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# A PROBABILISTIC ANALYSIS OF A STRING EDITING PROBLEM AND ITS VARIATION 

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# A PROBABILISTIC ANALYSIS OF A STRING EDITING PROBLEM AND ITS VARIATIONS* 

July 8, 1994

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#### Abstract

We consider a string editing problem in a probabilistic framework. This problem is of considerable interest to many facets of science, most notably molecular biology and computer science. A string editing transforms one string into another by performing a series of weighted edit operations of overall maximum (minimum) cost. The problem is equivalent to finding an optimal path in a weighted grid graph. In this paper, we provide several results regarding a typical behavior of such a path. In particular, we observe that the optimal path (i.e., edit distance) is almost surely (a.s.) equal to $\alpha n$ for large $n$ where $\alpha$ is a constant and $n$ is the sum of lengths of both strings. More importantly, we show that the edit distance is well concentrated around its average value. In the so called independent model in which all weights (in the associated grid graph) are statistically independent, we derive some bounds for the constant $\alpha$. As a by-product of our results, we also present a precise estimate of the number of alignments between two strings. To prove these findings we use techniques of random walks, diffusion limiting processes, generating functions, and the method of bounded difference.


[^0]
## 1. INTRODUCTION

String editing problem arises in many applications, notably in text editing, speech recognition, machine vision and, last but not least, molecular sequence comparison (cf. [36]). Algorithmic aspect of this problem has been studied rather extensively in the past (cf. [2], [30], [32], [33] and [36]). In fact, many important problems on words are special cases of string editing, including the longest common subsequence problem (cf. [1], [14]) and the problem of approximate pattern matching (cf. [12] and [34]).

In sequel we review the string editing problem, its importance, and its relationslip to the longest path problem in a special grid graph.

Let $\mathbf{b}$ be a string consisting of $\ell$ symbols on some alphabet $\Sigma$ of size $V$. There are three operations that can be performed on a string, namely deletion of a symbol, insertion of a symbol, and substitution of one symbol for another symbol in $\Sigma$. With each operation is associated a weight function. We denote by $W_{I}\left(b_{i}\right), W_{D}\left(b_{i}\right)$ and $W_{Q}\left(a_{i}, b_{j}\right)$ the weight of insertion and deletion of the symbol $b_{i} \in \Sigma$, and substitution of $a_{i}$ by $b_{j} \in \Sigma$, respectively. An edit script on $\mathbf{b}$ is any sequence $\omega$ of edit operations, and the total weight of $\omega$ is the sum of weights of the edit operations.

The string editing problem deals with two strings, say $\mathbf{b}$ of length $\ell$ (for $\ell$ ong) and a of length $s$ (for short), and consists of finding an edit script $\omega_{\max }\left(\omega_{\min }\right)$ of minimum (maximum) total weight that transforms a into $\mathbf{b}$. The maximum (minimum) weight is called the edil distance from $\mathbf{a}$ to $\mathbf{b}$, and its is also known as the Levenshtein distance. In molecular biology, the Levenshtein distance is used to measure similarity (homogeneity) of two molecular sequences, say DNA sequences (cf. [33]).

The string edit problem can be solved by the standard dynamic programming method. Let $C_{\max }(i, j)$ denote the maximum weight of transforming the prefix of $\mathbf{b}$ of size $i$ into the prefix of a of size $j$. Then, (cf. [2], [30], [36]).

$$
\begin{aligned}
C_{\max }(i, j)=\max \left\{C_{\max }(i-1, j-1)+W_{Q}\left(a_{i}, b_{j}\right)\right. & , C_{\max }(i-1, j)+W_{D}\left(a_{i}\right) \\
& \left., C_{\max }(i, j-1)+W_{I}\left(b_{j}\right)\right\}
\end{aligned}
$$

for all $1 \leq i \leq \ell$ and $1 \leq j \leq s$. We compute $C_{\max }(i, j)$ row by row to obtain finally the total cost $C_{\max }=C_{\max }(\ell, s)$ of the maximum edit script. A similar procedure works for the minimum edit distance.

The key observation for us is to note that interdependency among the partial optimal weights $C_{\max }(i, j)$ induce an $\ell \times s$ grid-like directed acyclic graph, called further a grid graph. In such a graph vertices are points in the grid and edges go only from ( $i, j$ ) point


Figure 1: Example of a grid graph of size $\ell=3$ and $s=2$.
to neighboring points, namely $(i, j+1),(i+1, j)$ and $(i+1, j+1)$. A horizontal edge from $(i, j-1)$ to $(i, j)$ carries the weight $W_{I}\left(b_{j}\right)$; a vertical edge from $(i, j-1)$ to $(i, j)$ has weight $W_{D}\left(a_{i}\right)$; and finally a diagonal edge from $(i-1, j-1)(i, j)$ is weighted according to $W_{Q}\left(a_{i}, b_{j}\right)$. Figure 1 shows an example of such an edit graph. The edit distance is the longest (shortest) path from the point $O=(0,0)$ to $E=(\ell, s)$.

In this paper, we analyze the string edit problem in a probabilistic framework. We adopt the Bernoulli model for a random string, that is, all symbols of a string are generated independently with probability $p_{i}$ for symbol $i \in \Sigma$. A standard probabilistic model assumes that both strings are generated according to the Bernoulli scheme (c.[. [3], [6], [7], [8], [14], [22], [35], [36]). We call it the string model. Such a framework, however, leads to statistical dependency of weights in the associated grid graph. To avoid this problem, most of the time we shall work within the framework of another probabilistic model which postulates that all weights in the associated grid graph are statistically independent. We call it independent model. This is closely related to a model in which only one string is random, say b, while the other one, say $\mathbf{a}$, is deterministic. Indeed, in such a situation all weights in a "horizontal" strip in the associated grid graph are independent, while weights in a "vertical" strip are dependent (e.g., if $\mathbf{a}=101$, and $\mathbf{b}$ is random, then the " 1 "s in the string a match independently all " 1 "s in b, but clearly the first " 1 " and the third "1" in a have to match " 1 "s in b at the same places). We call such a model semi-independent.

Most of the results in this paper deal either with the independent model or the string
model. We believe that better understanding of the independent model should be the first step to obtain valuable results for the semi-independent model. Certainly, results of the semi-independent model can be further used to deduce probabilistic behavior of the string model ( cf . Theorem 2.2). In passing, we note that the semi-independent model might be useful in some applications (e.g., when comparing a given string to all strings in a data base).

In the independent model the distributions of weights $W_{D}\left(a_{i}\right), W_{l}\left(b_{j}\right)$ and $W_{Q}\left(a_{i}, b_{j}\right)$ depend on the given string a. However, to avoid complicated notations we ignore this fact whenever the independent model is discussed - and consider a grid graph with weights $W_{H}$, $W_{D}$ and $W_{Q}$. In other words, we concentrate on finding the longest path in a grid graph with independent wcights $W_{I}, W_{D}$ and $W_{Q}$, not necessary equally distributed. By selecting properly these distributions, we can model several variations of the string editing problem. For example, in the standard setting the deletion and insertion weights are identical, and usually constant, while the substitution weight takes two values, one (high) when matching between a letter of $\mathbf{a}$ and a letter of $\mathbf{b}$ occurs, and another valuc (low) in the case of a mismatch (e.g., in the Longest Common Substring problem, one sets $W_{I}=W_{D}=0$, and $W_{Q}=1$ when a matching occurs, and $W_{Q}=-\infty$ in the other case).

Our results can be summarized as follows: Applying the Subadditive Ergodic Theorem we note that for the string model and the independent model $C_{\max } \sim \alpha n$ almost surely (a.s.), where $n=\ell+s$ (cf. Theorem 2.1 and Theorem 2.2). Our main contribution lies in establishing bounds for the constant $\alpha$ (cf. Theorem 2.7) for the independent model (cf. Theorem 2.2 for a possible extension to the string model). The upper bound is rather tight as verified by simulation experiments. More importantly, using the powerful and modern method of bounded differences (cf. [29]) we establish for all three models a sharp concentration of $C_{\text {max }}$ around its mean value under a mild condition on the tail of the weight distributions (cf. Theorem 2.3). This proves the conjecture of Chang and Lampe [13] who observed empirically such a sharp concentration of $C_{\text {max }}$ for a version of the string edit problem, namely the approximate string matching problem.

Our probabilistic results are proved in a unified manner by applying techniques of random walks (c. [18], [20]), generating functions (cf. [19], [26], [27]), and bounded differences (cf. [29]). In fact, these techniques allow us to establish further results of a more general interest. In particular, we present an asymptotic estimate for the number of paths in the grid graph (cf. Theorem 2.4), which coincides with the number of sequence alignments (cf. [15], [16], [36]). Finally, for the independent model we establish the limiting distribution of the total weight (cf. Theorem 2.5) and the tail distribution of the total weight (cf. Theorem
2.6) of a randomly selected path (edit script) in the grid graph.

The string edit problem and its special cases (e.g., the longest common subsequence problem and the approximate pattern matching) were studied quite extensively in the past, and are subject of further vigorous research due to their vital application in molecular biology. There are many algorithmic solutions to the problem, and we only mention here Apostolico and Guerra [1], Apostolico el al. [2], Chang and Lampe [13], Myeres [30], Ukkonen [34], and Waterman [36]. On the other hand, a probabilistic analysis of the problem was initiated by Chvatal and Sankoff [14] who analyzed the longest common subsequence problem. After an initial success in obtaining some probabilistic results for this problem, and its extensions by a rather straightforward applications of the subadditive ergodic theorem, a deadlock was reached due to a strong interdependency between weights in the grid graph. To the best of our knowledge, there is no much literature on the probabilistic analysis of the string edit problem and its variations with a notalle exception of a recent marvelous paper by Arratia and Waterman [7] (cf. [35]) who proved their own conjecture concerning phase transitions in a sequence malching.

There is, however, a substantial literature on probabilistic analysis of pattern matching. We mention here a series of papers by Arratia and Waterman (cf. [5], [6]) and with Gordon (cf. [3], [4]), as well as papers by Karlin and his co-authors (cf. [11], [21], [22]). Another approach for the probabilistic analysis of pattern matching with mismatches was recently reported by Atallah et al. in [8].

This paper is organized as follows. In the next section, we present our main results and discuss some of their consequences. Most of our proofs appear in Section 3.

## 2. MAIN RESULTS

We study a grid graph of size $\ell$ and $s(\ell \geq s)$ as shown in Figure 1. All of our results, however, will be expressed in terms of $n=\ell+s$ and $d=\ell-s$. We assign to every edge in such a graph a real number representing its weight. A family of such directed acyclic weighted graphs will be denoted by $\overrightarrow{\mathcal{G}}(n, d)$ or shortly $\overrightarrow{\mathcal{G}}(n)$. We write $\vec{G}(n) \in \overrightarrow{\mathcal{G}}(n, d)$ for a member of such a family.

For the independent model we assume that weights are independent from edge to edge. Let $F_{I}(\cdot), F_{D}(\cdot)$ and $F_{Q}(\cdot)$ denote distribution functions of $W_{I}, W_{D}$ and $W_{Q}$ respectively. Wc assume that the mean values $m_{I}, m_{D}$ and $m_{Q}$, and the variances $s_{I}^{2}, s_{D}^{2}$ and $s_{Q}^{2}$, respectively, are finite. The distribution functions are not necessary identical.

The edit distance can be viewed as an optimization problem on the grid graph. Indeed, let $\mathcal{B}(n, d)$ or shortly $\mathcal{B}(n)$ be the set of all directed paths from the starting point $O$ of the
grid graph to the end point $E$. (It corresponds, as we know, to a script in the original string edit problem.) The cardinality of $\mathcal{B}(n)$, that is, the total number of paths between $O$ and $E$, is denoted by $L(n, d)$. A particular path from $O$ to $E$ is denoted as $\mathcal{P}$, i.e., $\mathcal{P} \in \mathcal{B}(n, d)$. Note that the length $|\mathcal{P}|$ of a path $\mathcal{P}$ satisfies $\ell \leq|\mathcal{P}| \leq l+r=n$. Finally, let $N_{I}(\mathcal{P})$, $N_{D}(\mathcal{P})$ and $N_{Q}(\mathcal{P})$ denote the number of horizontal edges (say $I$-steps), vertical edges (say $D$-steps), and diagonal edges (say $Q$-steps) in a path $\mathcal{P}$.

With the above notation in mind, the problem at hand can be posed as follows:

$$
\begin{equation*}
C_{\max }=\max _{\mathcal{P} \in \mathcal{B}(n)}\left\{W_{n}(\mathcal{P})\right\} \quad, \quad C_{\min }=\min _{\mathcal{P} \in \mathcal{B}(n)}\left\{W_{n}(\mathcal{P})\right\} \tag{1}
\end{equation*}
$$

where $W_{n}(\mathcal{P})$ denotes the total weight of the path $\mathcal{P}$ which becomes

$$
\begin{equation*}
W_{n}(\mathcal{P})=\sum_{i=1}^{N_{I}(\mathcal{P})} W_{I}(i)+\sum_{i=1}^{N_{D}(\mathcal{P})} W_{D}(i)+\sum_{i=1}^{N_{Q}(\mathcal{P})} W_{Q}(i) \tag{2}
\end{equation*}
$$

We write $W_{n}$ to denote the total weight of a randomly selected path, that is,

$$
\begin{equation*}
\operatorname{Pr}\left\{W_{n}<x\right\}=\frac{1}{L(n, d)} \sum_{\mathcal{P} \in \mathcal{B}} \operatorname{Pr}\left\{W_{n}(\mathcal{P})<x\right\} . \tag{3}
\end{equation*}
$$

Our results crucially depend on the order of magnitude of $d$ with respect to $n$. We consider separately several cases. Below we define two of them that are analyzed in details in this paper:
$\operatorname{CASE}(\mathrm{A}): \quad d=O(\sqrt{n})$, and let $x=d \sqrt{\sqrt{2} / n}=\zeta d / \sqrt{n}$ where $\zeta=2^{1 / 4}$.
$\mathrm{C} \Lambda \mathrm{SE}(\mathrm{B}): \quad d=\Theta(n)$, and let $x=d / n$.
Three other case, described below, are discussed in our extended technical report [28]:
$\operatorname{CASE}(\mathrm{C}): \quad d=n-O\left(n^{1-\varepsilon}\right)$, that is, for some constant $x$ we have $d=n\left(1-x / n^{\varepsilon}\right)$.
$\operatorname{CASE}(\mathrm{D}): d=O(1)$ (we shall reduce this case to Case (A)).
CASE (E): $s=O(1)$ (we shall reduce this case to case (C)).
Now, we are in a position to present our results. To simplify further our presentation, we concentrate mainly on the longest path $C_{\text {max }}$. We start with a simple general result concerning the typical behavior of $C_{\max }$. The more refined results containing a computable upper bound for $E C_{\max }$ (in the independent model) are given at the end of this section (cf. Theorem 2.7).

Theorem 2.1. In the string model and the independent model, the following holds

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \frac{C_{\max }^{\prime}}{n}=\lim _{n \rightarrow \infty} \frac{E C_{\max }}{n}=\alpha \quad \text { (a.s.) } \tag{4}
\end{equation*}
$$

provided $\ell / s$ has a limit as $n \rightarrow \infty$.
Proof. Let us consider the $\ell \times s$ grid with starting point $O$ and ending point $E$ (cf. Fig. 1). Call it $\operatorname{Grid}(O, E)$. We also choose an arbitrary point, say $A$, inside the grid so that we can consider two grids, namely $\operatorname{Grid}(O, A)$ and $\operatorname{Grid}(A, E)$. Actuadly, point $A$ splits the edit distance problem into two subproblems with objective functions $C_{\max }(O, A)$ and $C_{\max }(A, E)$. Clearly, $C_{\max }(O, E) \geq C_{\max }(O, A)+C_{\max }(A, E)$. Thus, under our assumption regarding weights, the objective function $C_{\text {max }}$ is superadditive, and direct application of Superadditive Ergodic Theorem (cf. [25]) proves our result.

Remark 1. Observe that we cannot directly apply subadditive ergodic theorem to semiindependent model since weights are not stationary in this case. However, using an inductive argument, one can obtain similar results as above for the semi-independent model. In particular, since for the semi-independent model $E C_{\text {max }}$ is superadditive, we immediately prove that $E C_{\text {max }} \sim \alpha \ell$.

In the string and semi-independent models, weights depend on strings a and $\mathbf{b}$, hence the constant $\alpha$ is a function of $\mathbf{a}$ and $\mathbf{b}$. Furthermore, the string model can be reduced to the semi-independent model as follows. Let a be a given string (i.e., not random), and let $P(a)$ be the probability of a occurrence in our standard Bernoulli model (e.g., for the binary alphabet $\Sigma=\{a, b\}$ we have $P(a)=p^{|a|}(1-p)^{|b|}$ where $p$ is the probability of $a$ occurrence, and $|a|(|b|)$ is the number of $a$ 's ( $b$ 's) in the string a). Let $\alpha_{a}$ be the constant in the semi-independent model.

Theorem 2.2. In the string model, the constant $\alpha$ can be estimated as follows

$$
\begin{equation*}
\alpha=\sum_{\mathbf{a} \in \mathcal{H}} \alpha_{\mathbf{a}} P(\mathbf{a}) \tag{5}
\end{equation*}
$$

where $\mathcal{H}$ is the set of all possible strings a of length $s$ over the alphabet $\Sigma$.
Proof. Observe the following

$$
\sum_{\mathbf{a} \in \mathcal{H}} \alpha_{\mathbf{a}} P(\mathbf{a})=E \alpha_{\mathbf{a}}=E\left(\lim _{\ell \rightarrow \infty} \frac{E C_{\mathbf{n a x}}}{\ell} \cdot \frac{\ell}{n}\right)=\lim _{\ell \rightarrow \infty} \frac{E\left(E C_{\max }\right)}{\ell}=\alpha
$$

where the first equality is just definition of the expected value, the second follows from Remark $!$, while the last is a simple consequence of the bounded convergence theorem and $\ell / n \rightarrow 1$.

Finally, for the string and independent models we can report the following finding concerning the concentration of the edit distance. It proves the conjecture of Chang and Lampe
[13]. The proo[ of this result uses a powerful method of bounded differences or Azuma's type inequality (cr. [29]).

Theorem 2.3. (i) If all weights arc bounded random variables, say $\max \left\{W_{I}, W_{D}, W_{Q}\right\} \leq 1$, then for arbitrary $\varepsilon>0$ and large $n$

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|C_{\text {max }}-E C_{\max }\right|>\varepsilon E C_{\text {max }}\right\} \leq 2 \exp \left(-\varepsilon^{2} \alpha n\right) . \tag{6}
\end{equation*}
$$

(ii) If the weights are unbounded but such that for large $n, W_{\max }=\max \left\{W_{l}, W_{D}, W_{Q}\right\}$ satisfies the following

$$
\begin{equation*}
n \operatorname{Pr}\left\{W_{\max } \geq n^{1 / 2-\delta}\right\} \leq U(n) \tag{7}
\end{equation*}
$$

for some $\delta>0$ and a function $U(n) \rightarrow 0$ as $n \rightarrow \infty$, then

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|C_{\max }-E C_{\max }^{\prime}\right|>\varepsilon E C_{\text {max }}\right\} \leq 2 \exp \left(-\beta n^{\delta}\right)+U(n) \tag{8}
\end{equation*}
$$

for any $\varepsilon>0$ and some $\beta>0$.
Proof: We consider only the string model. Part (i) is a direct consequence of the following inequality of Azuma's lype (cf. [29]): Let $X_{i}$ be i.i.d. random variables such thal for some function $\int(-, \ldots, \cdot)$ the following is true

$$
\begin{equation*}
\left|f\left(X_{1}, \ldots, X_{i}, \ldots, X_{n}\right)-f\left(X_{1}, \ldots, X_{i}^{\prime}, \ldots, X_{n}\right)\right| \leq c_{i} \tag{9}
\end{equation*}
$$

where $c_{i}<\infty$ are constants, and $X_{i}^{\prime}$ has the same distribution as $X_{i}$. Then,

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|f\left(X_{1}, \ldots, X_{n}\right)-E f\left(X_{1}, \ldots, X_{n}\right)\right| \geq t\right\} \leq 2 \exp \left(-2 t^{2} / \sum_{i=1}^{n} c_{i}^{2}\right) \tag{10}
\end{equation*}
$$

for some $t>0$. The above technique is also called the method of bounded differences.
Now, for part (i) it suffices to set $X_{i}=b_{i}$ for $1 \leq i \leq \ell$, and $X_{i}=a_{i-\ell}$ for $\ell+1 \leq i \leq n$, where $a_{i}$ and $b_{i}$ are the $i$ symbols of the two strings a and $b$. Under our Bernoulli model, the $X_{i}$ are i.i.d. and (9) holds, with $f(\cdot)=C_{\text {max }}^{\prime}$. More precisely,

$$
\begin{equation*}
\left|C_{\max }\left(X_{1}, \ldots, X_{i}, \ldots, X_{n}\right)-C_{\max }\left(X_{1}, \ldots, X_{i}^{\prime}, \ldots, X_{n}\right)\right| \leq \max _{1 \leq i \leq n}\left\{W_{\max }(i)\right\} \tag{11}
\end{equation*}
$$

where $W_{\max }(i)$ is the $i$ th independent version of $W_{\max }$ defined in the theorem. Clearly, for part (i) we have $c_{i}=1$, thus we can apply (10). Inequality (6) follows from the above and $t=\varepsilon E C_{\text {max }}=O(n)$.

To prove part (ii), we start with (11). But, this time we have for some $c$

$$
\begin{aligned}
\operatorname{Pr}\left\{\left|C_{\max }^{\prime}-E C_{\text {max }}\right| \geq t\right\} & =\operatorname{Pr}\left\{\left|C_{\max }-E C_{\max }\right| \geq t, \max _{1 \leq i \leq n}\left\{W_{\max }(i)\right\} \leq c\right\} \\
& +\operatorname{Pr}\left\{\left|C_{\max }-E C_{\max }\right| \geq t, \max _{1 \leq i \leq n}\left\{W_{\max }(i)\right\}>c\right\} \\
& \leq 2 \exp \left(-2 t^{2} / n c^{2}\right)+n \operatorname{Pr}\left\{W_{\max }>c\right\} .
\end{aligned}
$$

Set now $t=\varepsilon E C_{\text {max }}=O(n)$ and $c=O\left(n^{1 / 2-\delta}\right)$, then

$$
\operatorname{Pr}\left\{\left|C_{\max }-E C_{\max }\right| \geq \varepsilon E C_{\max }\right\} \leq 2 \exp \left(-\beta n^{\delta}\right)+n \operatorname{Pr}\left\{W_{\max }>n^{1 / 2-\delta}\right\},
$$

for some constant $\beta>0$, and this implies (8) provided (7) holds.
Remark 2. Theorem 2.3 holds also for the semi-independent model if one replaces in the right-hand side of (6) $2 \exp \left(-\varepsilon^{2} n\right)$ by $2 \exp \left(-\varepsilon^{2} \ell\right)$ and set $E C_{\text {max }}=O(\ell)$.

Hereafter, we investigate only the independent model. For this model, we have obtained several new results regarding the probabilistic behavior of (longest) path in a weighted grid graph.

The next result presents limiting distribution of the total weight defined in (2). Its proof is quite complicated, however, it applies only standard techniques. Therefore, at the referee request, we omit completely the proof of this thcorem. It can be found in our extended technical report [28].

Theorem 2.4. The limiting distribution of the total weight satisfies

$$
\begin{equation*}
\frac{W_{n}-n \mu_{W}}{\sqrt{n} \sigma_{W}} \rightarrow \mathcal{N}(0, \mathrm{I}) \tag{12}
\end{equation*}
$$

where $\mathcal{N}(0,1)$ is the standard normal distribution, and

$$
\begin{align*}
\mu_{W} & =m_{I} \mu_{I}+m_{D} \mu_{D}+m_{Q} \mu_{Q}  \tag{13}\\
\sigma_{W}^{2} & =\mu_{I} s_{I}^{2}+\mu_{D} s_{D}^{2}+\mu_{Q} s_{Q}^{2}+\tilde{\sigma}_{Q}^{2}\left(m_{I}+m_{D}-m_{Q}\right)^{2} \tag{14}
\end{align*}
$$

where $\mu_{I}=E N_{I}(\mathcal{P}), \mu_{D}=E N_{D}(\mathcal{P}), \mu_{Q}=E N_{Q}(\mathcal{P})$ and $\tilde{\sigma}_{Q}^{2}=\operatorname{var} N_{Q}(\mathcal{P})$. Explicit formulas for these quantities can be found in the next section.

Our next resull enumerates the number of paths $L(n, d)$ in the grid graph. It is also of interest to some other problems since $L(n, d)$ represents the number of ways the string a can be transformed into $\mathbf{b}$, and this problem was already tackled by others (cr. [15], [16], [24], [36]) in the case of equal length strings (i.e., $\ell=s$ ). The formulation of this result depends on a parameter $u=d / n$ that takes different values for case (A) and (B), that is:
$\operatorname{CASE}(\mathrm{A}): \quad$ Set $d=x \sqrt{n / \sqrt{2}}$. Then, $u=x / \sqrt{\sqrt{2} n}=x /(\varsigma \sqrt{n})$.
CASE (B): Set $u=x$.
Theorem 2.5. Let $L(u)=L(n, d)$ be the number of paths in a grid graph $\vec{G} \in \overrightarrow{\mathcal{G}}(n)$. Then,

$$
\begin{equation*}
L(u)=\frac{C \psi_{2}\left(\beta_{2}(u)\right)^{n}}{\beta_{2}(u)^{n(1+u) / 2} \sqrt{2 \pi n V(u)}}(1+O(1 / n)) \tag{15}
\end{equation*}
$$

where

$$
\begin{align*}
& \beta_{2}(u)=\frac{1+3 u^{2}+u \sqrt{8\left(u^{2}+1\right)}}{1-u^{2}}  \tag{16}\\
& \psi_{2}(u)=\psi_{2}\left[\beta_{2}(u)\right]=\frac{2 u \beta_{2}(u)}{\beta_{2}(u)-1-u\left(1+\beta_{2}(u)\right)} \tag{17}
\end{align*}
$$

and $C$ is a constant that is found in Section 3 (cf. (79)). In the above, $V(u)$ is the variance obtained from the gencrating function $h(z)$ defined as $h(z)=\psi_{2}\left(z \beta_{2}(u)\right) / \psi_{2}\left(\beta_{2}(u)\right)$, that is, $V(u)=h^{\prime \prime}(1)-0.25\left(1-u^{2}\right)$ where $h^{\prime \prime}(z)$ is the second derivative of $h(z)$.

For most of our computations, we only necd the asymptotics of $L(u)$ in the following a less precise form

$$
\begin{equation*}
\log L(u)=n \rho(u)-0.5 \log n+O(1) \tag{18}
\end{equation*}
$$

where $\rho(u)$, for cases (A), (B) is respectively

$$
\begin{align*}
& \rho(u)=-\log (\sqrt{2}-1)  \tag{19}\\
& \rho(u)=\log \psi_{2}\left(\beta_{2}(u)\right)-\frac{1+u}{2} \log \beta_{2}(u) \tag{20}
\end{align*}
$$

The details of the above derivations can be found in Section 3.
Finally, in order to obtain an upper bound for the cost $C_{\text {max }}$, we need an estimate on the tail distribution of the total weight $W_{n}$ along a random path. Formula (2) suggests to apply Cramer's large deviation result (cf. Feller [18]) with some modifications (due to the fact that the total weight $W_{n}$ as in (2) is a sum of random number of weights). To avoid unnecessary complications, we consider in details only two cases, namely:
(a) all weights are identically distributed with mean $m=m_{I}=m_{D}=m_{Q}$ and the cumulant function $\Psi(s)=\log E e^{s(W-m)}$ for the common weight $W-m$;
(b) insertion weight and deletion weight are constant, say all equal to -1 (e.g., $W_{I}=$ $W_{D}=-1$ ), and the substitution weight $W_{Q}-m_{Q}$ has the cumulant function $\Psi_{Q}(s)=$ $\log E e^{s\left(W_{Q}-m_{Q}\right)}$. Such an assignment of weights is often encountered in the string editing problem.

Theorem 2.6. (i) In the case (a) of all identical weights, define $s^{*}$ as the solution of

$$
\begin{equation*}
a=\Psi^{\prime}\left(s^{*}\right), \tag{21}
\end{equation*}
$$

for a given $a>0$, and let

$$
\begin{align*}
& Z_{0}(a)=s^{*} \Psi^{\prime}\left(s^{*}\right)-\Psi\left(s^{*}\right)  \tag{22}\\
& E_{1}(a)=-\left(s^{*} m+\Psi\left(s^{*}\right)\right)  \tag{23}\\
& E_{2}^{2}(a)=\frac{\tilde{\sigma}_{Q}^{2} m^{2}+2 \widetilde{\sigma}_{Q}^{2} m a+\tilde{\sigma}_{Q}^{2}\left(\Psi^{\prime}\left(s^{*}\right)\right)^{2}+\left(1-\mu_{Q}\right) \Psi^{\prime \prime}\left(s^{*}\right)}{2\left(1-\mu_{Q}\right) \tilde{\sigma}_{Q}^{2} \Psi^{\prime \prime \prime}\left(s^{*}\right)} \tag{24}
\end{align*}
$$

where $\mu_{Q}=E N_{Q}$ and $\tilde{\sigma}_{Q}^{2}=\operatorname{var} N_{Q}$. Then,

$$
\begin{align*}
\operatorname{Pr}\left\{W_{n}>\right. & \left.\left(1-\mu_{Q}\right)(a+m) n\right\} \sim \\
& \frac{1}{2 s^{*} E_{2}(a) \tilde{\sigma}_{Q} \sqrt{\pi\left(1-\mu_{Q}\right) n \Psi^{\prime \prime}\left(s^{*}\right)}} \exp \left(-n\left(1-\mu_{Q}\right) Z_{0}(a)+n \frac{E_{1}^{2}(a)}{4 E_{2}^{2}(a)}\right) \tag{25}
\end{align*}
$$

(ii) In case (b) of constant I-weights and D-weights, we define $s^{*}$ as a solution of

$$
\begin{equation*}
a=\Psi_{Q}^{\prime}\left(s^{*}\right), \tag{26}
\end{equation*}
$$

and let

$$
\begin{align*}
& Z_{0}(a)=s^{*} \Psi_{Q}^{\prime}\left(s^{*}\right)-\Psi_{Q}\left(s^{*}\right),  \tag{27}\\
& E_{1}(a)=s^{*}\left(m_{Q}+2\right)+2 s^{*} a^{*}-\Psi\left(s^{*}\right),  \tag{28}\\
& E_{2}^{2}(a)=\frac{\tilde{\sigma}_{Q}^{2}\left(m_{Q}+2\right)\left(m_{Q}+2+2 a^{*}+4 s^{*} \Psi_{Q}^{\prime \prime}\left(s^{*}\right)\right)+\tilde{\sigma}_{Q}^{2} a^{*}\left(a^{*}+4 s^{*} \Psi_{Q}^{\prime \prime}\left(s^{*}\right)\right)+\mu_{Q} \Psi_{Q}^{\prime \prime}\left(s_{(2)}^{*}\right)}{2 \mu_{Q} \tilde{\sigma}_{Q}^{2} \Psi_{Q}^{\prime \prime}\left(s^{*}\right)}
\end{align*}
$$

Then,

$$
\begin{align*}
\operatorname{Pr}\left\{W_{n}>\right. & \left.\mu_{Q}\left(a+\beta / \mu_{Q}\right) n\right\} \sim \\
& \frac{1}{2 s^{*} E_{2}(a) \tilde{\sigma}_{Q} \sqrt{\pi \mu_{Q} n \Psi_{Q}^{\prime \prime}\left(s^{*}\right)}} \exp \left(-n \mu_{Q} Z_{0}(a)+n \frac{E_{1}^{2}(a)}{4 E_{2}^{2}(a)}\right) . \tag{30}
\end{align*}
$$

where $\beta=2 \mu_{Q}+m \mu_{Q}-1$.
Having the above estimates on the tail of the total cost of a path in the grid graph $\vec{G} \in \overrightarrow{\mathcal{G}}(n)$, we can provide a more precise information about the constant $\alpha$ in our Theorem 2.1, that is, we compute an upper bound $\bar{\alpha}$ and a lower bound $\underline{\alpha}$ of $\alpha$ for the independent model. We prove below the following result, which is one of our main finding.

Theorem 2.7 Assume the independent model.
(i) Consider first the identical weights case (a) above. Let $a^{*}$ be a solution of the following equation

$$
\begin{equation*}
\left(1-\mu_{Q}\right) Z_{0}\left(a^{m}\right)=\rho+\frac{E_{1}^{2}\left(a^{*}\right)}{4 E_{2}^{2}\left(a^{*}\right)}, \tag{31}
\end{equation*}
$$

where $\rho$ is defined in (19)-(20), and $Z_{0}, E_{1}$ and $E_{2}^{2}$ are defined in (22)-(24). Then, the upper bound $\bar{\alpha}$ of $\alpha$ becomes

$$
\begin{equation*}
\bar{\alpha}=\left(1-\mu_{Q}\right)\left(a^{*}+m\right)+O(\log n / n) \tag{32}
\end{equation*}
$$

In the case (b) of constant I and D weights, let $a^{*}$ be a solution of the equation

$$
\begin{equation*}
\mu_{Q} Z_{0}\left(a^{*}\right)=\rho+\frac{E_{1}^{2}\left(a^{*}\right)}{4 E_{2}^{2}\left(a^{m}\right)}, \tag{33}
\end{equation*}
$$

where $Z_{0}, E_{1}$ and $E_{2}^{2}$ are as in (27-30). Then,

$$
\begin{equation*}
\bar{\alpha}=\mu_{Q}\left(a^{*}+\beta / \mu_{Q}\right)+O(\log n / n), \tag{34}
\end{equation*}
$$

where $\beta$ is defined in Theorem 2.6(ii).
(ii) The lower bound $\underline{\alpha}$ of $\alpha$ can be obtained from a particular solution to our optimization problem (1). In particular, we have

$$
\begin{equation*}
\underline{\alpha}=\max \left\{\mu_{W}, \ell m_{D}+s m_{I}, \alpha_{g r}\right\}, \tag{35}
\end{equation*}
$$

where $\alpha_{g r}$ is constructed from a greedy solution of the problem, that is,

$$
\begin{equation*}
n \alpha_{g r}=(\ell+s(1-p)) m_{m a x} \tag{36}
\end{equation*}
$$

where $p=\operatorname{Pr}\left\{W_{Q}>W_{I}\right.$ and $\left.W_{Q}>W_{D}\right\}$, and $m_{\max }=E \max \left\{W_{I}, W_{D}, W_{Q}\right\}$.
Proof. We first prove part (i) provided Theorem 2.6 is granted (cf. Section 3 for the proof). Observe that by Boole's inequality we have for any real $x$

$$
\operatorname{Pr}\left\{C_{\text {max }}>x\right\} \leq \sum_{\mathcal{P} \in \mathcal{B}} \operatorname{Pr}\left\{W_{n}(\mathcal{P})>x\right\}=L(u) \operatorname{Pr}\left\{W_{n}>x\right\}
$$

where the last equality follows from (3). We now consider only case (a). Let $\beta(a)=$ $\left(1-\mu_{Q}\right) Z_{0}(a)$ and $\gamma(a)=E_{1}^{2}(a) /\left(4 E_{2}^{2}(a)\right)$. Then, by Theorem 2.5 and $2.6(\mathrm{i})$ we have

$$
\operatorname{Pr}\left\{C_{\text {max }}>\left(1-\mu_{Q}\right)(a+m) n\right\} \leq O(1 / n) \exp (n(\rho+\gamma(a)-\beta(a)))
$$

Setting in the above $a=a^{*}$ as defined in (31), we prove our result.
The lower bound can be established either by considering some particular paths $\mathcal{P}$ or applying a simple algorithm like a greedy one. The greedy algorithm selects in every step the most expensive edge, that is, the average cost per step is $m_{\text {max }}=E \max \left\{W_{D}, W_{I}, W_{Q}\right\}$. Let $p=\operatorname{Pr}\left\{W_{Q}>W_{I}, W_{Q}>W_{D}\right\}$. Observe that if there are $k D$-stcps, then necessarily, there are $s-k Q$-steps. But, the number of $Q$-steps is binomially distributed with parameters $p$ and $s$. Thus,

$$
n \alpha_{g r}=m_{\max } \sum_{k=1}^{s}(\ell+k)\binom{s}{k} p^{s-k}(1-p)^{k}=\ell+s(1-p) n_{\max }
$$

Table 1: Simulation results for exponentially distributed weights with means $m_{I}=m_{D}=$ $m_{Q}=1$ for case (B) with $d=0.6 n$.

| $\ell$ | $s$ | $\underline{\alpha}$ | $\alpha_{\operatorname{sim}}$ | $\bar{\alpha}$ |
| ---: | ---: | ---: | ---: | ---: |
| 200 | 50 | 1.588 | 1.909 | 2.45 |
| 400 | 100 | 1.588 | 1.808 | 2.45 |
| 600 | 150 | 1.588 | 1.899 | 2.45 |
| 800 | 200 | 1.588 | 1.926 | 2.45 |
| 1000 | 250 | 1.588 | 1.922 | 2.45 |

and this proves our result.
We compared our bounds for $C_{\max }$ with some simulation experiments. In the simulation we restricted our analysis to uniformly and exponentially distributed weights, and here we only report the latter results. They are shown in Table 1. It is plausible that the normalized limiting distribution for $C_{\text {max }}$ is double exponential (i.e., $e^{-c^{-x}}$ ), however, the normalizing constants are quit hard to find.

The editing problem can be generalized, as it was recently done by Pevzner and Waterman [32] for the longest common subsecfuence problem. In terms of the grid graph, their generalization boils down to adding new edges in the grid graph that connect no-neighboring vertices. In sucl a situation our Theorem 2.1 may not hold. In fact, based on recent results of Newman [31] concerning the longest (unweighted) path in a general acyclic graph, we predict that a phase transition can occur, and $C_{\text {max }}$ may switch from $\Theta(n)$ to $\Theta(\log n)$. This was already observed by Arratia and Waterman [7] for another string problen, namely, for the score in the pattern matching problem.

## 3. ANALYSIS THROUGH THE RANDOM WALK APPROACH

In tlis section, we only analyze the independent model. To recall, we consider an $\ell \times s$ grid graph with independent weights $W_{I}, W_{D}$ and $W_{Q}$. We represent a path in the grid graph $\vec{G}$ as a random walk. First of all, it is convenient to append our $\ell \times s$ graph to a full $\ell \times \ell$ grid graph, with all steps possible, as shown in Figure 2. It should be noted that in our new representation, a $Q$-step is twice as long as $I$-step and $D$-step, and therefore the increments in such a random walk are not independent (e.g., after the first diagonal move, the second one comes with probability one).


Figure 2: An extended $\ell \times \ell$ grid graph
We first analyze a path without weights in the grid graph shown in Figure 2. We call it an unweighted random walk (in short: R.W.) and denote as $Y(\cdot)$. To model a path $\mathcal{P}$ in our original problem, we must assure that the random walk $Y(\cdot)$ coincides with the script path $\mathcal{P}$, we require that the random walk $Y(\cdot)$ in Figure 2 ends at the point $E$ of the grid graph after $n$ steps where $n=2 \ell-(\ell-s)=\ell+s$. Thus, we impose the following constraint

$$
\begin{equation*}
Y(n)=d \tag{37}
\end{equation*}
$$

where $d=\ell-s$.
We first consider an unconstraint random walk $\hat{Y}(\cdot)$ such that the condition (37) does not hold, and that the probabililies of $I$-step, $D$-step and $Q$-step $\tau=\sqrt{2}-1, \tau$ and $\tau^{2}$ respectively, as shown in Figure 3. These probabilities are chosen in such a way that all paths with the same length reccive the same probability (e.g., a two-step path $I \& D$ has probability $r^{2}$, the same as one-step path $Q$ of length two).

### 3.1 Case (A): $d=O(\sqrt{n})$

Consider first the unconstraint random walk $\hat{Y}(\cdot)$ (cl. Fig. 3). We make the following scale changes $\ell=\frac{i}{n}, y=j \frac{\zeta}{\sqrt{n}}$ with $\zeta=\sqrt{\sqrt{2}}$ to establish the following theorem, where $\underset{n \rightarrow \infty}{\Rightarrow}$ represents the weak convergence of random functions in the space of all right continuous functions having left limit and endowed with the Skorohod metric (see Billingsley [9] Ch.III).


Figure 3: Probabilities of $I$-step, $D$-step and $Q$-step in the unconstraint random walk $\dot{Y}$
Theorem 3.1. The unconstraint R.W. $\hat{Y}(\cdot)$ possesses the following limiting behavior

$$
\frac{\zeta \hat{Y}([n t])}{\sqrt{n}} \Rightarrow B(l), \quad n \rightarrow \infty
$$

where $B(\cdot)$ is a classical Brownian Motion (B.M.), and $\zeta=\sqrt{\sqrt{2}}$.
Proof. Let $p_{i}(j)=\operatorname{Pr}\{\hat{Y}(i)=j\}$. Then, $p_{i+1}(j)=\tau p_{i}(j-1)+\tau p_{i}(j+1)+\tau^{2} p_{i-1}(j)$ for $i \geq 2$, and our result lollows from standard arguments (cf. [28]).

Now, we take into account the constraint (37), that is, we set $Y(n)=d=O(\sqrt{n})$. Let

$$
\begin{equation*}
x=d \frac{\zeta}{\sqrt{n}} \tag{38}
\end{equation*}
$$

with $\zeta=2^{1 / 4}$. To handle this constraint, we recompute the probabilities of the $I, D$, and $Q$ steps so that $E Y(n)=d$ lolds, and later we relax it so that our primary constraint (37) is true (cf. (48)). We define these new one-step probabilities as follows: $p_{I}=\operatorname{Pr}\{$ first move is $I \mid Y(n)=d\}, p_{D}=\operatorname{Pr}\{$ first move is $D \mid Y(n)=d\}$, and $p_{Q}=$ $\operatorname{Pr}\{$ first move is $Q \mid Y(n)=d\}$. Note that these probabilities depend on $n$, but we do not show explicitly this dependency.

Lemma 3.2a. The new one-step probabilities become

$$
\begin{align*}
& p_{I}=\tau\left(1+\frac{\zeta}{\sqrt{n}} x\right)+O\left(\frac{1}{n}\right),  \tag{39}\\
& p_{D}=\tau\left(1-\frac{\zeta}{\sqrt{n}} x\right)+O\left(\frac{1}{n}\right),  \tag{40}\\
& p_{Q}=\tau^{2}+O\left(\frac{1}{n}\right) \tag{41}
\end{align*}
$$

Proof. We know from Theorem 3.1 that $f(t, x)=\frac{e^{-\frac{x^{2}}{2 t}}}{\sqrt{2 \pi t}}$. This, and the above definitions of $p_{I}, p_{D}$ and $p_{Q}$ lead to the following

$$
p_{I}=\tau \frac{p_{n-1}(d-1)}{p_{n}(d)}=\frac{\tau}{p_{n}(d)}\left(p_{n}(d)-\frac{1}{n} \partial_{t} f-\frac{\zeta}{\sqrt{n}} \partial_{x} f+\frac{\zeta^{2}}{2 n} \frac{\partial^{2} f}{\partial x^{2}}+O\left(\frac{1}{n^{3 / 2}}\right)\right)
$$

But $\partial_{x} f=-\frac{x}{t} f$, hence $p_{I}$ is now readily computed by setting $t=1$ in the above. The two other probabilities are derived in a similar manner.

Remark 3. In fact, using similar arguments to the ones in the proof of Theorem 3.1, we can prove much stronger result. Namely, the constraint random walk $Y(\cdot)$ characterized by the probabilities $p_{I}, p_{D}, p_{Q}$ las the limiting density given by $f(y, v)=\exp \left(-\frac{(y-x)^{2}}{2 v}\right) / \sqrt{2 \pi} v$, which is exactly the density of a B.M. with drift $x$ and variance $v$.

To estimate the large deviation of the total weight $W_{n}$ we need a precise evaluation of the random variables $N_{I}, N_{D}$ and $N_{Q}$ representing the number of $I$-steps, $D$-steps and $Q$-steps in a path $\mathcal{P}$. First of all, we compute the limiting distribution of the sum $N_{I}+N_{D}+N_{Q}$. Using the renewal theory (cf. Feller [18], p. 321, 341, and Iglehart [20] Theorem 4.1) we can easily prove that

$$
\begin{equation*}
N_{I}+N_{D}+N_{Q} \sim \mathcal{N}\left(\frac{n}{\bar{d}}, n \frac{\bar{\sigma}^{2}}{\dot{d}^{3}}\right)+O(1), \quad n \rightarrow \infty \tag{42}
\end{equation*}
$$

where $\mathcal{N}\left(m, \sigma^{2}\right)$ is a classical Gaussian variable with mean $m$ and variance (VAR) $\sigma^{2}$. In the above, $\bar{d}$ is the average move step, that is, from Lemma 3.2 we have

$$
\begin{equation*}
\bar{d}=p_{I}+p_{D}+2 p_{Q}=1+p_{Q} \tag{43}
\end{equation*}
$$

so that $\bar{d}=2(2-\sqrt{2})+O(1 / n)$, and

$$
\begin{equation*}
\bar{\sigma}^{2}=p_{Q}\left(1-p_{Q}\right) \tag{44}
\end{equation*}
$$

hence $\bar{\sigma}^{2}=\sqrt{2}(10-7 \sqrt{2})+O(1 / n)$. Let

$$
\begin{equation*}
\alpha=\frac{1}{\bar{d}}=\frac{1}{1+p_{Q}}=\frac{2+\sqrt{2}}{4}+O\left(\frac{1}{n}\right) \tag{45}
\end{equation*}
$$

and

$$
\begin{equation*}
\kappa=\frac{\bar{\sigma}^{2}}{\bar{d}^{3}}=\frac{\sqrt{2}}{16}+O\left(\frac{1}{n}\right) \tag{46}
\end{equation*}
$$

Then, from (42), we obtain $N_{I}+N_{D}+N_{Q} \sim \mathcal{N}(n \alpha, n \kappa)+O(1)$.
From the expression (2) on the total weight $W_{n}$, it should be clear that we need the joint distribution of $N_{I}, N_{D}$ and $N_{Q}$ (cf. Louchard et al. [27]). For this, we must consider two
constraints on $N .:^{1}$ one on the total number of steps, and the other related to $Y(n)=d$. More precisely, together with (38) we have the following constraint on the number of steps

$$
\begin{align*}
N_{I}+N_{D}+2 N_{Q} & =n  \tag{47}\\
N_{I}-N_{D} & =d=x \frac{\sqrt{n}}{\zeta} \tag{48}
\end{align*}
$$

We first consider only the constraint (47). This will allow us to compute the asymptotic joint distribution of $N_{I}, N_{D}, N_{Q}$, as stated in the next theorem. The proof can be found in Appendix A.

Theorem 3.3a. The number of $I, D$ and $Q$ steps, $N_{I}, N_{D}, N_{Q}$ are asymptotically Gaussian, with mean $n \mu_{I}, n \mu_{D}, n \mu_{Q}$ respectively, where

$$
\begin{align*}
& \mu_{I}=\bar{p}_{I}(2 \alpha-1)=p_{I} /\left(1+p_{Q}\right)=\frac{\sqrt{2}}{4}\left(1+\frac{\zeta}{\sqrt{n}} x+O\left(\frac{1}{n}\right)\right),  \tag{49}\\
& \mu_{D}=\bar{p}_{D}(2 \alpha-1)=p_{D} /\left(1+p_{Q}\right)=\frac{\sqrt{2}}{4}\left(1-\frac{\zeta}{\sqrt{n}} x+O\left(\frac{1}{n}\right)\right),  \tag{50}\\
& \mu_{Q}=(1-\alpha)=p_{Q} /\left(1+p_{Q}\right)=\frac{2-\sqrt{2}}{4}+O\left(\frac{1}{n}\right) \tag{51}
\end{align*}
$$

where $\bar{p}_{I}=p_{I} /\left(p_{I}+p_{D}\right), \bar{p}_{D}=p_{D} /\left(p_{I}+p_{D}\right)$. Moreover, $\alpha$ is given by ( 15$)$. The asymptolic covariance matrix is given by

$$
n \cdot \begin{gather*}
I  \tag{52}\\
D \\
Q
\end{gather*}\left(\begin{array}{ccc}
\sigma_{I}^{2} & C_{I D} & -2 \bar{p}_{I} \kappa \\
C_{I D} & \sigma_{D}^{2} & -2 \bar{p}_{D} \kappa \\
C_{I Q} & -2 \bar{p}_{D \kappa} & \kappa
\end{array}\right)
$$

with

$$
\begin{aligned}
\sigma_{I}^{2} & =(2 \alpha-1) \bar{p}_{I}\left(1-\bar{p}_{I}\right)+4 \kappa \bar{p}_{I}^{2}=\frac{3 \sqrt{2}}{16}+\frac{2^{3 / 4}}{8} x+O\left(\frac{1}{n}\right) \\
\sigma_{D}^{2} & =(2 \alpha-1) \bar{p}_{D}\left(1-\bar{p}_{D}\right)+4 \kappa \bar{p}_{D}^{2}=\frac{3 \sqrt{2}}{16}-\frac{2^{3 / 4}}{8} x+O\left(\frac{1}{n}\right) \\
C_{I D} & =2 \kappa-(2 \alpha-1) \bar{p}_{I} \bar{p}_{D}-2 \kappa\left(\bar{p}_{I}^{2}+\bar{p}_{D}^{2}\right)=\frac{\sqrt{2}}{16}+O\left(\frac{1}{n}\right)
\end{aligned}
$$

where $\kappa$ is given in (46).
To complete our study of the number of steps in the grid graph, we must take into account the constraint (48). Set $\eta .=(N .-n \mu.) / \sqrt{n}$. Observe that by Lemma 3.2 a and Theorem 3.3a, $E\left(N_{I}-N_{D}\right)=\frac{\sqrt{n}}{6} x+O(1)$ as it sloould be, so (48) and (47) imply respectively that $\eta_{I}=\eta_{D}$ and $\eta_{I}=-\eta_{Q}$.

[^1]To derive the constrained density of $\eta_{Q}$, we first write the joint asymptotic density $f\left(n_{I}, n_{Q}\right)$ of ( $\eta_{Q}, \eta_{I}$ ), which by Theorem 3.3a becomes

$$
\begin{equation*}
f\left(n_{I}, n_{Q}\right)=\frac{\exp \left\{-\frac{1}{2\left(1-R^{2}\right)}\left(\frac{n_{I}^{2}}{\sigma_{I}^{2}}-2 R \frac{n_{I} n_{O}}{\sigma_{I} \sigma_{Q}}+\frac{n_{Q}^{2}}{\sigma_{Q}^{2}}\right)\right\}}{2 \pi \sigma_{I} \sigma_{Q} \sqrt{1-\overline{R^{2}}}} \tag{53}
\end{equation*}
$$

with $R=\frac{C_{I Q}}{\sigma_{I} \sigma_{Q}}$. Setting $\eta_{I}=-\eta_{Q}$, we finally obtain the asymptotic density of $\eta_{Q}$, as stated below.

Lemma 3.4a Under constraint (48), we have

$$
\begin{equation*}
\eta_{Q} \sim \mathcal{N}\left(0, \tilde{\sigma}_{Q}^{2}\right)+O\left(\frac{1}{\sqrt{n}}\right) \tag{54}
\end{equation*}
$$

with

$$
\begin{equation*}
\tilde{\sigma}_{Q}^{2}=\left(1-R^{2}\right)\left(\frac{1}{\sigma_{I}^{2}}+2 \frac{C_{I Q}}{\sigma_{I}^{2} \sigma_{Q}^{2}}+\frac{1}{\sigma_{Q}^{2}}\right)^{-1}=\sqrt{2} / 16+O(1 / n), \tag{55}
\end{equation*}
$$

where all the quantities in the above were defined before.
We delay the discussion of the number of paths $L(n, d)$ (cf. Theorem 2.5) until the next subsection since the recurrence on $L(n, d)$ is of the same kind as the one needed to study the behavior of $W_{n}$ in the case (B). It will turn out that the asymptotics of $L(n, d)$ for (A) can be deduced from the asymptotics of $L(n, d)$ obtained in case (B).

Finally, we prove our last result concerning the large deviation of the total weight distribution (cf. Theorem 2.6). As discussed in Section 2, we only consider two cases, namely: (a) identically distributed weights, that is, $W_{I}={ }^{d} W_{D}={ }^{d} W_{Q}=W$ where $={ }^{d}$ means equal in distribution; and (b) constant $D$-weight and $I$-weight, i.e., $W_{D}=W_{I}=-1$.

Let us first establish notation needed to express a large deviation result. Define $S_{n}=$ $\sum_{i=1}^{n} W(i)$ where $W(i)$ is an independent copy of $W$. Let $\Psi(z)=\log E e^{z(W-m)}$ be the cumulant function of $W-m$ where $m=E W$, and let $s$ be the unique solution, if exists, of the following equation

$$
a=\Psi^{\prime}(s)
$$

for any $a>0$. Finally, let $Z(a)=-\left(\Psi(s)-s \Psi^{\prime}(s)\right)$. Then (cf. Feller [18])

$$
\begin{equation*}
\operatorname{Pr}\left\{S_{n} \geq n(a+m)\right\} \sim \frac{1}{s \sqrt{2 \pi n \Psi^{\prime \prime}(s)}} \exp (-n Z(a)) \tag{56}
\end{equation*}
$$

In our case, the total weight $W_{n}\left(N_{Q}\right)$ of a random path in a grid graph with exactly $N_{Q}$ diagonal edges becomes $W_{n}\left(N_{Q}\right)=\sum_{i=1}^{N_{I}} W_{I}(i)+\sum_{i=1}^{N_{D}} W_{D}(i)+\sum_{i=1}^{N_{Q}} W_{Q}(i)=\sum_{i=1}^{n-N_{Q}} W(i)$ (cf. (47)). Note that $N_{Q}$ is a random variable, hence the unconditional total weight $W_{n}$
can be computed from an estimate of the conditional total weight $W_{n}\left(N_{Q}\right)$ and the limiting distribution of $N_{Q}$ (cf. Lemma 3.4a). But, $N_{Q}=n \mu_{Q}+\eta_{Q} \sqrt{n}$ and by Lemma 3.4a $\eta_{Q}$ is asymptotically normal with mean 0 and variance $\tilde{\sigma}_{Q}^{2}$. We must now translate (56) into our new situation. Let $\tilde{n}=\gamma n$ where $\gamma=1-\mu_{Q}$. Define $\tilde{a}$ such that

$$
\begin{aligned}
\tilde{n} a+m \sqrt{n} \eta_{Q} & =\left(\tilde{n}-\sqrt{n} \eta_{Q}\right) \tilde{a}, \quad \text { i.e. } \\
\tilde{a} & =a+\frac{(m+a)}{\gamma} \frac{\eta_{Q}}{\sqrt{n}}+\frac{m+a}{\gamma^{2}} \frac{\eta_{Q}^{2}}{n}+O\left(\frac{\eta_{Q}^{3}}{n^{3 / 2}}\right)
\end{aligned}
$$

Let also $s^{*}$ and $s$ be solutions of the following equations $a=\Psi^{\prime}\left(s^{*}\right)$ and $\tilde{a}=\Psi^{\prime}(s)$. Using Taylor's expansion of $\Psi(s)$ and $\Psi^{\prime}(s)$ around $s^{*}$, we obtain

$$
\begin{equation*}
s=s^{*}+\frac{a^{*}+m}{\gamma \psi^{\prime \prime}\left(s^{*}\right)} \frac{\eta_{Q}}{\sqrt{n}}-\frac{1}{2} \frac{\left(a^{*}+m\right)\left[-2\left(\psi^{\prime \prime}\left(s^{*}\right)\right)^{2}+(a+m) \psi^{\prime \prime \prime}\left(s^{*}\right)\right]}{\gamma^{2}\left(\psi^{\prime \prime}\left(s^{*}\right)\right)^{3}} \frac{\eta_{Q}^{2}}{n}+O\left(\frac{1}{n^{3 / 2}}\right) . \tag{57}
\end{equation*}
$$

With the notation as above, we reduce the problem to the following one

$$
\operatorname{Pr}\left\{W_{n} \geq \gamma(a+m) n\right\}=\int_{-\infty}^{\infty} \operatorname{Pr}\left\{\sum_{i=1}^{\bar{n}-\sqrt{n} \eta}(W(i)-m) \geq(\tilde{n}-\sqrt{n} \eta) \tilde{a} \mid \eta_{Q}=\eta\right\} d F_{\eta_{Q}}(\eta)
$$

where $F_{n Q}(\eta)=\Phi(\eta)(1+O(1 / \sqrt{n}))$ (cf. Lemma 3.4a) $\Phi(\cdot)$ stands for the normal distribution with mean zero and variance $\tilde{\sigma}_{Q}^{2}$. The probability under the above integral can be estimated as in (56). Using, in addition, the well known formula

$$
\int_{-\infty}^{\infty} \exp \left(-p^{2} x^{2} \pm q x\right) d x=\frac{\sqrt{\pi}}{p} \exp \left(\frac{q^{2}}{4 p^{2}}\right)
$$

after tedious algebra, we obtain our result (25) presented in Theorem 2.6.
In a similar manner we deal with the second case (b). However, this time the starting equation is $W_{n}\left(N_{Q}\right)=\sum_{i=1}^{N_{Q}} W_{Q}(i)-\left(n-2 N_{Q}\right)$. The details are left to the reader.

### 3.2 Case (B): $d=O(n)$

The main purpose of this section is to derive the limiting distribution of the total weight for a given path $\mathcal{P}$ in a grid graph $\vec{G} \in \overrightarrow{\mathcal{G}}$, and the asymptotics for the number of paths $L(n, d)$. As in the previous subsection, we proceed in three steps: at first, we consider an unweighted unconstraint random walk, then we derive probabilities $p_{I}, p_{D}$ and $p_{Q}$ for the constraint unweighted random walk, and finally we deal with the total weight $W_{n}$ -

Consider the unweighted random walk $Y(\cdot)$ in the grid graph as in Figure 2 such that $Y(n)=d=n x$ for some $x<1$. Naturally, in this domain of $d$ and $n$ we cannot use the normal approximation, which works only up to $O(\sqrt{n})$. We have to appeal to the large
deviation arguments to obtain the probability distribution of the random walk $Y(\cdot)$. We proceed along the lines of arguments suggested by Louchard [26].

We consider the constraint random walk $Y(n)=n x$, however, it is convenient to generalize our constraint to the following one

$$
\begin{equation*}
Y(m)=m u \tag{58}
\end{equation*}
$$

One can imagine that the random walk $Y(\cdot)$ at step $m$ has to be at position $m u$, where $m$ and $u$ are functions of $n$ and $x$ (e.g., we shall assume later that $m u=n x$ ).

As in the case (A), the analysis of the number of steps $N_{I}, N_{D}$ and $N_{Q}$ is crucial for the total weight. Note that, under our constraint (58), we have $N_{I}+N_{D}+2 N_{Q}=m$ and $N_{I}-N_{D}=m u$. The above can be translated to the following constraint: $N_{I}+N_{Q}=$ $\frac{m}{2}(1+u)$. Bearing this in mind, we transformed the random walk $Y(\cdot)$ into another random walk $\tilde{Y}(\cdot)$ that is defined in Figure 4 below (i.e., its one-step moves are shown in Fig. 4). Our interest lies in estimating $\operatorname{Pr}\{Y(m) \in m d u\}$ or in terms of the new random walk $\tilde{Y}(\cdot)$


Figure 4: Definition of the new random walk $\tilde{Y}(\cdot)$.
we evaluate the following

$$
\begin{equation*}
\operatorname{Pr}\{Y(m) \in m d u\} \rightleftharpoons \operatorname{Pr}\left\{\frac{\tilde{Y}(m)-m / 2}{m} \in \frac{d u}{2}\right\} \tag{59}
\end{equation*}
$$

To analyze $\bar{Y}(\cdot)$, we compute the probability $p_{i}(j)=\operatorname{Pr}\{\bar{Y}(i)=j\}$. It is easy to see that this probability satisfies the following recurrence

$$
\begin{equation*}
p_{i+1}(j)=\tau p_{i}(j)+\tau p_{i}(j-1)+\tau^{2} p_{i-1}(j-1), \quad i \geq 1 \tag{60}
\end{equation*}
$$

We solve this recurrence by the mean of generating function approach. Let $g_{i}(z)=$ $\sum_{j=0}^{\infty} z^{j} p_{i}(j)$. Clearly

$$
\begin{equation*}
g_{0}(z)=1, \quad g_{1}(z)=\tau(1+z) \tag{61}
\end{equation*}
$$

$$
g_{i+1}(z)=\tau g_{i}(z)+\tau z g_{i}(z)+\tau^{2} z g_{i-1}(z), \quad i \geq 1
$$

Let now $\varphi(\theta, z)=\sum_{i=0}^{\infty} \theta^{i} g_{i}(z)$, and after some algebra one obtains

$$
\begin{equation*}
\varphi(\theta, z)=\frac{1}{1-(1+z) \theta \tau-z \theta^{2} \tau^{2}} \tag{62}
\end{equation*}
$$

The roots of the denominator of the above become

$$
\theta_{1,2}(z)=-\frac{(1+z) \tau \pm \tau \sqrt{w_{1}(z)}}{2 z \tau^{2}}
$$

where $w_{1}(z)=1+6 z+z^{2}$. Then,

$$
\begin{equation*}
\varphi(\theta, z)=\left(\frac{\alpha_{1}(z)}{\theta-\theta_{1}(z)}+\frac{\alpha_{2}(z)}{\theta-\theta_{2}(z)}\right) . \tag{63}
\end{equation*}
$$

where $\alpha_{2}(z)=-\left(\tau \sqrt{w_{1}(z)}\right)^{-1}$, and $\alpha_{1}(z)=-\alpha_{2}(z)$.
To extract the generating function $g_{i}(z)$ from (63), we expand $\varphi(\theta, z)$ in the powers of $\theta$ to obtain

$$
\begin{equation*}
g_{m}(z)=-\frac{\alpha_{1}(z)}{\theta_{1}(z)}\left(\frac{1}{\theta_{1}(z)}\right)^{m}-\frac{\alpha_{2}(z)}{\theta_{2}(z)}\left(\frac{1}{\theta_{2}(z)}\right)^{m} \tag{64}
\end{equation*}
$$

Since we are interested in large values of $m$, we deduce from (64) that the leading term of the asymptotics can be extracted from the following

$$
\begin{equation*}
g_{m}(z) \sim \frac{1}{\theta_{2}^{m}(z)}=\psi_{1}^{m}(z), \quad m \rightarrow \infty \tag{65}
\end{equation*}
$$

with $\psi_{1}(z)=1 / \theta_{2}(z)$. In the above, we omitted the function $\alpha_{1}(z) / \theta_{1}(z)$ since it only contributes a constant in the final asymptotics.

Our aim now is to assess asymptotically the probability $p_{m}(k)=\operatorname{Pr}\{\tilde{Y}(m)=k\}$. Clearly, it can be estimated as $p_{m}(k) \sim\left[\psi_{1}^{m}(z)\right]_{k}$ where $[f(z)]_{k}$ is the coefficient of $z^{k}$ in the power expansion of $f(z)$. Hence, we have to deal with evaluating the $k$ th coefficient of of $\psi_{1}^{m}(z)$, where $k=m(1+u) / 2$. To obtain such asymptotics we shall use the classical "shift of the mean" technique (cf. Feller [18] p. 548 and Greene and Knuth [19] p.79). For the reader convenience, we discuss briefly this technique below. We follow the approach of Greene and Knuth [19].

Let $g(z)$ be the generating function of a random variable with mean equal to $\mu$ and the variance equal to $\sigma^{2}$. Then, $g^{n}(z)$ represents the generating function of the sum of $n$ such independent random variables. We estimate the coefficient of $z^{\mu n+r}$ in $g^{n}(z)$ for such $r$ that $\mu n+\tau$ is an integer. Call such a coefficient $A_{n, r}$. By the Cauchy formula Greene and Knuth [19] derive the following

$$
\begin{equation*}
A_{n, r}=\frac{1}{\sigma \sqrt{2 \pi n}} \exp \left(\frac{-\tau^{2}}{2 \sigma^{2} n}\right)+O\left(n^{3 \varepsilon-1}\right) \tag{66}
\end{equation*}
$$

where $\varepsilon$ is arbitrary small positive number. The reader should notice that this asymptotics is valid only for $r=O(\sqrt{n})$.

In our case, we need the $k$ th coefficient of $\psi_{1}^{m}(z)$, where $k=m(1+u) / 2$. Therefore, we cannot directly apply (66) since we are not in the range $O(\sqrt{m})$. A solution to this dilemma is proposed in [19] by a simple and elegant application of the "shift of the mean" technique, which we discuss below.

Let us return to Greene and Knuth \{19], and assume that one needs the $k$ th coefficient of $g^{n}(z)$. The shift of the mean technique computes the $k$ th coefficient as follows

$$
\begin{equation*}
\left[g^{n}(z)\right]_{k}=\frac{g(\beta)^{n}}{\beta^{k}}\left[\left(\frac{g(\beta z)}{g(\beta)}\right)^{n}\right]_{k}, \tag{67}
\end{equation*}
$$

where the parameter $\beta$ allows to shift the mean of the distribution to a value close to $k / n$, and hence allows to apply the asymptotics (66). The choice of $\beta$ is specified by the following equation

$$
\begin{equation*}
\frac{\beta g^{\prime}(\beta)}{g(\beta)}=\frac{k}{n} . \tag{68}
\end{equation*}
$$

Now, we are ready to derive our asymptotics. Since we seek the $k=m(1+u) / 2$ coefficient of $\psi_{1}^{m}(z)$, we first apply (68) to shift the mean. Define $\beta_{1}(u)$ as

$$
\begin{equation*}
\frac{\beta_{1} \psi_{1}^{\prime}\left(\beta_{1}\right)}{\psi_{1}\left(\beta_{1}\right)}=\frac{1+u}{2} \tag{69}
\end{equation*}
$$

Finally, applying (66), we obtain our main result.
Theorem 3.1b We have proved

$$
\begin{aligned}
\Psi_{1}(m, u) & =\operatorname{Pr}\{Y(m) \in m d u\}=\operatorname{Pr}\left\{\bar{Y}(m)-\frac{m}{2} \in \frac{m d u}{2}\right\} \\
& \sim \frac{\left(\psi_{1}\left(\beta_{1}\right)\right)^{m}}{\beta_{1}^{m(1+u) / 2} \sqrt{2 \pi m V(u)}} \frac{m d u}{2}(1+O(1 / n))
\end{aligned}
$$

for all $m$, where

$$
\begin{align*}
\frac{\left(\beta_{1}(u)-1\right) \psi_{1}(u)}{2 \tau \beta_{1}(u)+\left(1+\beta_{1}(u)\right) \psi_{1}(u)} & =u  \tag{70}\\
\beta_{1}(u) & =\frac{1+3 u^{2}+u \sqrt{8\left(u^{2}+1\right)}}{1-u^{2}}  \tag{71}\\
\psi_{1}(u) & =\psi_{1}\left[\beta_{1}(u)\right]=\frac{2 u \tau \beta_{1}(u)}{\beta_{1}(u)-1-u\left(1+\beta_{1}(u)\right)} \tag{72}
\end{align*}
$$

for all $u<1$.

Theorem 3.1b allows to analyze the constraint random walk $Y(n)=d$. In particular, as for case (A), we can compute the probabilities $p_{I}, p_{D}$, and $p_{Q}$ of one-step moves. Setting in Theorem 3.1b, $m=n t, u=\frac{x}{l}$ so that $m u=n x$, we obtain

$$
\begin{aligned}
\Psi_{2}(x, t) d x & =\operatorname{Pr}\{Y(n t) \in n d x\} \\
& \sim \exp \left\{n t\left[\log \psi_{1}\left(\frac{x}{t}\right)-\frac{1}{2} \log \beta_{1}\left(\frac{x}{t}\right)\right]-\frac{n x}{2} \log \beta_{1}\left(\frac{x}{t}\right)\right\} \frac{\sqrt{n} \lambda\left[\beta_{1}\left(\frac{x}{t}\right)\right]}{2 \sqrt{2 \pi t V\left(\frac{x}{t}\right)}} \cdot d x
\end{aligned}
$$

This implies, for example, that

$$
p_{I} \sim \frac{1}{3}\left[\frac{\Psi_{2}\left(x-\frac{1}{n}, t-\frac{1}{n}\right)}{\Psi_{2}(x, t)}\right]_{t=1}
$$

and in a similar fashion for $D$ and $Q$. After some algebra, we finally derived the following lemma.

Lemma 3.2b. The probabilities $p_{I}, p_{D}$ and $p_{Q}$ become

$$
\begin{aligned}
p_{I} & =\frac{\tau \beta_{1}(x)}{\psi_{1}(x)}+O(1 / n) \\
p_{D} & =\frac{\tau}{\psi_{1}(x)}+O(1 / n) \\
p_{Q} & =\frac{\tau^{2} \beta_{1}(x)}{\psi_{1}^{2}(x)}+O(1 / n)
\end{aligned}
$$

for all $x<1$.
Concerning the limiting joint distribution of the number of $I$-steps, $D$-steps and $Q$-steps, we proceed as before. We use the same notation as in Theorem 3.3 with appropriate values for probabilities $p_{I}, p_{D}$ and $p_{Q}$ from Lermma 3.2 b . This leads to the following results.

Theorem 3.3b. The number of $I, D$ and $Q$ steps, $N_{I}, N_{D}, N_{Q}$ respectively, are asymptotically Gaussian, with mean $n \mu_{I}, n \mu_{D}, n \mu_{Q}$ respectively, where these quantities are computed according to (49)-(51), (52) with new probabilities $p_{I}, p_{D}$ and $p_{Q}$, as in Lemma 3.2b.

Lemma 3.4b We have $\eta_{Q}=\mathcal{N}\left(0, \tilde{\sigma}_{Q}^{2}\right)$ with $\tilde{\sigma}_{Q}^{2}$ given by (55) with probabilities $p_{I}, p_{D}$ and $p_{Q}$ as in Lemma 3.2b.

Finally, we prove Theorem 2.5 that enumerates the total number of path $L(u)$. As discussed in Section 2, this estimate is necessary to evaluate our upper bound $\bar{\alpha}$ in Theorem 2.7. We start the analysis with setting up a recurrence for $L(u)$. Let $f_{i}(j)$ be the total number of paths from $O$ to $j$ in $i$ steps of the associated random walk in our grid graph
$\vec{G}$. Then, $L(u)=f_{n}(d)$. Hereafter, we set $d=u n$. Clearly, $f_{i}(j)$ satisfies the following recurrence

$$
\begin{equation*}
f_{i+1}(j)=f_{i}(j)+f_{i}(j-1)+f_{i-1}(j-1) \tag{73}
\end{equation*}
$$

with $f_{1}(1)=1$. This recurrence was already studied by Laquer [24] for $d=O(n)$. Observe that the above recurrence is similar to the one consider before, and we can use the same technique to attack it. Set $g_{i}(z)=\sum_{j=0}^{\infty} z^{j} f_{i}(j)$ and let $\varphi(\theta, z)=\sum_{i=0}^{\infty} \theta^{i} g_{i}(z)$. After the some algebra we obtain

$$
\begin{equation*}
\varphi(\theta, z)=\frac{\alpha_{1}(z)}{\theta-\theta_{1}(z)}+\frac{\alpha_{2}(z)}{\theta-\theta_{2}(z)} \tag{74}
\end{equation*}
$$

with $w_{2}(z)=1+z^{2}+6 z$ and

$$
\begin{equation*}
\theta_{1,2}(z)=\frac{1+z \pm \sqrt{w_{2}(z)}}{-2 z} \tag{75}
\end{equation*}
$$

where $\alpha_{2}(z)=-1 / \sqrt{w_{2}(z)}$ and $\alpha_{1}(z)=-\alpha_{2}(z)$. As $n$ becomes larger, the dominant contribution comes from (74), and asymptotically we have

$$
g_{n}(z) \sim \lambda(z)\left[\frac{1}{\theta_{2}(z)}\right]^{n}=\lambda(z) \psi_{2}^{n}(z)
$$

with $\lambda(z)=-\alpha_{2}(z) / \theta_{2}(z)$.
To extract the coefficient of $g_{n}(z)$ we shall apply the "shift of the mean" method, as described before. We first consider only the coefficient at $g_{n}(z) / \lambda(z)=\theta_{2}^{-n}(z)$. Call it $l(u)$. Applying equation (67), as in (69), we estimate the new mean value with $\psi_{1}(z)$ replaced by $\psi_{2}(z)$ and the new $\beta_{2}(u)$ becomes

$$
\begin{equation*}
\beta_{2}(u)=\frac{1+3 u^{2}+u \sqrt{8\left(u^{2}+1\right)}}{1-u^{2}} \tag{76}
\end{equation*}
$$

and then

$$
\psi_{2}(u)=\psi_{2}\left[\beta_{2}(u)\right]=\frac{2 u \beta_{2}(u)}{\beta_{2}(u)-1-u\left(1+\beta_{2}(u)\right)} .
$$

Let $V(u)$ be the variance related to the generating function $\frac{\psi_{2}\left(\beta_{2 z}\right)}{\psi_{2}\left(\beta_{2}\right)}$ as defined in Theorem 2.5. With this in mind, it is easy to see that

$$
\begin{equation*}
l(u)=\frac{\psi_{2}\left(\beta_{2}(u)\right)^{n}}{\beta_{2}(u)^{\mathfrak{n}(1+u) / 2} \sqrt{2 \pi n V(u)}}(1+O(1 / n))=\frac{\exp (n \rho(u))}{\sqrt{2 \pi n V(u)}}(1+O(1 / n)) \tag{77}
\end{equation*}
$$

where $\rho(u)$ is a function of $u$ and it is given by (20).
Now, we compute the coefficient at $g_{n}(z)=\lambda(z) \theta_{2}^{-n}(z)$, that is, we include the correction coming from $\lambda(z)$. Note that $\lambda_{1}(z)=\lambda(z) / \lambda(1)$ can be viewed as the generating function of a random variable. Let its probability distribution be denoted by $p_{\lambda}(i)$. Since the product
of two generating functions translates into the convolution of the appropriate coefficients, we have $L(u)=\sum_{i=0}^{\infty} p_{\lambda}(i) l(u-i / n)$. By (77) we finally obtain

$$
\begin{align*}
L(u) \sqrt{2 \pi n V(u)} & =\lambda(1) \sum_{i=0}^{\infty} p_{\lambda}(i) \exp \left(n\left(\rho(u)-\rho^{\prime}(u) i / n+O\left(n^{-2}\right)\right)\right)\left(1+O\left(n^{-1}\right)\right) \\
& =\lambda\left(e^{-\rho^{\prime}(u)}\right) \exp (n \rho(u))\left(1+O\left(n^{-1}\right)\right) \tag{78}
\end{align*}
$$

where $\rho^{\prime}(u)$ is the derivative of $\rho(u)$. From the above, we conclude that the constant $C$ in Theorem 2.5 becomes

$$
\begin{equation*}
C=\lambda\left(e^{-\rho^{\prime}(u)}\right) \tag{79}
\end{equation*}
$$

This completes the proof of Theorem 2.5 for case (A) and (B) (in case (A) $\rho(u)$ is given by (19)).

## APPENDIX A: Proof of Theorem 3.3a

From (42) and (47) after setting $N_{T}=N_{I}+N_{D}$, we see that

$$
\begin{align*}
& \frac{n}{2}+\frac{N_{T}}{2}=n-N_{Q} \sim \mathcal{N}(n \alpha, n \kappa)+O(1)  \tag{80}\\
& \mu_{T}=E\left(N_{T}\right) / n=n(2 \alpha-1)+O(1)  \tag{81}\\
& \sigma_{T}^{2}=\operatorname{VAR}\left(N_{T}\right) / n \sim 4 \kappa
\end{align*}
$$

But given $N_{T}$, the number of $I$-steps $N_{I}$ is a binomial random variable with parameter $\bar{p}_{I}$, and mean $N_{T} \bar{p}_{I}$ and the variance $N_{T} \bar{p}_{I} \bar{q}_{I}$ (where $\bar{q}_{I}=1-\bar{p}_{I}$ ). By (81) we have $E\left(N_{I}\right)=n \mu_{T} \bar{p}_{I}=n \bar{p}_{I}(2 \alpha-1)$. We also obtain $E\left(N_{I}^{2} \mid N_{T}\right)=N_{T} \bar{p}_{I} \bar{q}_{I}+N_{T}^{2} \bar{p}_{I}^{2}$ and $E\left(N_{I}^{2}\right)=$ $n \mu_{T} \bar{p}_{I} \bar{q}_{I}+\bar{p}_{I}^{2}\left(n \sigma_{T}^{2}+n^{2} \mu_{T}^{2}\right)$. This finally leads to

$$
\sigma_{I}^{2}=n\left(\mu_{T} \bar{p}_{I} \bar{q}_{I}+\bar{p}_{I}^{2} \sigma_{T}^{2}\right) \sim n\left((2 \alpha-1) \bar{p}_{I} \bar{q}_{I}+\bar{p}_{I}^{2} 4 \kappa\right)
$$

The number of $D$-steps is analyzed in a similar fashion. To compute the covariance $C_{I D}$ between $N_{I}$ and $N_{D}$, note that

$$
n \sigma_{T}^{2}=\operatorname{VAR}\left(N_{T}\right)=\operatorname{VAR}\left(N_{I}+N_{D}\right)=n\left(\sigma_{I}^{2}+\sigma_{D}^{2}+2 C_{I D}\right)
$$

or $4 \kappa \sim 2(2 \alpha-1) \bar{p}_{I} \bar{q}_{I}+4 \kappa\left(\bar{p}_{I}^{2}+\bar{p}_{D}^{2}\right)+2 C_{I D}$. Finally,

$$
\operatorname{COV}\left(N_{I} N_{T}\right)=\bar{p}_{I} E\left(N_{T}^{2}\right)-E\left(N_{I}\right) E\left(N_{T}\right)=n \bar{p}_{I} \sigma_{T}^{2} \sim \tilde{p}_{I} 4 \kappa n,
$$

and with (80), we obtain $\operatorname{COV}\left(N_{I} N_{Q}\right)=-\frac{1}{2} \operatorname{COV}\left(N_{I} N_{T}\right) \sim-\bar{p}_{I} 2 \kappa n$.

To complete the proof, it suffices to check the asymptotic Gaussian property of $N_{I}, N_{D}, N_{Q}$. For $N_{Q}$, this follows from (80). For $N_{I}$, which is binomially distributed with parameter $\bar{p}_{I}$, we obtain, conditioning on $N_{T}$

$$
\begin{align*}
E\left[e^{i \theta N_{I} / \sqrt{n}}\right] & =E\left\{\left[\bar{p}_{I} e^{i \theta / \sqrt{n}}+\bar{q}_{I}\right]^{N_{T}}\right\}=E\left\{\left[1+\bar{p}_{I} i \frac{\theta}{\sqrt{n}}-\bar{p}_{I} \frac{\theta^{2}}{2 n}+O\left(\frac{1}{n^{3 / 2}}\right)\right]^{N_{T}}\right\} \\
& =E\left\{\exp N_{T}\left[i \bar{p}_{I} \frac{\theta}{\sqrt{n}}-\frac{\theta^{2}}{2 n} \bar{p}_{I} \ddot{q}_{I}+O\left(\frac{1}{n^{3 / 2}}\right)\right]\right\} \tag{82}
\end{align*}
$$

But, by (80),

$$
E\left[e^{i \rho N_{T}}\right]=e^{i \pi \mu_{T} \rho-\frac{1}{2} \pi \sigma_{T}^{2} \rho^{2}+O\left(\rho^{3} \frac{n^{3 / 2}}{\sqrt{n}}\right)}
$$

hence by (82) we obtain

$$
E\left[e^{i \theta N_{I} / \sqrt{n}}\right]=\exp \left\{i n \mu_{T} \frac{\theta}{\sqrt{n}} \bar{p}_{I}-\frac{\theta^{2}}{2}\left(\bar{p}_{I} \bar{q}_{I} \mu_{T}+\sigma_{T}^{2} \bar{p}_{I}^{2}\right)+O\left(\frac{1}{\sqrt{n}}\right)\right\}
$$

which proves the asymptotic Gaussian property of $N_{I}$.

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[^1]:    ${ }^{1}$ To simplify our notation, we often write $X$. to denote any of $X_{I}, X_{D}$ or $X_{Q}$.

