1\right) \epsilon(n)>\frac{l 2 \epsilon(n)-2 L m^{2}}{2 c}-\frac{l \epsilon(n)}{2 c}\right\} .
\end{aligned}
$$

Here $k_{*}:=\left\ulcorner\frac{2 L m^{2}+c}{2 c}\right\urcorner$ and $k^{*}:=\left\llcorner\frac{l-2 L m^{2}-c+1}{2 c}\right\lrcorner$. Thus $k^{*}-k_{*} \leq \frac{l-4 L m^{2}+1}{2 c}<\frac{l}{2 c}, k^{*}-k_{*}+1<l$. Note, by definition $X_{i} \in \sigma\left(\xi_{j} \mid j=i-c, i-c+1, \ldots, i+c-1\right)$. Thus, $X_{k}$ and $X_{k 2 c}$ are independent. This means, for each $j$ we can apply Höffding's inequality. Thus, for each $j$

$$
\begin{aligned}
P\left(\sum_{k=k_{*}}^{k^{*}}\left(X_{k 2 c-j}-\epsilon(n)\right)>\frac{l \epsilon(n)-2 L m^{2}}{2 c}\right) & \leq P\left(\sum_{k=k_{*}}^{k^{*}}\left(X_{k 2 c-j}-E X_{k 2 c-j}\right)>\frac{l \epsilon(n)-2 L m^{2}}{2 c}\right) \\
& \leq \exp \left[-\frac{\left(l \epsilon(n)-2 L m^{2}\right)^{2}}{c\left(k^{*}-k_{*}\right)}\right] \leq \exp \left[-\frac{l \epsilon^{2}(n)}{2 c}\right]
\end{aligned}
$$

provided $l$ is big enough to satisfy $l \epsilon(n)-2 L m^{2} \geq l \frac{\epsilon(n)}{2}$. Hence,

$$
\begin{equation*}
P\left(E_{\mathrm{OK}_{a}^{c}}^{c}\right) \leq 2 c \exp \left(\frac{-\epsilon^{2}(n) l}{2 c}\right) \leq \exp \left(-a_{1}(n) l\right) \tag{4.1}
\end{equation*}
$$

for some $a_{1}(n)>0$, provided $l$ is big enough.
We estimate $P\left(E_{\mathrm{OK}_{b}}{ }^{c}\right)$ by the same argument. Define

$$
E_{\mathrm{OK} b}^{i *}:=\left\{\left|\mathcal{I}_{I I}^{i}\left(\xi^{m l}\right)\right| \geq l\left(1-\exp \left(-m^{0.8}\right)\right)\right\}, \quad i=1,2 .
$$

Clearly, for $n$ big enough,

$$
\begin{equation*}
E_{\mathrm{OK} b}^{1 *} \cap E_{\mathrm{OK} b}^{2 *} \subset E_{\mathrm{OK} b} \quad \text { and } \mathrm{P}\left(E_{\mathrm{OK}}^{c}\right) \leq P\left(E_{\mathrm{OK} b}^{1 *}\right)+P\left(E_{\mathrm{OK} b}^{2 *}\right) \tag{4.2}
\end{equation*}
$$

Let us estimate $P\left(E_{\mathrm{OK} b}^{2 *}\right)$.

Let $Y_{i}$ be Bernoulli random variable that is $1 \mathrm{iff} \xi_{i}$ has not empty neighborhood. Let us estimate $P\left(Y_{i}=1\right)$. If $d_{i-1}-L m^{2} \geq 0$ and $d_{i}+L m^{2} \leq l m$, then

$$
\begin{aligned}
P\left(Y_{i}=1\right) & =\left(\exists j \in\left[d_{i-1}-L m^{2}, d_{i}+L m^{2}\right]: \xi(j)=\cdots=\xi\left(j+m^{0.9}\right)\right) \\
& \leq\left(2 L m^{2}+m+1\right)(0.5)^{m^{0.9}} \leq \exp \left(-m^{0.85}\right)
\end{aligned}
$$

in $m$ is big. Otherwise, by definition, $Y_{i}=1$. Let $N$ be such that the inequality above holds as well as (4.2) if $n>N$. Note that $E_{\text {OKb }}^{2 *}$ is independent of $\xi^{L m^{3}}$.

Clearly $Y_{i} \in \sigma\left(\xi_{i-L m}, \ldots, \xi_{i+L m}\right)$. Hence $Y_{i}$ and $Y_{i+2+2 L m}$ are independent. Let $k=2(1+L m)$. Now with $i^{*}=\left\llcorner\frac{l-2 L m^{2}-k+1}{k}\right\lrcorner$ and $i^{*} \leq \frac{l}{k}$ we get

$$
\begin{aligned}
E_{\mathrm{OK} b}^{2 *}{ }^{c} & =\left\{\sum_{i=2 L m^{2}}^{l} Y_{i}>l \exp \left(-m^{0.8}\right)\right\} \subset\left\{\sum_{i=2 L m^{2}}^{l-2 L m^{2}} Y_{i}>l \exp \left(-m^{0.8}\right)-2 L m^{2}\right\} \\
& \subset \bigcup_{j=0}^{k-1}\left\{\sum_{i=0}^{i^{*}} Y_{2 L m^{2}+j+i k}>\frac{l \exp \left(-m^{0.8}\right)-2 L m^{2}}{k}\right\} \\
& \subset \bigcup_{j=0}^{k-1}\left\{\sum_{i=0}^{i^{*}} Y_{2 L m^{2}+j+i k}-i^{*} \exp \left(-m^{0.85}\right)>\frac{l\left(\exp \left(-m^{0.8}\right)-\exp \left(-m^{0.85}\right)\right)-2 L m^{2}}{k}\right\} \\
& \subset \bigcup_{j=0}^{k-1}\left\{\sum_{i=0}^{i^{*}}\left(Y_{2 L m^{2}+j+i k}-E Y_{2 L m^{2}+j+i k}\right)>\frac{l\left(\exp \left(-m^{0.8}\right)-\exp \left(-m^{0.85}\right)\right)-2 L m^{2}}{k}\right\}
\end{aligned}
$$

Denote $\exp \left(-m^{0.8}\right)-\exp \left(-m^{0.85}\right)=: e(m)$ and apply Höffdings inequality

$$
P\left(\sum_{i=0}^{i^{*}}\left(Y_{2 L m^{2}+j+i k}-E Y_{2 L m^{2}+j+i k}\right) \geq \frac{l e(m)-2 L m^{2}}{k}\right) \leq \exp \left[-\frac{2(l e(m)-2 L m)^{2}}{l k}\right] \leq \exp \left[-a_{2}(m) l\right]
$$

for some $a_{2}(m)>0$, if $l$ is sufficiently big. Now, for big $l$,

$$
P\left(E_{\mathrm{OK} b}^{2 *}{ }^{c}\right) \leq 2(k+1) \exp \left(-a_{2}(m) l\right) \leq 2(m+1) \exp \left(-a_{2}(m) l\right) \leq \exp \left(-a_{3}(m) l\right)
$$

for some $a_{3}(m)>0$.
Similarly we estimate $P\left(E_{\mathrm{OKb}_{b}}^{1 *}\right)$.
Let $Z_{i}$ be Bernoulli random variable that is 1 iff $\xi_{i}$ is not isolated. If $i \geq l-L m$, then, by definition $Z_{i}=1$. Thus

$$
E_{\mathrm{OK} b}^{1 *{ }^{c}}=\left\{\sum_{i=2 L m^{2}}^{l} Z_{i}>l \exp (-m)\right\} \subset\left\{\sum_{i=2 L m^{2}}^{l-L m} Z_{i}>l \exp (-m)-L m\right\}
$$

Again, $E_{\mathrm{OK} b}^{1 *}$ is independent on $\xi^{L m^{3}}$. Note, if $\sum_{i+2 L m^{2}}^{l} Z_{i}>l \exp (-m)-L m$, then among the vectors $\left\{\xi_{2 L m^{2}-L m-1}, \xi_{2 L m^{2}-L m}, \ldots, \xi_{l}\right\}$ there exists at least $\frac{1}{2}(l \exp (-m)-L m-1)$ intervals $\xi_{i}$ without fence. Let $Z_{i}^{\prime}$ Bernoulli random variable that is 1 iff the srandom vector (but not the cell) $\xi \mid\left(d_{i-1}, d_{i}\right)$ does not contain a fence. Since the intervals $\left(d_{i-1}, d_{i}\right)$ and $\left(d_{j-1}, d_{j}\right)(i \neq j)$ are disjoint, $Z_{i}^{\prime}$ are iid random variables. Hence, with $j^{*}=2 L m^{2}-L m-1$, we get

$$
P\left(E_{\mathrm{OKb}}^{1 *}\right) \leq P\left(\sum_{j=j^{*}}^{l} Z_{j}^{\prime}>\frac{1}{2}(l \exp (-m)-L m-1)\right)
$$

Clearly

$$
P\left(Z_{i}^{\prime}=1\right)=P\left(\xi \mid\left(d_{i-1}, d_{i}\right) \text { contains no fence }\right) \leq\left(1-(0.5)^{2 L-1}\right)^{\frac{m-2}{2 L}}<e^{-c m}
$$

for some $c>0$. Now Höffding's inequality yields

$$
\begin{align*}
& P\left(\sum_{j=j^{*}}^{l} Z_{j}^{\prime} \geq \frac{1}{2}\left(l e^{-m^{0.8}}-L m\right)\right) \leq P\left(\sum_{j=1}^{l} Z_{j}^{\prime}-l e^{-c m} \geq \frac{1}{2}\left(l e^{-m^{0.8}}-L m\right)-l e^{-c m}\right)= \\
& P\left(\sum_{j=1}^{l}\left(Z_{j}^{\prime}-E Z_{j}^{\prime}\right)>\frac{1}{2} l\left(e^{-m^{0.8}}-2 e^{-c m}\right)-\frac{L}{2} m\right) \leq \exp \left[-\frac{\left(l\left(e^{-m^{0.8}}-2 e^{-c m}\right)-L m\right)^{2}}{2 l}\right] \tag{4.3}
\end{align*}
$$

The right side of (4.3) is bounded by $\exp \left(-l a_{4}(m)\right)$, for some $a_{4}(m)>0$, provided $l$ is big enough. Now, there exists $a_{5}(m)>0$ such that for big $l$,

$$
\begin{equation*}
P\left(E_{\mathrm{OK}}^{b}{ }_{b}^{c}\right) \leq \exp \left(-a_{3} l\right)+\exp \left(-a_{4} l\right) \leq \exp \left(-a_{5} l\right) \tag{4.4}
\end{equation*}
$$

Now, by (2.2), (4.1), (4.4)

$$
P\left(E_{\mathrm{OK}}^{c}\right) \leq P\left(E_{\mathrm{OK}_{a}^{c}}^{c}\right)+P\left(E_{\mathrm{OK}_{b}^{c}}^{c}\right) \leq \exp \left(-l a_{1}\right)+\exp \left(-l a_{5}\right) \leq \exp (-l a)
$$

for some $a(m)>0$ and $\operatorname{big} l$.

### 4.2 Proof of Proposition

By definition,

$$
E_{\text {min }}(i) \in \sigma\left(S(t)-S(t-1) \mid t \in\left[1,\left(s_{i}-r_{i}\right) m\right]\right)
$$

This means, if $E_{\text {min }}(i) \neq \emptyset$, then $P_{\psi}\left(E_{\text {min }}(i)\right) \geq\left(p_{\text {min }}\right)^{\left(s_{i}-r_{i}\right) m}$. We shall show that $E_{\text {min }}(i) \neq \emptyset$.
Let $i \in\{1, \ldots, k\}$. Let us describe an admissible path $R_{i} \in \mathcal{R}\left(\left(s_{i}-r_{i}\right) m\right)$ such that simultaneously satisfies (2.13), (2.14), (2.15). If such an path exists then and (2.16) holds.

Consider an arbitrary index-interval $\left[l_{2 i-1}, l_{2 i}\right], i>1$. It corresponds to the location-interval $\left[r_{i}, s_{i}\right]$. Let $C_{1}<\cdots<C_{q}$ be the big clusters of $\psi$ in $\left[s_{i}, r_{i}\right]$. Denote by $c_{j}, d_{j}, j=1, \ldots, q$ the beginnings and ends of big clusters, respectively. Hence, $C_{j} \subset\left[c_{j}, d_{j}\right]$. The path $R_{i}$ should read the big clusters as one block, i.e. along the reading-path.
Moreover, let $B_{1}<B_{2}<\cdots<B_{p}$ be the blocks of $\psi$ in the set $\left[s_{i}, r_{i}\right] \backslash\left(\cup_{j=1}^{q}\left[c_{j}+2, d_{j}-2\right]\right)$ that are bigger than $m^{2} / 2 \bar{v}$. By definition, $l\left(B_{j}\right)<m^{3}, j=1, \ldots, p$. Indeed, if $l\left(B_{j}\right) \geq m^{3}$, then $B_{j}$ would be a (part of) big cluster. We refer to a $B_{j}$ as a small block. The small blocks should be red as shortly as possible, i.e. along the reading path.
Finally let $A_{1}<A_{2}<\ldots<A_{K}, K=p+q$ be the ordered big clusters and small blocks. Let $a_{j}, b_{j}$ denote (an arbitrary) reading-beginning and reading-end of $A_{j}$.
Since $i>1$, it holds $l_{2 i-1} \in \mathcal{I}_{I I}$. Then $D_{2 i-1}$ has empty neighborhood, hence $\left[r_{i}, r_{i}+L m^{2}\right]$ is empty (for $\psi)$ and, therefore, does not contain any small blocks. Also $D_{2 i-1}$ is isolated. This implies that there is no point in $\left[r_{i}, r_{i}+L m^{2}\right]$ that is connected with any point in $\left[r_{i}+L m^{2}+m, s_{i}\right]$. In particular, all objects of interest, $A_{1}, \ldots, A_{K}$ are outside of $\left[r_{i}, r_{i}+L m^{2}\right]$ or, formally, $a_{1}>r_{i}+L m^{2}$.
If $s_{i}-r_{i} \leq 2 L m^{2}$, then the interval does not contain blocks that are bigger than $m^{0.9}$. In this case the path $R_{i}$ starts at $r_{i}$, i.e. $R(0)=r_{i}$ and goes to the point $s_{i}$ with $\left(l_{2 i}-l_{2 i-1}+1\right) m^{2}$ step without generating more than $m \bar{v}$ consecutive same colors in observations. This is clearly possible.
If $s_{i}-r_{i}>2 L m^{2}$, then we define the minimum-blocks path $R_{i}$ for interval $\left[r_{i}, s_{i}\right]$ backwards. More precisely, we define or prescribe a path $R^{*}$ that starts at $s_{i}$ and goes to $r_{i}$ with $\left(s_{i}-r_{i}\right) m^{2}$ steps. The prescription of $R^{*}$ is the following: start at $s_{i}$, i.e. $R^{*}(0)=s_{i}$. Then move stepwise to $b_{K}$ (recall, this is a reading-end of the last small block or the last big cluster in $\left.\left[r_{i}, s_{i}\right]\right)$. Recall $s_{i}=l_{2 i} m$ If $s_{i} \neq l$, then $l_{2 i} \in \mathcal{I}_{I I}$ and $\left[s_{i}-L m^{2}, s_{i}+m\right]$ is empty and $\left[s_{i}-L m^{2}-m, s_{i}-L m^{2}\right]$ contains a fence. As explained above, this implies that $b_{K} \leq s_{i}-L m^{2}$. So, by moving stepwise from $s_{i}$ to $b_{K}$, it is not possible that $S$ generates more than $m^{0.9} \bar{v}$ same colors in the beginning.
After reaching $b_{K}$ move along the reading path to $a_{K}$. Then move stepwise to $b_{K-1}$. Continue so until $a_{1}$
and then stepwise until $r_{i}+L m^{2}$. Since $a_{1}>r_{i}+L m^{2}$, for such a path less than $\left(\left(s_{i}-r_{i}\right)-L m^{2}\right) \bar{v}$ steps are needed. This means that the path has more than $\left(s_{i}-r_{i}\right)(m-\bar{v})+L m^{2} \bar{v}$ steps to cover the interval [ $r_{i}, r_{i}+L m^{2}$ ] with length $L m^{2}$ without generating more than $m \bar{v}$ consecutive same colors in observations and satisfying $R^{*}\left(\left(s_{i}-r_{i}\right) m\right)=r_{i}$. This is obviously possible, because the interval does not contain more than $m^{0.9}$ consecutive same colors. Finally define $R_{i}$ as $R^{*}$ backwards, i.e. $R_{i}(0)=R^{*}\left(\left(s_{i}-r_{i}\right) m\right)=$ $r_{i}, R_{k}(1)=R^{*}\left(\left(s_{i}-r_{i}\right) m-1\right), \ldots, R_{i}(j)=R^{*}\left(\left(s_{i}-r_{i}\right) m-j\right), \ldots, R_{k}\left(\left(s_{i}-r_{i}\right) m\right)=R^{*}(0)=s_{i}$ (recall, $S$ is symmetric).
Such definition of $R_{i}$ ensures that (2.13) and (2.15) hold. Let us show that (2.14) holds as well.
Note that the number of big blocks in $\psi \circ R_{i}$ is equal with the number of big clusters in $\left[r_{i}, s_{i}\right]$. Let this number be $M$. That means

$$
\hat{q}_{V}\left(\psi \circ R_{i}\right)=q_{V}\left(\left[r_{i}, s_{i}\right]\right)=M
$$

where $V:=l_{2 i}-l_{2 i-1}+1$. Let

$$
T(j):=\inf \left\{k: q_{k}\left(\psi \mid\left[r_{i}, s_{i}\right]\right)=j\right\}, \quad \hat{T}(j):=\inf \left\{k: \hat{q}_{k}\left(\psi \circ R_{i}\right)=j\right\} \quad j=1, \ldots, M
$$

Clearly, (2.14) is violated if there exists $j \in\{1, \ldots, M\}$ such that $\hat{T}(j)<T(j)$. Fix a $j \in\{1, \ldots, M\}$. The inequality $\hat{T}(j)<T(j)$ means that after reading the $j$-th big cluster, $R_{i}$ has more than $(V-T(j)+1) m^{2}$ steps to go to $s_{i}$. However, the path $R_{i}$ is constructed such that after reaching to the $b_{j}$ we have at most $(V-T(j)+1) m \bar{v}$ step to go $s_{i}$. That proves $(2.14)$.
Finally consider the first interval $\left[r_{1}, s_{1}\right]=\left[0, s_{1}\right]$ (obviously, $r_{1}=0$ ). Since $l_{1}=1 \notin \mathcal{I}_{I I}$, the interval $\left[0, L m^{2}\right]$ is not necessarily empty. And $\left[L m^{2}, L m^{2}+m\right]$ does not necessarily contain a fence. This means that it might be not possible to go from $a_{1}$ to 0 without generating more than $m \bar{v}$ consecutive same colors in observations and satisfying $R^{*}\left(\left(s_{1}\right) m\right)=0$. However, it is clearly possible to go from $a_{1}$ to 0 without generating any big block in observations. So, for $R_{1}$, the description of reverse-path, $R^{*}$ ends: go from $a_{1}$ to 0 without generating any big block in the observations. For example, if $\psi(0)=\psi(1)=\cdots=\psi\left(L m^{3}\right)=$ 1 , then the reverse of the minimum-block path, $R^{*}$, states that $S$ goes to 0 (with suitable many steps, satisfying $R^{*}\left(s_{1} m^{2}\right)=0$ ) by generating only one's. Thus, if $R_{1}$ and $\psi(0)=\psi(1)=\cdots=\psi\left(L m^{3}\right)=1$ hold, then $\psi \circ R_{1}$ starts with at least $m^{3}$ consecutive ones but it does not start with a big block. This means that (2.14) still holds.
Hence, $E_{\text {min }}^{*}(i) \neq \emptyset$ for each $i=1, \ldots, k$.

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[^0]:    * The words $w_{z}^{3}(k)$ provide us some information about (unknown) $z$. This information is captured

