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# AVERAGE PROFILE OF THE LEMPEL-ZIV PARSING SCHEME FOR A MARKOVIAN SOURCE* 

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#### Abstract

For a Markovian source, we analyze the Lempel-Ziv parsing scheme that partitions sequences into phrases such that a ncw phrase is the shortest phrase not seen in the past. We consider three models: In the Markov Independent model, several sequences are generated independently by Markovian sources, and the $i$ th phrase is the shortest prefix of the $i$ th sequence that was not seen before as a phrase (i.c., a prefix of previous ( $i-1$ ) sequences). In the other two models, only a single sequence is generated by a Markovian source. In the second model, for which we coin the name Gilbert-Kadota model, a fixed number of phrases is generated according to the Lempel-Ziv algorithm, thus producing a sequence of a variable (random) length. In the last model, known also as the Lempel-Ziv model, a string of fixed length is partitioned into a variable (random) number of phrases. These three models can be efficiently represented and analyzed by digital search trees that are of interest to other algorithms such as sorting, searching and pattern matching. In this paper, we concentrate on analyzing the average profile (i.e., the average number of phrases of a given length), the typical phrase length, and the length of the last phrase. We obtain asymptotic expansions for the mean and the variance of the phrase length, and we prove that appropriately normalized phrase length in all thrce models tends to the standard normal distribution which lead to bounds on the average redundancy of the Lempel-Ziv code. For Markov Independent model, this finding is established by analytic methods (i.e., generating functions, Mellin transform and depoissonization), while for the other two models we use a combination of analytic and probabilistic analyses.


Index Terms: Lempel-Ziv scheme, Markov source, digital search trees, data compression, phrase length, depth in a tree, Poisson transform, Mellin transform, analytic depoissonization, stochastic comparisons.

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## 1 Introduction

The heart of many lossless data compression schemes is the incremental parsing algorithm due to Lempel and Ziv [24]. It partitions a sequence into variable phrases such that a new phrase is the shortest substring not seen in the past as a phrase. Revealing its intrinsic bchavior should lead to a better understanding of the internal structure of sequences, and this is of a broader interest to molecular biology, languages, coding, efficient data transmission, estimation of entropy, discrimination between information sources, test of randomness, estimation of the statistical model for individual sequences, multimedia compression, and so forth. Fundamental information about the algorithm is contained in such paramcters as the number of phrases, the phrase length, the number of phrases of a given size, and the longest phrase. Herc, for Markovian sources we study the length of a randomly selected phrase (which is equivalent to the so called average profile defined as the average number of phrases of a given size) and the length of the last phrase.

In the past, mostly first order analysis of these parameters were available with the exception of $[7,11,12,18]$ where largely memoryless sources were analyzed. The first order analysis provides the first order asymptotics (e.g., is the redundancy of a code $o(n)$ ?). The second order analysis attempts to establish the rate of convergence, or even a full asymptotic expansion, large deviations behavior, deviation from the mean (e.g., central limit theorems), and so forth. We present here a second order analysis of the (typical) phrase length for the Lempel-Ziv parsing scheme in a Markovian setting.

One can still wonder why do we need a second order analysis or a second order approximation of information systems. Gilbert and Kadota in [4] and J. Ziv in his 1997 Shannon Lecture provided some convincing arguments for the need of such investigations. In fact, J. Ziv presented compelling arguments for "backing off" to a certain degree from the first-order asymptotic analysis of information systems in order to predict the behavior of real systems where we always face finite, and often small, lengths (of sequences, files, codes, etc.) One way of overcoming these difficulties is to increase the accuracy of asymptotic analysis by replacing first-order analysis by full asymptotic expansions and more accurate analysis so that the approximate value of a quantity of interest is closer to the true value even for moderate and small lengths. On the other hand, Kadota and Gilbert used a numerical evaluation (instead of a crude first order asymptotic) to obtain qualitative insights into the behavior of the Lempel-Ziv algorithm. Some of their results were analytically recovered in [7, 11] where second order asymptotics were obtained for the quantities studied in [4]. In [4, 7, 11] only memoryless sources were analyzed, and in this paper we extend the analysis to Markovian
sources.
In this paper, we shall analyze three models of the Lempel-Ziv scheme in the Markovian settings. In the first one, called Markov Independent model or shortly MI model, we assume that there are $m$ independent Markov sources defined on the same underlying probability space. The parsing is done with respect to the previous sequences. Namely, the zeroth phrase is an empty phrase, while the first phrase is a one character prefix of the first sequence. The $i$ th phrase ( $i \leq m$ ) is defined as the shortest prefix of the $i$ th sequence not seen as a phrase (prefix) of the previous ( $i-1$ ) sequences. For example, for $m=4$ sequences: $X(1)=000000 \ldots, X(2)=1010101 \ldots, X(3)=1001101 \ldots$ and $X(4)=001100111 \ldots$ we can construct the following Lempel-Ziv sequence: $(\epsilon)(0)(1)(10)(00)$ where $\epsilon$ is an empty phrase, and all phrases are shown in parentheses. We shall study two parameters, namely the length, $D_{m}$, of a randomly selected phrase, and the length $I_{m}$ of the last phrase. In addition, one may investigate the length $L_{m}$ of the Lempel-Ziv sequence. In the example above we have $D_{4}=1 \frac{1}{2}, I_{4}=2$ and $L_{4}=6$.

The next two models deal with a single sequences generated by a Markovian source. In the fixed number of phrases model, we partition the sequence according to the Lempel-Ziv algorithm until we obtain $m$ full phrases (thus producing a variable and random length of the Lempel-Ziv sequence). For example, for $X=11001010001000100 \ldots$ we can construct $m=5$ phrases as follows: $(\epsilon)(1)(10)(0)(101)(00)$. Such a model was also considered by Gilbert and Kadota [4], so we call it the Gilbert-Kadota model or shortly GK model. As before, we will be interested in the typical phrase length $D_{m}$ and the last phrase length $I_{m}$. In the above example, we have $D_{5}=1 \frac{1}{5}, I_{5}=2$, and in addition the length of the Lempel-Ziv sequence is $L_{5}=9$.

Finally, in the traditional Lempel-Ziv model or fixed length model, a sequence of fixed length, say $n$ symbols, is partitioned according to the Lempel-Ziv algorithm. For example, the string $X=110010100010$ of length $n=12$ is parsed as $(\epsilon)(1)(10)(0)(101)(00)(01)(0)$. We shall study the length $\Delta_{n}$ of the randomly selected phrase (see Section 2 for a precise definition) and the length $J_{n}$ of the last full phrase. The number of full phrases $M_{n}$ is of significant intercst for this model, but we will not investigate it here. In the example above, $\Delta_{12}=1 \frac{5}{6}, J_{12}=2$ and $M_{12}=6$.

The above three models can be efficiently analyzed and uniformly represented by a digital search tree, a data structure that have been studied by its own right for more than thirty years (cf. [10, 14]). This tree is used to store strings in its nodes and can be described as follows: We consider $m$, possibly infinite, strings of symbols over a finite alphabet $\mathcal{A}=\{1,2, \ldots, V\}$ (however, we often restrict our discussions to a binary alphabet $\mathcal{A}=\{0,1\}$ ). The root


Markov Independent Model

## Lempel-Ziv Model

Figure 1: Digital tree representations for the MI model $\left(X(1)=00000, X(2)=01111, X_{3}=\right.$ 101010, $X(4)=111000, X(5)=110111, X(6)=111111)$ and the LZ model $(X=$ $11001010001000100 \ldots$. . of the Lempel-Ziv algorithm.
contains the empty string $\epsilon$. The first string occupics the right or the left child of the root depending whether its first symbol is " 1 " or " 0 ". The remaining strings are stored in available nodes (that are directly attached to nodes already existing in the tree). The scarch for an available node follows the prefix structure of a string. The rule is simple: if the next symbol in a string is " 1 " we move to the right, otherwise move to the left. The resulting tree has $m$ internal nodes. It corresponds to the MI model and the GK model, however, in the latter the strings are substrings (phrases) of one infinite string We can call such a digital search tree a suffix search tree (cf. Figure 1).

In the LZ model, we construct an analogous (suffix) digital tree except that the number of nodes varics and equals to the number of phrases $M_{n}$. More precisely, the empty phrase is stored in the root, and all other phrases are located in nodes. When a new phrase is created, the search starts at the root and proceeds down the tree as directed by the input symbols exactly in the same manner as in the digital search tree construction. For example, for the binary alphabet, " 0 " in the input string means move to the left and " 1 " means proceed to the right. The search is completed when a branch is taken from an existing tree node to a new node that has not been visited before. Then, an edge and a new node are added to the
tree. Phrases created in such a way are stored directly in nodes of the tree (cf. [11]). This is illustrated in Figure 1.

As mentioned before, in this paper we present second order analysis of the above three models of the Lempel-Ziv algorithm for a Markovian source. Among others, we compute precise asymptotic formulæ for the mean and the variance of the phrase length in the MI model. We also show that the appropriately normalized phrase length tends to a normal distribution with the rate of convergence of $O(1 / \sqrt{\ln m})$. These results - which are at the heart of our findings - are established by analytic methods. The line of the attack can be brielly described as follows: We first derive a set of recurrence equations for the ordinary generating functions of the average profile (conditioned on the first symbol). These recurrence equations are too complicated to be solved directly, hence we derive a set of differentialfunctional equations on the so called Poisson transform of the average profile. In the Poisson model, the number of sequences $m$ becomes a random variable $N$ distributed as a Poisson with mean $m$. This process of replacing the deterministic input $m$ by a Poisson variable is called poissonization. We shall use analytic poissonization since we replace $m$ by a complex variable $z$. A typical sct of differential-functional equations we have to deal with is of the following form

$$
\frac{\partial \widetilde{B}^{i}(u, z)}{\partial z}+\widetilde{B}^{i}(z, u)=u\left(\widetilde{B}^{1}\left(u, p_{i, 1} z\right)+\cdots+\tilde{B}^{V}\left(u, p_{i, V} z\right)\right)+a(z, u), \quad i=1,2, \ldots, V,
$$

where $\widetilde{B}^{i}(z, u)$ is the Poisson transform (cf. [7, 21]) of the average profile when all strings start with symbol $i \in \mathcal{A}=\{1,2, \ldots, V\}, a(z, u)$ is a given function, and $P=\left\{p_{i j}\right\}_{i, j=1}^{Y}$ is the underlying Markov chain. These differential-functional equations are reduced to a simple matrix functional equations of the Mellin transform $B_{i}^{*}(s)$ with respect to $z$ of $\widetilde{B}^{i}(z, u)$ (cf. [3, 21]). A typical equation of the Mellin transform looks like

$$
B_{i}^{*}(s)-(s-1) B_{i}^{*}(s-1)=B_{1}^{*}(s) p_{1, i}^{-s}+\cdots+B_{V}^{*}(s) p_{i, V}^{-s}+a^{*}(s), \quad i=1,2, \ldots, V
$$

We can solve exactly this matrix equation in a form of an infinite product of matrices. However, we develop a method to obtain relevant asymptotics without an explicit solution. It turns out that such asymptotics depend on singularity points of the matrix $Q(s)=(1-P(s))^{-1}$ where $\mathrm{P}(s)=\left\{p_{i j}^{-s}\right\}_{i, j=1}^{V}$ for some complex $s$. Then, through the inverse Mellin transform we obtain asymptotics of the Poisson transform $\tilde{B}^{i}(z, u)$ for large $z$. We need to translate it into the asymptotics of the original generating function $B_{m}^{i}(u)$. This process is called depoissonization, and we shall use recent results of Jacquet and Szpankowski [9] on analytic depoissonization. The program just described was recently dubbed analytic information theory since it applies analytic methods to solve problems of information theory.

To translate the results of the MI model to GK model and LZ model we shall use a combination of analytic, combinatorial and probabilistic methods. In particular, we construct two MI models that upper bound and lower bound stochastically the GK model. This will allow us to conclude the central limit theorem for the phrase length in the GK model, which will further lead to a similar result for the LZ model.

Finally, we should mentioned that our MI model is equivalent to the Markov model of digital search trees studied extensively in computer science. In fact, digital trecs appear in a variety of computer and communications applications including searching, sorting, dynamic hashing, codes, conflict resolution protocols for multiaccess communications, and data compression (cf, [10, 14]) Thus, better understanding of their behavior is desirable and could lead to some algorithmic improvements. One parameter that is of interest to these applications is the depth of a randomly selected node (i.e., the length of the path from the root to the chosen node), and depth of insertion, which may represent the search time. Clearly, the depth and the depth of insertion are equivalent to the typical phrase length and the last phrase length in the MI model. The average profile of the MI model is the same as the average number of nodes at a given level in the associated digital tree.

Digital trees (which include tries, PATRICIA tries and digital search trees) have been studied extensively in the past for memoryless source (cf. [10, 7, 11, 13, 14, 17, 20]). Extension to Markovian sources are scarce, and to the best of our knowledge only trics were analyzed (cf. [2, 6]). Lempel-Ziv model for memoryless sources was discussed in [7, 11, 12], while second order analysis for Markovian sources are basically non-existing. Savari [18] proposed the redundancy analysis of the LZ code for Markovian sources, but redundancy analysis requires rather a minor extension of the first order analysis. Wyner [23] derived the limiting distribution of the phrase length in the other Lempel-Ziv scheme (i.e., LZ'77) which is known to be considerable simpler for analysis.

This paper is organized as follows. In the next section we present our main results for all three models, and discuss some of their consequences. The proof for the MI model can be found in Section 3, while Section 4 presents our analysis of the GK model. The proof of the LZ model is discussed after Theorem 3 in Section 2.

## 2 Main Results

We present here our main results for all three models, namely Markov Independent model, Gilbert-Kadota (fixed number of plrases) model, and Lempel-Ziv model. Most of the proofs are delayed till the next section. Throughout, we assume that a sequence, say
$X=\left(X_{0}, X_{1}, \ldots\right)$, is generated by a Markov source over a finite alphabet $\mathcal{A}=\{1,2, \ldots, V\}$. More precisely:

## (M) Markov Source

There is a Markovian dependency between consecutive symbols in a sequence, that is, the probability $p_{i j}=\operatorname{Pr}\left\{X_{k+1}=j \mid X_{k}=i\right\}>0$ for all $k \geq 0$ describes the conditional probability of sampling symbol $j \in \mathcal{A}$ immediately after symbol $i \in \mathcal{A}$. We denote by $\mathrm{P}=\left\{p_{i j}\right\}_{i, j=1}^{V}$ the transition matrix, and by $\pi=\left(\pi_{1}, \ldots, \pi_{V}\right)$ the stationary vector satisfying $\pi \mathrm{P}=\pi$. We say that the Markov chain is stationary if $\operatorname{Pr}\left\{X_{k}=i\right\}=\pi_{i}$ for all $k \geq 0$ and $i \in \mathcal{A}$. In gencral, $X_{k+1}$ may dependent on last $r$ symbols, and then we have $r$ th order Markov chains, however, hereafter we mostly restraint ourselves to $r=1$.

### 2.1 Markov Independent Model

Hereafter, we assume that $m$ independent Markov sources generate $m$ sequences which are parsed with respect to previous ones according to the Lempel-Ziv algorithm, as described in the Introduction. Equivalently, we build a digital scarch tree from these $m$ sequences, as shown in Figure 1. Actually, it is more convenient to think in terms of this associated digital search tree (in short: DST). In particular, the $i$ th phrase length $I_{i}$ is also the depth of the $i$ th node in such a tree (where the depth of a node is understood as the number of nodes from the root to the $i$ th node). When $i=m$ we shall refer to $I_{m}$ as the depth of insertion or the last phrase length. The typical depth (typical phrase length) $D_{m}$ is defined as the length of a randomly selected depth, that is

$$
\operatorname{Pr}\left\{D_{m}=k\right\}=\frac{1}{m} \sum_{i=1}^{m} \operatorname{Pr}\left\{I_{i}=k\right\}
$$

Finally, we defined the average profile (in short: profile) $B_{m}^{k}$ as the average number of nodes at level $k$ of the DST or the average number of phrases of length $k$. Obscrve that $B_{0}^{k}=0$ for all $k \geq 0$

There are simple relationships between just defined parameters. First of all, we notice that (cf. $[10,11,20]$ )

$$
\begin{equation*}
\operatorname{Pr}\left\{D_{m}=k\right\}=\frac{B_{m}^{k}}{m} \tag{1}
\end{equation*}
$$

This, and the definition of the typical depth, immediately imply

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{m+1}=k\right\}=B_{m+1}^{k}-B_{m}^{k} \tag{2}
\end{equation*}
$$

with $\operatorname{Pr}\left\{I_{0}=0\right\}=1$ and $\operatorname{Pr}\left\{I_{0}=k\right\}=0$ for all $k \geq 1$.
Throughout, we shall work with generating functions of the above quantities and the so called Poisson transforms that we define next. The ordinary generating functions are

$$
\begin{array}{rlr}
D_{m}(u) & =\mathbf{E}\left[u^{D_{m}}\right]=\sum_{k \geq 0} \operatorname{Pr}\left\{D_{m}=k\right\} u^{k}, & D_{0}(u)=1, \\
I_{m}(i) & =\mathbf{E}\left[u^{I_{m}}\right]=\sum_{k \geq 0} \operatorname{Pr}\left\{I_{m}=k\right\} u^{k}, & I_{0}(u)=1, \\
B_{m}(u) & =\sum_{k \geq 0} B_{m}^{k} u^{k} & B_{0}(u)=0
\end{array}
$$

for a complex $u$ such that $|u|<1$. The Poisson transforms are defined as follows

$$
\begin{aligned}
\tilde{D}(z, u) & =\sum_{m \geq 0} D_{m}(u) \frac{z^{m}}{m!} e^{-z}, \\
\tilde{B}(z, u) & =\sum_{m \geq 0} B_{m}(u) \frac{z^{m}}{m!} e^{-z}, \\
\tilde{I}(z, u) & =\sum_{m \geq 0} I_{m}(u) \frac{z^{m}}{m!} e^{-z} .
\end{aligned}
$$

The Poisson transform can be interpreted as the generating function in the so called Poisson model in which the deterministic number of sequences $m$ is replaced by a random number of sequences distributed according to Poisson with mean $z=m$. We shall assume that $z$ is a complex variable, and $\widetilde{B}(z, u)$ as well as $\tilde{I}(z, u)$ are defined on the whole complex plane. We should also observe that by (2)

$$
\begin{equation*}
\frac{\partial \widetilde{I}(z, u)}{\partial z}+\tilde{I}(z, u)=\frac{\partial \widetilde{B}(z, u)}{\partial z} \tag{3}
\end{equation*}
$$

Since also $D_{m}(u)=B_{m}(u) / m$, we can recover all results on the depth of insertion $I_{m}$ as well as on the typical depth from the average profile $B_{m}^{k}$. Therefore, hereafter we concentrate on the analysis of the average profile.

To start the analysis, we derive a system of recurrence equations for the generating function of the average profile. We first need one more notation. Let $B_{m}^{i}(u)$ for $i \in \mathcal{A}$ be the ordinary generating function of the average profile when all sequences start with symbol $i$. Let also $\mathbf{p}=\left(p_{1}, \ldots, p_{V}\right)$ be the initial probability vector of the underlying Markov chain, that is, $\operatorname{Pr}\left\{X_{0}=i\right\}=p_{i}$. (For the stationary Markov chain we have $\mathrm{p}=\pi$.) Consider now the generating function $B_{m+1}(u)$ of the DST in which the root contains an empty string and the other $m$ independent Markov sequences are stored in $V$ subtrees, which are digital search trees by themselves but of smaller size. Indeed, the probability that the first subtree contains
$j_{1}$ sequences, the second subtree has $j_{2}$ sequences, and so on until the $V$ subtree stores $j_{V}$ sequences (out of $m$ sequences) is equal to the multinomial distribution, that is,

$$
\binom{m}{j_{1}, \ldots j_{V}} p_{1}^{j_{1}} \cdots p_{V}^{j_{V}} .
$$

But, the $i$ th subtree is again a digital search tree of size $j_{i}$ containing only those sequences that start with symbol $i$. Hence, its average profile generating function must be $B_{j_{1}}^{i}(u)$. This leads to the following recurrence equation assuming $B_{0}(u)=0$

$$
\begin{equation*}
B_{\pi+1}(u)=u \sum_{|\mathrm{j}|=m}\binom{m}{\mathbf{j}} p_{1}^{j_{1}} \cdots p_{V}^{j v}\left(B_{j_{1}}^{1}(u)+\cdots+B_{j v}^{V}(u)\right)+1 \tag{4}
\end{equation*}
$$

where $\mathbf{j}=\left(j_{1}, \ldots, j_{V}\right), \underline{\mathbf{j}} \mid=j_{1}+\cdots+j_{V}$ and for simplicity $\binom{m}{\mathbf{j}}=\binom{m}{j_{1}, \ldots j_{V}}$. Clearly, we can set up similar recurrences for the subtrees. That is,

$$
\begin{equation*}
B_{m+1}^{i}(u)=u \sum_{|\mathbf{j}|=m}\binom{m}{\mathbf{j}} p_{i 1}^{j_{1}} \ldots p_{i V}^{j V}\left(B_{j_{1}}^{1}(u)+\ldots+B_{j_{V}}^{V}(u)\right)+1, \quad \text { for all } \quad i \in \mathcal{A} \tag{5}
\end{equation*}
$$

where $B_{0}^{i}(u)=0$ for $i \in \mathcal{A}$.
If we can solve the above recurrences, then we can compute all moments and the distribution of the average profile, and consequently the characteristics of the typical depth and the depth of insertion. Indeed, after observing that $B_{m}(1)=m$, the average depth becomes $\mathrm{E}\left[D_{m}\right]=B_{m}^{\prime}(1)$ and

$$
\operatorname{Var}\left[D_{\pi n}\right]=\frac{B_{m}^{\prime \prime}(1)}{m}+\frac{B_{m}^{\prime}(1)}{m}-\left(\frac{B_{m}^{\prime}(1)}{m}\right)^{2}
$$

where $B_{m}^{\prime}(1)$ and $B_{m}^{\prime \prime}(1)$ are the first and the second derivatives of the generating function $B_{m}(u)$ calculated at $u=1$. In passing, we should observe that $B_{m}^{\prime}(1)$ and $B_{m}^{\prime \prime}(1)$ satisfy recurrences equations similar to the ones derived for $B_{m}(u)$, and we shall discuss them in details in the next section.

One must say, however, that the above recurrence equations are not easy to solve. Even, if in principle, one can write an explicit solution (cf. [11, 20] for memoryless sources), it is too complicated to gain any insights. Therefore, we must retreat to the asymptotic analysis. To accomplish this, we shall derive a functional-differential equations on the Poisson transforms $\widetilde{B}^{i}(z, u)$, which seem to have a simpler, or at least more compact, form. These functionaldifferential equations are next changed into a simple matrix recurrence in terms of the Mellin transform (cf. [3, 14, 21]). After solving this matrix equation (in fact, for the asymptotic analysis we do not even need to solve it explicitly), we apply the inverse Mellin transform
to recover the Poisson transform $B^{i}(z, u)$ for $z \rightarrow \infty$ in a cone around the real axis. This suffices, since by analytic depoissonization (cf. [7, 9]) we can extract asymptotic expression for the average profile $B_{m}^{i}$ for $m \rightarrow \infty$, which further leads to our final results.

Before we present out findings, we must introduce some more notation. Let $s$ be complex, and then

$$
\mathrm{Q}(s)=1-\mathrm{P}(s), \quad \text { where } \quad \mathrm{P}(s)=\left\{p_{i j}^{-s}\right\}_{i, j=1}^{V}
$$

and $I$ is the identity matrix. Let now $\mathrm{Q}^{*}(s)=\operatorname{adj}[\mathrm{Q}(s)]$ be the adjoint matrix of $\mathrm{Q}(s)$, that is, $Q^{*}(s)=(-1)^{i+j}\left\{Q^{j, i}(s)\right\}_{i, j \in \mathcal{A}}$ where $Q^{j, i}(s)$ is the $(j, i)$ cofactor of $Q(s)$ defined as $\mathrm{Q}^{-1}(s)=\mathrm{Q}^{*}(s) / \operatorname{det} \mathrm{Q}(s)$ (cf. [16]). Furthermore,

$$
\begin{aligned}
\beta & :=\left.\left[\operatorname{det} \mathrm{Q}^{\prime \prime}(s)\right]\right|_{s=-1}, \\
\dot{\mathrm{Q}}^{-} & :=\left.\dot{\mathrm{Q}}^{*}(s)\right|_{s=-1}, \\
\dot{\mathrm{x}}(-2) & :=\sum_{i=1}^{\infty}\left(\left.\mathrm{Q}^{-1}(-2) \cdots \mathrm{Q}^{-1}(-i)\left(\mathrm{Q}^{-1}(s)\right)^{\prime}\right|_{s=-i-1} \mathrm{Q}^{-1}(-i-2) \cdots\right) \mathrm{K}, \\
\vartheta & :=\pi \dot{\mathrm{x}}(-2),
\end{aligned}
$$

where

$$
\begin{equation*}
\mathrm{K}:=\left(\prod_{i=0}^{\infty} Q^{-1}(-2-i)\right)^{-1} \psi \tag{6}
\end{equation*}
$$

$\psi=[1,1, \cdots, 1]_{V \times 1}^{T}$ is the column vector consisting of all 1 s . In the above, we use the following

$$
\dot{\mathbf{x}}(-2)=\left.\frac{d}{d s} \mathbf{x}(s)\right|_{s=-2}
$$

Finally

$$
\omega:=\operatorname{det}\left[\begin{array}{cccc}
1 & -p_{12} & \ldots & -p_{1 V} \\
1 & 1-p_{22} & \ldots & -p_{2 V} \\
\vdots & \vdots & \ddots & \vdots \\
1 & -p_{V 2} & \ldots & 1-p_{V V}
\end{array}\right]
$$

In addition, we use the standard notation for entropy of a Markov source. In particular,

$$
h=-\sum_{i=1}^{V} \pi_{i} \sum_{j=1}^{V} p_{i j} \ln p_{i j}
$$

and for a probability vector $\mathrm{p}=\left(p_{1}, \ldots, p_{V}\right)$

$$
h_{\mathbf{p}}=-\sum_{i=1}^{V} p_{i} \ln p_{i} .
$$

Also, we often use $\mathbf{p}(s)=\left[\pi_{1}^{-s}, \pi_{2}^{-s}, \ldots, \pi_{V}^{-s}\right]$ which becomes $\pi$ when $s=-1$.

In Section 3.1 we prove the following main result for MI model for stationary Markov sources (i.e., $\mathbf{p}=\boldsymbol{\pi}$ ).

Theorem 1 Consider a Markov stationary source with transition probabilities $P=\left\{p_{i j}\right\}_{i, j=1}$, , that is, $\operatorname{Pr}\left\{X_{t}(\ell)=k\right\}=\pi_{k}$ for all $t=0,1, \ldots$ and $\ell=1,2, \ldots, m$.
(i) [ Typical Depth/Phrase Length] For latge m the following holds

$$
\begin{align*}
\mathrm{E}\left[D_{m}\right] & =\frac{1}{h}\left(\ln m+\gamma-1+h-h_{\pi}-\frac{\beta}{2 \omega h}-\vartheta+\delta_{1}(\ln m)\right)+O\left(\frac{\ln m}{m}\right)  \tag{7}\\
\operatorname{Var}\left[D_{m}\right] & =\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{\mathrm{Q}}^{*} \psi-h^{2}\right) \ln m+O(1), \tag{8}
\end{align*}
$$

$a n d$

$$
\begin{equation*}
\frac{D_{m}-\mathrm{E}\left[D_{m}\right]}{\sqrt{\operatorname{Var} D_{m}}} \rightarrow N(0,1) \tag{9}
\end{equation*}
$$

where $\vartheta=\pi \dot{\mathbf{x}}(-2)$ and $\gamma=0.577 \ldots$ is the Euler constant. The function $\delta_{1}(x)$ is a fluctuating function with a small amplitude when

$$
\begin{equation*}
\frac{\ln p_{i j}+\ln p_{1 i}-\ln p_{1 j}}{\ln p_{11}} \in \mathbb{Q}, \quad i, j=1,2, \ldots, V \tag{10}
\end{equation*}
$$

where $\mathbb{Q}$ is the set of rational numbers. If (10) does not hold, then $\lim _{x \rightarrow \infty} \delta_{1}(x)=0$.
One can strengthen (9) as follows. If $\mu_{m}=\mathbf{E}\left[D_{m}\right]$, and $\sigma_{m}=\sqrt{\operatorname{Var} D_{m}}$, then for a complex $\tau$ the generating function $D_{m}(u)=\mathbf{E}\left[u^{D_{m}}\right]$ becomes

$$
\begin{equation*}
e^{-\tau \mu_{\mathrm{m}} / \sigma_{m}} D_{m}\left(e^{\tau / \sigma_{\mathrm{m}}}\right)=e^{\frac{\tau^{2}}{2}}\left(1+O\left(\frac{1}{\sqrt{\ln m}}\right)\right) \tag{11}
\end{equation*}
$$

as $m \rightarrow \infty$, thus the rate of convergence to the normal distribution is $O(1 / \sqrt{\ln m})$. This implies the existence of positive constants $A$ and $\alpha<1$ such that

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|\frac{D_{m}-\mathrm{E}\left[D_{m}\right]}{\sqrt{\operatorname{Var} D_{m}}}\right| \geq k\right\} \leq A \alpha^{k} \tag{12}
\end{equation*}
$$

uniformly in $k$.
(ii) [Depth of Insertion/Last Phrase Length] The depth of insertion (or equivalently, the last phrase length) $I_{m}$ behaves asymptotically as the typical phrase $D_{m}$. More precisely, for some $A>0$ and $\alpha<1$

$$
\begin{align*}
\mathbf{E}\left[I_{m}\right] & =\frac{1}{h}\left(\ln m+\gamma+h-h_{\pi}-\frac{\beta}{2 \omega h}-\vartheta+\delta_{2}(\ln m)\right)+O\left(\frac{\ln m}{m}\right)  \tag{13}\\
\operatorname{Var}\left[I_{m}\right] & =\operatorname{Var}\left[D_{m}\right]+O(1)  \tag{14}\\
e^{-\tau \mu_{\mathrm{m}} / \sigma_{\mathrm{m}}} I_{m}\left(e^{\tau / \sigma_{\mathrm{m}}}\right) & =e^{\frac{\boldsymbol{\tau}^{2}}{2}}\left(1+O\left(\frac{1}{\sqrt{\ln m}}\right)\right) \tag{15}
\end{align*}
$$

where $\delta_{2}(x)$ is a fluctuating function with the same property as $\delta_{1}(x)$. In addition, there exist positive constants $A$ and $\alpha<1$ such that

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|\frac{I_{m}-\mathbf{E}\left[I_{m}\right]}{\sqrt{\operatorname{Var} I_{m}}}\right| \geq k\right\} \leq A \alpha^{k} \tag{16}
\end{equation*}
$$

Remarks. (i) Alternative Representation. We can present main results of Theorem 1 in a different form which is particularly useful for the proof of the limiting distribution and, more importantly, can lead to some further generalizations. This new derivation can be found in Appendix A. For matrix $\mathrm{P}(s)$, we define the principal left eigenvector $\pi(s)$, the principal right eigenvector $\psi(s)$ associated with the largest eigenvalue $\lambda(s)$ as

$$
\begin{align*}
\pi(s) \mathrm{P}(s) & =\lambda(s) \pi(s)  \tag{17}\\
\mathrm{P}(s) \psi(s) & =\lambda(s) \psi(s) \tag{18}
\end{align*}
$$

where $\pi(s) \psi(s)=1$. Observe that $\pi(-1)=\pi=\left(\pi_{1}, \ldots, \pi_{V}\right), \psi(-1)=\psi=(1, \ldots, 1)$, and $\lambda(-1)=1$. Also, for an vector $\mathbf{x}(s)$ we write $\dot{\mathbf{x}}(s)=\frac{d}{d s} \mathbf{x}(s)$ and $\ddot{\mathbf{x}}(s)=\frac{d^{2}}{d s^{2}} \mathbf{x}(s)$. In Appendix A we shall prove that

$$
\begin{aligned}
& \dot{\lambda}(-1)=\pi \dot{\mathrm{P}}(-1) \psi=h \\
& \ddot{\lambda}(-1)=\pi \ddot{\mathrm{P}}(-1) \psi+2 \dot{\pi}(-1) \dot{\mathrm{P}}(-1) \psi-2 \dot{\lambda}(-1) \dot{\pi}(-1) \psi .
\end{aligned}
$$

Then, (7)-(8) of Theorem 1 can be alternatively written as

$$
\begin{align*}
\mathrm{E}\left[D_{m}\right] & =\frac{1}{\dot{\lambda}(-1)}\left(\ln m+\gamma-1+\dot{\lambda}(-1)+\frac{\ddot{\lambda}(-1)}{2 \dot{\lambda}^{2}(-1)}-\vartheta-\pi \dot{\psi}(-1)+\delta_{1}(\ln m)\right) \\
& +O\left(\frac{\ln m}{m}\right),  \tag{19}\\
\operatorname{Var}\left[D_{m}\right] & =\frac{\ddot{\lambda}(-1)-\dot{\lambda}^{2}(-1)}{\dot{\lambda}^{3}(-1)} \ln m+O(1) . \tag{20}
\end{align*}
$$

In a similar fashion, we can write for $I_{m n}$.
(ii) Memoryless Source. Let us compare the findings of Theorem 1 to those obtained for a memoryless source (cf. [11, 20]). The Markov source becomes a memoryless source if we assume $p_{j i}=\pi_{i}$ for $i_{,} j=1,2, \ldots, V$. Observe that then $\omega=1, \beta=-\sum_{i=1}^{V} \pi_{i} \ln ^{2} \pi_{i}, h \pi=h$, and

$$
\begin{aligned}
\mathrm{Q}(s) & =1-\psi \otimes \mathbf{p}(s) \\
\mathbf{Q}^{-1}(s) & =\frac{1}{1-\mathbf{p}(s) \psi}[(1-\mathbf{p}(s) \psi)!+\psi \otimes \mathbf{p}(s)] \\
\mathbf{Q}(-j) \psi & =(1-\mathbf{p}(-j) \psi) \psi
\end{aligned}
$$

where $\mathbf{p}(s)=\left(\pi_{1}^{-s}, \ldots, \pi_{V}^{-s}\right)$, and $\otimes$ is the tensor product of vectors (e.g., the product $\psi \otimes \mathbf{p}(s)$ is a matrix with the $i$ th column equal to $\left.\left(\pi_{i}^{-s}, \ldots, \pi_{i}^{-s}\right)^{T}\right)$. Thus,

$$
\dot{\mathrm{Q}}^{*} \psi=\left(-\mathrm{p}^{\prime}(s) \psi \mid+\psi \times \mathrm{p}^{\prime}(s)\right) \psi=0
$$

We can also prove the following commutation laws

$$
\mathrm{Q}^{-1}(-i) \mathrm{Q}^{-1}(-j)=\mathrm{Q}^{-1}(-j) \mathrm{Q}^{-1}(-i), \dot{\mathrm{Q}}^{-1}(-i) \mathrm{Q}^{-1}(-j)=\mathrm{Q}^{-1}(-j) \dot{\mathrm{Q}}^{-1}(-i)
$$

for any $i, j \geq 2$. As a result, we find

$$
\begin{aligned}
\dot{\mathbf{x}}(-2) & =\sum_{i=2}^{\infty} \dot{\mathrm{Q}}^{-1}(-i) \mathrm{Q}(-i) \psi=\sum_{i=2}^{\infty} \dot{\mathrm{Q}}^{-1}(-i)(1-\mathbf{p}(-i) \psi) \psi \\
& =\sum_{i=2}^{\infty} \frac{\mathbf{p}^{\prime}(-i) \psi}{1-\mathbf{p}(-i) \psi} \psi,
\end{aligned}
$$

and finally

$$
\vartheta=\pi \dot{\mathrm{x}}(-2)=-\sum_{k=1}^{\infty} \frac{\sum_{i=1}^{V} \pi_{i}^{k+1} \ln \pi_{i}}{1-\sum_{i=1}^{V} \pi_{i}^{k+1}}
$$

which coincides with the findings of [20]. In summary, our results for the Markovian source reduce to those of [20] when the source becomes memoryless.
(iii) Fluctuating Function $\delta(x)$. A few words of discussion about the fluctuating function $\delta(x)$ is in order. The amplitude of this function is very small, however, it increases with $V$. For example, for the unbiased memoryless source $\left|\delta_{1}(x)\right| \leq 10^{-6}$ for $V=2$ (cf. [10, 14]). While this valuc may be savely ignored in the first order analysis, it is of prime interest to second order analysis. For example, the fluctuating function $\delta_{1}(x)$ determines the behavior of the Lempel-Ziv redundancy (cf. [12]). In view of this, one may ask for which Markov sources condition (10) holds. We know that for memoryless sources (10) becomes

$$
\frac{\ln \pi_{i}}{\ln \pi_{j}} \in \mathbb{Q} \quad i \in \mathcal{A} .
$$

But, can we find a non-degenerate Markov source (i.e., which is not a memoryless) that satisfies (10)? The answer is positive, and here is an example. Let $M(b)=\left\{e^{-2 \pi k_{i j} / b}\right\}_{i, j=1}^{V}$ for some integers $k_{i j}$ and a positive $b$ where $i, j \in \mathcal{A}$. The matrix $M(b)$ is positive definite and its main eigenvalue $\lambda(b)$ is real positive with positive right eigenvector $\mathbf{r}(b)=\left(r_{1}\left(b 0, \ldots, r_{V}(b)\right)\right.$. Since $\lambda(b) \rightarrow 0$ as $b \rightarrow 0$ and $\lambda(b) \rightarrow V$ as $b \rightarrow \infty$, there exists $b_{0}$ such that $\lambda\left(b_{0}\right)=1$. Define now

$$
p_{i j}=\frac{\tau_{j}\left(b_{0}\right)}{r_{i}\left(b_{0}\right)} e^{-2 \pi k_{i j} / b_{0}}
$$

for $i, j \in \mathcal{A}$. Observe that

$$
\sum_{j \in \mathcal{A}} p_{i j}=\frac{1}{r_{i}\left(b_{0}\right)} \sum_{j \in \mathcal{A}} r_{j}\left(b_{0}\right) e^{-2 \pi k_{i j} / b_{0}}=\frac{r_{i}\left(b_{0}\right)}{r_{i}\left(b_{0}\right)}=1,
$$

since $r\left(b_{0}\right)$ is the right eigenvector of $M$ with $\lambda\left(b_{0}\right)=1$. There $P=\left\{p_{i j}\right\}_{i j \in \mathcal{A}}$ generates a non-degenerated Markov source for which (10) holds.

We now extend the above results into two directions, namely a non-stationary Markov source and the MI model with binomial ( $m, r$ ) number of independent sources. Both extensions are crucial for our derivation of results for GK model (i.e., with fixed number of phrases).

Let us start with a non-stationary Markov source. Observe that our basic set of recurrences (5) for the conditional generating functions $B_{m}^{i}(u)$ stays the same, and the only change in our global recurrence (4) for the cumulative generating function $B_{m}(u)$ reduces to replacing the stationary probability $\pi$ by the initial distribution vector $p$. As we shall sce in Section 3, the asymptotics of the average profile largely depend on the asymptotics of the conditional average profile. This will translate in the same leading terms of the asymptotic expansions of the average depth (phrase length) $D_{r n}(\mathbf{p})$, and the depth of insertion (last phrase length) $I_{m}(\mathbf{p})$. In fact, the difference is exhibited only in the $O(1)$ term.

We summarize our finding in the following corollary.
Corollary 1 [Non-Stationary Markov Source] Consider a Markov source with initial probability vector $\mathbf{p}=\left(p_{1}, \ldots, p_{V}\right)$. Then for large $m$

$$
\begin{align*}
\mathrm{E}\left[D_{m}(\mathbf{p})\right] & =\frac{1}{h}\left(\ln m+\gamma-1+h-h_{\mathrm{p}}-\frac{\beta}{2 \omega h}-\vartheta+\delta_{3}(\ln m)\right)+O\left(\frac{\ln m}{m}\right),  \tag{21}\\
\mathrm{E}\left[I_{m}(\mathbf{p})\right] & =\frac{1}{h}\left(\ln m+\gamma+h-h_{\mathrm{p}}-\frac{\beta}{2 \omega h}-\vartheta+\delta_{4}(\ln m)\right)+O\left(\frac{\ln m}{m}\right)  \tag{22}\\
\operatorname{Var}\left[I_{m}(\mathbf{p})\right] & =\operatorname{Var}\left[D_{m}(\mathbf{p})\right]+O(1)=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{\mathrm{Q}}^{*} \psi-h^{2}\right) \ln m+O(1) \tag{23}
\end{align*}
$$

with the notation as in Theorem 1, where $\delta_{3}(x)$ and $\delta_{4}(x)$ are fluctuating functions with small amplitudes. In addition,

$$
\begin{gather*}
\frac{D_{m}(\mathbf{p})-\mathbf{E}\left[D_{m}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} D_{m}(\mathbf{p})}} \rightarrow N(0,1),  \tag{24}\\
\frac{I_{m}(\mathbf{p})-\mathbf{E}\left[I_{m}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} I_{m}(\mathbf{p})}} \rightarrow N(0,1) \tag{25}
\end{gather*}
$$

with the rate of convergence $O(1 / \sqrt{\ln m})$. Moreover, there exist positive constants $A$ and $\alpha<1$ such that

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|\frac{D_{m}(\mathbf{p})-\mathbf{E}\left[D_{m}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} D_{m}(\mathbf{p})}}\right| \geq k\right\} \leq A \alpha^{k} \tag{26}
\end{equation*}
$$

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|\frac{I_{m}(\mathrm{p})-\mathrm{E}\left[I_{m}(\mathrm{p})\right]}{\sqrt{\operatorname{Var} I_{m}(\mathrm{p})}}\right| \geq k\right\} \leq A \alpha^{k} \tag{27}
\end{equation*}
$$

uniformly in $k$.

Finally, we consider the MI model in which the number of sources $M$ is a random variable distributed as $B(m, r):=\operatorname{binomial}(m, r)$, that is,

$$
\operatorname{Pr}\{M=k\}=\binom{m}{k} \tau^{k}(1-\tau)^{m-k}
$$

Let $D_{m}^{B}$ and $I_{m}^{B}$ (or $D_{m}^{B(r)}$ and $I_{m}^{B(r)}$ ) denote, respectively, the typical depth and the depth of insertion in such a model.

Corollary 2 [Random number of non-stationary Markov Sources] Consider a Markov source with initial probability vector $\mathbf{p}=\left(p_{1}, \ldots, p_{V}\right)$ and random number, $M$, of sources distributed as the binomial $(m, r)$. Then for large $m$

$$
\begin{align*}
\mathbf{E}\left[D_{m}^{B}(\mathbf{p})\right] & =\frac{1}{h}\left(\ln (m r)+\gamma-1+h-h_{\mathbf{p}}-\frac{\beta}{2 \omega h}-\vartheta+\delta_{5}(\ln m)\right)+O\left(\frac{\ln m}{m}\right),  \tag{28}\\
\mathbf{E}\left[I_{m}^{B}(\mathbf{p})\right] & =\frac{1}{h}\left(\ln (m r)+\gamma+h-h_{\mathbf{p}}-\frac{\beta}{2 \omega h}-\vartheta+\delta_{6}(\ln m)\right)+O\left(\frac{\ln m}{m}\right),  \tag{29}\\
\operatorname{Var}\left[I_{m}^{B}(\mathbf{p})\right] & =\operatorname{Var}\left[D_{m}^{B}(\mathbf{p})\right]+O(1)=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{\mathbf{Q}}^{-} \psi-h^{2}\right) \ln (m r)+O(1) \tag{30}
\end{align*}
$$

where $\delta_{5}(x)$ and $\delta_{6}(x)$ are fuctuating functions with small amplitudes. In addition,

$$
\begin{gather*}
\frac{D_{m}^{B}(\mathbf{p})-\mathrm{E}\left[D_{m}^{B}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} D_{m}^{B}(\mathbf{p})}} \rightarrow N(0,1)  \tag{31}\\
\frac{I_{m}^{B}(\mathbf{p})-\mathbf{E}\left[I_{m}^{B}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} I_{m}^{B}(\mathbf{p})}} \rightarrow N(0,1) \tag{32}
\end{gather*}
$$

with the rate of convergence $O(1 / \sqrt{\ln m})$. Finally, there exist positive constants $A$ and $\alpha<1$ such that

$$
\begin{align*}
& \operatorname{Pr}\left\{\left|\frac{D_{m}^{B}(\mathbf{p})-\mathbf{E}\left[D_{m}^{B}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} D_{m}^{B}(\mathbf{p})}}\right| \geq k\right\} \leq A \alpha^{k}  \tag{33}\\
& \operatorname{Pr}\left\{\left|\frac{I_{m}^{B}(\mathbf{p})-\mathrm{E}\left[I_{m}^{B}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} I_{m}^{B}(\mathbf{p})}}\right| \geq k\right\} \leq A \alpha^{k} \tag{34}
\end{align*}
$$

uniformly in $k$.

Proof. Let us only consider the typical depth $D_{m}^{B}$. The proof follows immediately from the fact that the generating function $D_{m}^{B}(u)$ satisfies

$$
D_{m}^{B}(u)=\sum_{k=0}^{m}\binom{m}{k} r^{k}(1-r)^{m-k} D_{k}(u)
$$

where $D_{k}(u)$ is the generating function of the typical depth in the MI model with $k$ Markov sources. Observe now that the Poisson transform of $D_{m}^{B}$ satisfies

$$
\tilde{D}^{B}(z, u)=\tilde{D}(z r, u) e^{-z r}
$$

where $\tilde{D}(z, u)$ is the Poisson transform of the MI model with fixed number of sources (and already presented in Theorem 1 while the analysis can be found in Section 3). The moments can be also recovered from the following formula recently proved in [8] (interestingly, analytic depoissonization was used to derive it, too)

$$
\sum_{k=0}^{m}\binom{m}{k} r^{k}(1-r)^{m-k} \ln k=\ln (m r)-\frac{1-r}{2 m r}+\sum_{i \geq 2} \frac{a_{i}}{n^{i}}
$$

where the cocfficients $a_{k}$ are explicitly computable.

### 2.2 Fixed Number of Phrases Model - Gilbert-Kadota Model

In this subsection, we present our main findings for the Gilbert-Kadota model in which a single Markovian source generates (possibly infinite) sequence that is partitioned according to the Lempel-Ziv algorithm until $m$ full phrases are obtained. As before, we study the typical phrase length $D_{m}$ and the last phrase length $I_{m}$. To avoid confusions, we often append an upper index $M I$ or $G K$ to $D_{m}$ and $I_{m}$ to denote the typical phrase length and last phrase length in the MI model and the GK model, respectively. Furthermore, as before, it is convenient to build a digital scarch tree out of these $m$ phrases, as shown in Figure 1. We observe, however, that this time the DST is built from suffizes of a single Markovian sequence, thus we might call it a suffix digital search tree. Clearly, the typical phrase length $D_{m}^{G K}$ becomes the typical depth, and the last phrase length $I_{m}^{G K}$ corresponds to the depth of insertion in the associated DST.

The GK model introduces additional trickly statistical dependency between phrases. The recurrence (4) and the differential-functional equation (5) do not hold any more, however, the relationship (3) between the typical depth and the depth of insertion is still true. To analyze GK model, we use stochastic dominance, that is, we (asymptotically) bound in a stochastic sense define below the depth of insertion $I_{m}^{G K}$ by the depth of insertion in the
modified MI model. More precisely, in the GK model, we delete $K$ phrases, thus making a "gap" of significant size so that the newly inserted phrase resembles the one in the MI model, hence results of MI model can be applied.

To present more succinctly our analysis, we introduce some new notation. We say that $I_{m}^{\prime}$ stochastically dominates $I_{m}$ and write $I_{m} \leq_{s t} I_{m}^{\prime}$ if for every $k$ we have

$$
\operatorname{Pr}\left\{I_{m} \geq k\right\} \leq \operatorname{Pr}\left\{I_{m}^{\prime} \geq k\right\}
$$

In our investigation, however, we also need the so called asymptotic stochastic dominance that we denote as $I_{m} \preceq_{s l} I_{m}^{\prime}$ and define precisely below.

Definition 1 (i) Let $X$ and $Y$ be two integer random variables, and $\varepsilon>0$. We say that $X$ is at distance $\varepsilon$ from $Y$ and write it as $d(X, Y) \leq \varepsilon$ if for all integers $k$ :

$$
\begin{equation*}
|\operatorname{Pr}\{X \geq k\}-\operatorname{Pr}\{Y \geq k\}|<\varepsilon . \tag{35}
\end{equation*}
$$

(ii) We say that the sequence of random variables $X_{m}$ asymptotically dominates $Y_{m}$ or shortly

$$
X_{m} \preceq_{\mathrm{st}} Y_{m}
$$

if

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \max _{k}\left(\operatorname{Pr}\left\{X_{m} \geq k\right\}-\operatorname{Pr}\left\{Y_{m} \geq k\right\}\right)=0 \tag{36}
\end{equation*}
$$

The last definition is illustrated well by the following simple result.
Lemma 1 If $X_{m} \leq_{s t} Y_{m}^{\prime}$ and $\lim _{m \rightarrow \infty} d\left(Y_{m}, Y_{m}^{\prime}\right)=0$, then $X_{m} \preceq_{s t} Y_{m}$.
Proof. By assumptions, for all integers $k$ and $m$ we have $\operatorname{Pr}\left\{X_{m} \geq k\right\} \leq \operatorname{Pr}\left\{Y_{m}^{\prime} \geq k\right\}$ and $\lim _{m \rightarrow \infty} \max _{k}\left|\operatorname{Pr}\left\{Y_{m}^{\prime} \geq k\right\}-\operatorname{Pr}\left\{Y_{m} \geq k\right\}\right|=0$. Thus, (36) follows.

In the next section, we establish certain inequalities between the MI model and the GK model, that we review briefly here. For some $K<m$ we denote by $I_{m-K+1}$ the depth of insertion to a DST tree that is built from any subset of size $m-K$ of $m$ original phrases. It is easy to see that in both models we have the following (deterministic) inequality

$$
\begin{equation*}
I_{m-K+1} \leq I_{m+1} \leq I_{m-K+1}+K, \tag{37}
\end{equation*}
$$

provided the same phrase is inserted. The left-hand size is quite obvious, while the right-hand size is a consequence of the fact that a now phrase can be incremented at most by one symbol. In other words, the DST tree does not have unary nodes (i.e., nodes with degree one).

In view of this, we can work on $I_{m-K}$ in which $K$ phrases are (conveniently) dcleted smoothing down dependencies. We consider now the MI model such that all phrases start with a given, but otherwise arbitrary symbol, say $a \in \mathcal{A}$. In other words, we consider a non-stationary model with an initial vector $\mathrm{p}_{a}$ concentrated on symbol $a$, i.e. contains all zeros except a 1 on the position corresponding to symbol $a$. We denote $I_{m}^{M I}$ ( pa ) the depth of insertion in this model. We also consider the GK model conditioned on the fact that the $m$ th phrase starts with symbol $a$. We denote $I_{m, K}^{G K}\left(\mathbf{p}_{a}\right)$ the depth of insertion of the $m$ th phrase when $K$ phrases are deleted before it. We shall prove in Section 4 that there exists $K=O(1)$ such that

$$
I_{m-K}^{M I B(r)}\left(\mathbf{p}_{a}\right) \preceq_{\mathrm{st}} I_{m, K}^{G K}\left(\mathbf{p}_{a}\right) \preceq_{\mathrm{st}} I_{m-K}^{M I}\left(\mathbf{p}_{a}\right)+K
$$

where $I_{m-K}^{M I B(r)}\left(\mathbf{p}_{a}\right)$ is the depth of insertion in the MI model with the binomial $(r, m-K)$ number of phrases for some $0<r<1$. Thus, based on our results from the previous section, we shall be able to prove the following theorem.

Theorem 2 Consider a Markov source with initial probability vector $\mathbf{p}$. Then for large $m$

$$
\begin{align*}
\mathbf{E}\left[D_{m}^{G K}(\mathbf{p})\right] & =\mathrm{E}\left[I_{m}^{G K}(\mathbf{p})\right]+O(1)=\frac{1}{h} \ln m+O(1)  \tag{38}\\
\operatorname{Var}\left[I_{m}^{G K}(\mathbf{p})\right] & =\operatorname{Var}\left[D_{m}^{G K}(\mathbf{p})\right]+O(1)=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{\mathbf{Q}}^{*} \psi-h^{2}\right) \ln m+O(1) \tag{39}
\end{align*}
$$

with the notation as in Theorem 1, and

$$
\begin{gather*}
\frac{D_{m}^{G K}(\mathbf{p})-\mathbf{E}\left[D_{m}^{G K}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} D_{m}^{G K}(\mathbf{p})}} \rightarrow N(0,1)  \tag{40}\\
\frac{r_{m}^{G K}(\mathbf{p})-\mathbf{E}\left[I_{m}^{G K}(\mathbf{p})\right]}{\sqrt{\operatorname{Var} I_{m}^{G K}(\mathbf{p})}} \rightarrow N(0,1) \tag{41}
\end{gather*}
$$

with the rate of convergence $O(1 / \sqrt{\ln m})$. In addition, the normalized $D_{m}^{G K}(\mathrm{p})$ and $I_{m}^{G K}(\mathbf{p})$ converge in moments to the corresponding moments of the standard normal distribution.

### 2.3 Lempel-Ziv ModeI

Finally, we deal with the Lempel-Ziv model in which a Markov sequence of fixed length $n$ is partitioned in (a random number) $M_{n}$ of (full) phrases. As before, $I_{i}$ represents the $i$ th phrase for $1 \leq i \leq M_{n}$. We write $J_{n}$ for the last full phrase which also becomes $J_{n}=I_{M_{n}}$. The typical phrase length $\Delta_{n}$ is defined as follows:

$$
\begin{equation*}
\operatorname{Pr}\left\{\Delta_{n}=k\right\}=\sum_{m=M_{\min }}^{M_{\max }} \frac{1}{m} \sum_{i=1}^{m} \operatorname{Pr}\left\{I_{i}=k \& M_{n}=m\right\} \tag{42}
\end{equation*}
$$

where $M_{\text {min }}=O(\sqrt{n})$ is the minimum number of phrases and $M_{\max }=O\left(n / \log _{2} n\right)$ is the maximum number of phrases (cf. [11]). In passing, we should observe that there is a relationship between the phrase length $I_{i}$ and the number of phrases $M_{n}$. Indeed,

$$
M_{n}=\max \left\{m: \sum_{i=1}^{m} I_{i}^{G K} \leq n\right\}
$$

where in the above we explicitly show that the phrase length $I_{i}^{G K}$ is the one corresponding to the phrase length in the GK model.

Using Theorem 2, we shall prove below the following result.
Theorem 3 Let a Markov source generates a single sequence of length $n$. Then, for large $n$

$$
\begin{equation*}
\frac{\Delta_{n}-\mathrm{E}\left[\Delta_{n}\right]}{\sqrt{\operatorname{Var} \Delta_{n}}} \rightarrow N(0,1) \tag{43}
\end{equation*}
$$

In addition, $\Delta_{n}$ converges in moments, and in particular

$$
\begin{align*}
\mathrm{E}\left[\Delta_{n}\right] & \sim \mathbf{E}\left[J_{n}\right] \sim \frac{1}{h} \ln (n h / \ln n)  \tag{44}\\
\operatorname{Var}\left[\Delta_{n}\right] & \sim \operatorname{Var}\left[J_{n}\right] \sim \frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{\mathrm{Q}}^{-} \psi-h^{2}\right) \ln (n h / \ln n) \tag{45}
\end{align*}
$$

provided the number of phrases $M_{n}$ converges exponentially to its mean.
Proof. Let

$$
\begin{align*}
& \mu(\pi)=\frac{1}{h} \ln \left(\frac{n h}{\ln n}\right)  \tag{46}\\
& \sigma(n)=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{Q}^{*} \psi-h^{2}\right) \ln \left(\frac{n h}{\ln n}\right) \tag{47}
\end{align*}
$$

We are going to prove that for any $\varepsilon>0$ and for all set of integers $B$

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \max _{B}\left(\operatorname{Pr}\left\{\Delta_{n} \in B\right\}-\operatorname{Pr}\left\{D_{\lfloor(1+\varepsilon) \mu(n)\rfloor} \in B\right\}\right)=0 \tag{48}
\end{equation*}
$$

and

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} \max _{B}\left(\operatorname{Pr}\left\{D_{\lfloor(1-\varepsilon) \mu(n)\rfloor} \in B\right\}-\operatorname{Pr}\left\{\Delta_{n} \in B\right\}\right)=0 \tag{49}
\end{equation*}
$$

First of all, observe that for any $\varepsilon>0$ (cf. [24])

$$
\lim _{n \rightarrow \infty} \operatorname{Pr}\left\{M_{n} \notin((1-\varepsilon) \mu(n),(1+\varepsilon) \mu(n))\right\}=0 .
$$

We rewrite (42) as

$$
\operatorname{Pr}\left\{\Delta_{n} \in B\right\}=\sum_{m=1}^{n} \frac{1}{m} \sum_{\ell=1}^{m} \operatorname{Pr}\left\{I_{\ell}^{G K} \in B \& M_{n}=m\right\}
$$

for any set of integers $B$. Then

$$
\operatorname{Pr}\left\{\Delta_{n} \in B\right\} \leq \delta_{n}+\sum_{m=\lfloor(1-\varepsilon) \mu(n)\rfloor}^{m=\lfloor(1+\varepsilon) \mu(n))\rfloor} \frac{1}{m} \sum_{\ell=1}^{m} \operatorname{Pr}\left\{I_{\ell}^{G K} \in B \& M_{n}=m\right\}
$$

with $\delta_{n}=\operatorname{Pr}\left\{M_{n} \notin((1-\varepsilon) \mu(n),(1+\varepsilon) \mu(n))\right\}$. We have the following chain of trivial inequalities

$$
\begin{aligned}
\operatorname{Pr}\left\{\Delta_{n} \in B\right\} & \leq \delta_{n}+\sum_{m=\lceil(1-\varepsilon) \mu(n)\rceil}^{m=\lfloor(1+\varepsilon) \mu(n))]} \frac{1}{m} \sum_{\ell=1}^{\ell=i\lfloor(1+\varepsilon) \mu(n)\rfloor} \operatorname{Pr}\left\{I_{\ell} \in B \& M_{n}=m\right\} \\
& \leq \delta_{n}+\sum_{m=\lceil(1-\varepsilon) \mu(n) \mu(n))]}^{m} \frac{1}{(1-\varepsilon) \mu(n)} \sum_{\ell=1}^{\ell=\lfloor(1+\varepsilon) \mu(n)\rfloor} \operatorname{Pr}\left\{I_{\ell}^{G K} \in B \& M_{n}=m\right\} \\
& \leq \delta_{n}+\left(\frac{1+\varepsilon}{1-\varepsilon}\right) \sum_{m=\lceil(1-\varepsilon) \mu(n)\rceil}^{m=\lfloor(1+\varepsilon) \mu(n))]} \operatorname{Pr}\left\{D_{\lfloor(1+\varepsilon) \mu(n)]}^{G K} \in B \& M_{n}=m\right\} \\
& \leq \delta_{n}+\left(\frac{1+\varepsilon}{1-\varepsilon}\right) \delta_{n}+\left(\frac{1+\varepsilon}{1-\varepsilon}\right) \operatorname{Pr}\left\{D_{\lfloor(1+\varepsilon) \mu(n)]}^{G K} \in B\right\}
\end{aligned}
$$

In a similar manner, we prove a lower bound

$$
\operatorname{Pr}\left\{\Delta_{n} \in B\right\} \geq\left(\frac{1+\varepsilon}{1-\varepsilon}\right) \delta_{n}+\left(\frac{1+\varepsilon}{1-\varepsilon}\right) \operatorname{Pr}\left\{D_{\lfloor(1-\varepsilon) \mu(n)\rfloor}^{G K} \in B\right\}
$$

The above two inequalities prove (48) and (49). For the convergence in moments we need $\delta_{n}=O\left(e^{-\alpha n}\right)$ for some $\alpha>0$, which is assumed to hold.

Remark. Merhav's result [15] allows to conclude that for Markov sources

$$
\left.\operatorname{Pr}\left\{M_{n} \geq \mu(n)(1+\varepsilon)\right)\right\} \leq(1+o(1)) e^{-\alpha n}
$$

for a constant $\alpha>0$ and $\varepsilon=O(1 / \sqrt{\log n})$. Unfortunately, we do not sce an easy way to convert Merhav's proof for the left tail of $M_{n}$.

As a consequence of Theorems 2 and 3, we can derive bounds on the average redundancy rate $\vec{R}_{n}$ of the Lempel-Ziv code for Markovian sources. To recall, consider a Markovian sequence of length $n$ for which the Lempel-Ziv code is $\ell_{n}$. Then, the redundancy rate is defined as

$$
R_{n}=\frac{\ell_{n}-n h}{n}
$$

We denote by $\vec{R}_{n}=\mathrm{E}\left[R_{n}\right]$ the average redundancy rate. Using the approach recently proposed in Louchard and Szpankowski [12], we obtain from Theorem 2 the following bounds

$$
\begin{array}{ll}
\frac{2-\ln r-\gamma+h_{\mathrm{r}}+\frac{\beta}{2 \omega h}+\vartheta-\delta_{3}(\ln n)}{\ln n} & (1+o(1)) \leq \frac{\bar{R}_{n}}{h} \\
& \leq \frac{2-\gamma-h+\min _{\mathrm{p}} h_{\mathrm{P}}+\frac{\beta}{2 \omega h}+\vartheta-\delta_{2}(\ln n)}{\ln n}(1+o(1))
\end{array}
$$

where $r=\sum_{a \in \mathcal{A}} \min _{i}\left\{p_{i a}\right\}$ and $\mathbf{r}$ is a vector of size $V$ whose $j$ th component is equal to $\min _{i}\left\{p_{i j}\right\} / r$ (cf. Lemma 11).

## 3 Analysis of Markov Independent Model

As mentioned before, the analysis of MI model is at the heart of our contribution to analytic information theory. In view of this, we present here a detailed proof. It is based on such analytic techniques as: analytic poissonization, Mellin transform, singularities of a complex matrix, and analytic depoissonization.

### 3.1 Poissonization and Mellin Transforms: Analysis of Moments

We first consider the stationary Markov source. The generating function $B_{m}(u)$ of the average profile satisfics (4) with the initial vector $\mathbf{p}=\pi$. Observe that the conditional generating functions $B_{m}^{i}(u)$ fulfill the system of recurrence equations (5). We shall first deal with (5). There is no easy way to solve these recurrences, and thercfore, we transform them to the Poisson model in which $m$ is replaced by a Poisson random variable with mean (complex) $z$ which becomes $m$ when $z$ is restricted to positive integers. Let

$$
\tilde{B}^{i}(z, u)=\sum_{n=1}^{\infty} \tilde{B}_{n}^{i}(u) \frac{z^{n}}{n!} e^{-z}, \quad i \in \mathcal{A}
$$

be the Poisson transform of $B_{m}^{i}(u)$. In addition, we shall write $\tilde{B}_{z}^{i}(z, u):=\frac{\partial}{\partial z} \tilde{B}^{i}(z, u)$ for the derivative of $\tilde{B}^{i}(z, u)$ with respect to $z$. After some simple algebra, we have the following Poissonized differential-functional equations of recurrences (4) and (5)

$$
\begin{equation*}
\widetilde{B}_{z}(z, u)+\tilde{B}(z, u)=u\left[\widetilde{B}^{1}\left(u, \pi_{1} z\right)+\cdots+\widetilde{B}^{V}\left(u, \pi_{V} z\right)\right]+1 \tag{50}
\end{equation*}
$$

and

$$
\begin{equation*}
\widetilde{B}_{z}^{i}(z, u)+\widetilde{B}^{i}(z, u)=u\left[\tilde{B}^{1}\left(u, p_{i 1} z\right)+\cdots+\tilde{B}^{V}\left(u, p_{i V} z\right)\right]+1 \quad \text { for all } \quad i \in \mathcal{A} \tag{51}
\end{equation*}
$$

Let us now concentrate on the evaluation of the first two moments of the depth, that is, we need the first two derivatives of $\widetilde{B}(z, u)$ with respect to $u$ at $u=1$. We derive the following two systems of functional equations after taking into account that $\widetilde{B}^{i}(z, 1)=z$, $\pi_{1}+\cdots+\pi_{V}=1$ and $\sum_{j=1}^{V} p_{i j}=1$

$$
\begin{equation*}
\tilde{B}_{z u}(z, 1)+\tilde{B}_{u}(z, 1)=z+\left[\tilde{B}_{u}^{1}\left(1, \pi_{1} z\right)+\cdots+\tilde{B}_{u}^{V}\left(1, \pi_{v} z\right)\right], \tag{52}
\end{equation*}
$$

$$
\begin{array}{ccc}
\widetilde{B}_{z u}^{1}(z, 1)+\widetilde{B}_{u}^{1}(z, 1) & = & z+\left[\widetilde{B}_{u}^{1}\left(1, p_{11} z\right)+\cdots+\widetilde{B}_{u}^{V}\left(1, p_{1 V} z\right)\right] \\
\cdots & = & \cdots \\
\tilde{B}_{z u}^{V}(z, 1)+\widetilde{B}_{u}^{V}(z, 1) & = & z+\left[\tilde{B}_{u}^{1}\left(1, p_{V 1} z\right)+\cdots+\widetilde{B}_{u}^{V}\left(1, p_{V V} z\right)\right],
\end{array}
$$

and

$$
\begin{aligned}
\tilde{B}_{z u u}(z, 1)+\tilde{B}_{u u}(z, 1) & =2\left[\widetilde{B}_{u}^{1}\left(1, \pi_{1} z\right)+\cdots+\widetilde{B}_{u}^{V}\left(1, \pi_{V} z\right)\right]+\left[\widetilde{B}_{u u}^{1}\left(1, \pi_{1} z\right)+\cdots+\widetilde{B}_{u u}^{V}\left(1, \pi_{V} z\right)\right], \quad(53) \\
\widetilde{B}_{z u u}^{1}(z, 1)+\widetilde{B}_{u u}^{1}(z, 1) & =2\left[\widetilde{B}_{u}^{1}\left(1, p_{11} z\right)+\cdots+\widetilde{B}_{u}^{V}\left(1, p_{1 V} z\right)\right]+\left[\widetilde{B}_{u u}^{1}\left(1, p_{11} z\right)+\cdots+\widetilde{B}_{u u}^{v}\left(1, p_{1 V} z\right)\right] \\
\cdots & =\cdots \\
\tilde{B}_{z u u}^{V}(z, 1)+\widetilde{B}_{u u}^{V}(z, 1) & =2\left[\widetilde{B}_{u}^{1}\left(1, p_{V 1} z\right)+\cdots+\widetilde{B}_{u}^{V}\left(1, p_{V V} z\right)\right]+\left[\widetilde{B}_{u u}^{1}\left(1, p_{V 1} z\right)+\cdots+\tilde{B}_{u u}^{V}\left(1, p_{V V} z\right)\right] .
\end{aligned}
$$

Our goal is now to solve asymptotically (as $z \rightarrow \infty$ in a cone around $\Re(z)>0$ ) the above two set of functional equations. It is well known that equations like these are amiable to attack by the Mellin transform (cf. [3]). To recall, for a function $f(x)$ of real $x$, we define its Mellin transform $F^{*}(s)$ as

$$
F^{*}(s)=\mathcal{M}[f(t) ; s]=\int_{0}^{\infty} f(t) t^{s-1} d t
$$

In some of our arguments we could use either Mellin transform of a complex variable function $f(z)$ or an analytical continuation argument. It is known (cf. [7]) that as long as $\arg (z)$ belongs to some cone around the real axis, the Mellin transform $F(s)$ of a function $f(x)$ of a real argument and its corresponding function of a complex argument is the same. Therefore, we work most of the time with the Mellin transform of a function of real variable as defined above. In our case, a direct solution through Mellin transform does not work well, and therefore we factorize the Mellin transforms of the above functions as follows:

$$
\begin{align*}
B_{i}^{*}(s) & :=\mathcal{M}\left[B_{u}^{i}(z, 1) ; s\right]=\Gamma(s) x_{i}(s), \quad i \in \mathcal{A}  \tag{54}\\
B^{\mathbf{z}}(s) & :=\mathcal{M}\left[B_{u}(z, 1) ; s\right]=\Gamma(s) x(s),  \tag{55}\\
C_{i}^{*}(s) & :=\mathcal{M}\left[B_{u u}^{i}(z, 1) ; s\right]=\Gamma(s) v_{i}(s), \quad i \in \mathcal{A}  \tag{56}\\
C^{*}(s) & :=\mathcal{M}\left[B_{u u}(z, 1) ; s\right]=\Gamma(s) v(s), \tag{57}
\end{align*}
$$

where $\Gamma(s)$ is the Euler gamma function, and $x_{i}(s), x(s), v_{i}(s)$ and $v(s)$ are unknown. The lemma below cstablishes the existence of the above Mellin transforms.

Lemma 2 The Mellin transforms $B_{i}^{*}(s), B^{*}(s)$ and $C_{i}^{*}(s), C^{*}(s)$ exist for $\Re(s) \in(-2,-1)$. In addition,

$$
\begin{array}{ll}
x_{i}(-2)=1, & x(-2)=1 \\
v_{i}(-2)=0, & v(-2)=0 .
\end{array}
$$

Proof. The proof is quite standard and replies on the Lemma 2 from [13]. We leave the details to the interested reader.

Now, we are ready to compute the Mellin transforms of $B_{u}^{i}(z, 1), B_{u u}^{i}(z, 1)$ (cf. (52) and (53), respectively) with respect to $z$. We obtain

$$
\begin{align*}
-(s-1) B^{*}(s-1)+B^{*}(s) & =B_{1}^{*}(s) \pi_{1}^{-s}+\cdots+B_{V}^{*}(s) \pi_{V}^{-s}  \tag{58}\\
-(s-1) B_{1}^{*}(s-1)+B_{1}^{*}(s) & =B_{1}^{*}(s) p_{11}^{-s}+\cdots+B_{V}^{*}(s) p_{1 V}^{-s} \\
\cdots & =\cdots \\
-(s-1) B_{V}^{*}(s-1)+B_{V}^{*}(s) & =B_{1}^{*}(s) p_{V 1}^{-s}+\cdots+B_{V}^{*}(s) p_{V V}^{-s}
\end{align*}
$$

and

$$
\begin{align*}
-(s-1) C^{*}(s-1)+C^{*}(s) & =2\left[B_{1}^{*}(s) \pi_{1}^{-s}+\cdots+B_{V}^{*}(s) \pi_{V}^{-s}\right]+\left[C_{1}^{-}(s) \pi_{1}^{-s}+\cdots+C_{V}^{*}(s) \pi_{V}^{-s}\right],  \tag{59}\\
-(s-1) C_{1}^{*}(s-1)+C_{1}^{*}(s) & =2\left[B_{1}^{*}(s) p_{11}^{-s}+\cdots+B_{V}^{*}(s) p_{1 V}^{-s}\right]+\left[C_{1}^{*}(s) p_{11}^{-s}+\cdots+C_{V}^{*}(s) p_{1 V}^{-s}\right], \\
\cdots & =\cdots \\
-(s-1) C_{V}^{-}(s-1)+C_{V}^{*}(s) & =2\left[B_{1}^{-}(s) p_{V 1}^{-s}+\cdots+B_{V}^{*}(s) p_{V}^{-s}\right]+\left[C_{1}^{-}(s) p_{V 1}^{-s}+\cdots+C_{V}^{*}(s) p_{V V}^{-s}\right] .
\end{align*}
$$

In the above, we used the following two properties of the Mellin transform (cf. [3]):

$$
\begin{aligned}
\mathcal{M}[f(a x) ; s] & =a^{-s} F^{\mathbf{*}}(s), \\
\mathcal{M}\left[f^{\prime}(x) ; s\right] & =-(s-1) F^{*}(s-1) .
\end{aligned}
$$

Unfortunately, the above systems of Mcllin transforms do not have simple explicit solutions. But, we may obtain ones in terms of the functions $x_{i}(s)$ and $v_{i}(s)$ defined in (54) and (56) due to the following property of the gamma function: $\Gamma(s)=(s-1) \Gamma(s-1)$. To present these solution in a compact form, we use from now on matrix notation. Let

$$
\mathbf{x}(s)=\left[\begin{array}{c}
x_{1}(s)  \tag{60}\\
x_{2}(s) \\
\vdots \\
x_{V}(s)
\end{array}\right], \quad \mathbf{v}(s)=\left[\begin{array}{c}
v_{1}(s) \\
v_{2}(s) \\
\vdots \\
v_{V}(s)
\end{array}\right]
$$

In addition, we define

$$
\mathbf{b}(s)=\left[\begin{array}{c}
B_{1}^{*}(s)  \tag{61}\\
B_{2}^{*}(s) \\
\vdots \\
B_{V}^{*}(s)
\end{array}\right], \quad \mathbf{c}(s)=\left[\begin{array}{c}
C_{1}^{*}(s) \\
C_{2}^{*}(s) \\
\vdots \\
C_{V}^{*}(s)
\end{array}\right]
$$

Now, the system of equations (58) and (59) become

$$
\begin{aligned}
\mathbf{x}(s)-\mathbf{x}(s-1) & =\mathrm{P}(s) \mathbf{x}(s) \\
\mathbf{v}(s)-\mathbf{v}(s-1) & =2 \mathrm{P}(s) \mathbf{x}(s)+\mathrm{P}(s) \mathbf{v}(s)
\end{aligned}
$$

where $\mathrm{P}=\left\{p_{i j}^{-s}\right\}_{i, j \in \mathcal{A}}$. Thus,

$$
\begin{align*}
& \mathbf{x}(s)=\mathrm{Q}^{-1}(s) \mathbf{x}(s-1)=\left(\prod_{i=0}^{\infty} \mathrm{Q}^{-1}(s-i)\right) \mathrm{K}  \tag{62}\\
& \mathbf{v}(s)=2 \mathrm{Q}^{-1}(s) \mathrm{P}(s) \mathbf{x}(s)+\mathrm{Q}^{-1}(s) \mathbf{v}(s-1) \tag{63}
\end{align*}
$$

where $Q=I-P$ and $I$ is the identity matrix, and $K$ is defined in (6). The formula on $K$ follows from Lemma 2 (i.e., $\left.\mathbf{x}(-2)=(1, \ldots, 1)^{T}\right)$ and (62). In the next section we prove the convergence of the above infinite product (cf. Lemma 4), however, we shall not use this explicit infinite product solution anywhere in our further analysis.

Thus far we have obtained the Mellin transforms of the conditional generating functions $\widetilde{B}_{i}(z, 1)$. In order to obtain the composite Mellin transform $B^{*}(s)$ and $C^{*}(s)$ of $\widetilde{B}_{u}(z, 1)$ and $B_{u u}(z, 1)$, respectively, we refer to (58) and (59) and, after some algebra, we finally obtain

$$
\begin{align*}
& B^{*}(s)=\mathbf{p}(s) \mathbf{b}(s)+\Gamma(s) x(s-1)  \tag{64}\\
& C^{*}(s)=2 \mathbf{p}(s) \mathbf{b}(s)+\mathbf{p}(s) \mathbf{c}(s)+\Gamma(s) v(s-1) \tag{65}
\end{align*}
$$

where $\mathbf{p}(s)=\left(\pi_{1}^{-s}, \ldots, \pi_{V}^{-s}\right)$ in the stationary case and $\mathbf{p}(s)=\left(p_{1}^{-s}, \ldots, p_{V}^{-s}\right)$ in the nonstationary case. We shall see that the dominant asymptotics of $B^{*}(s)$ and $C^{*}(s)$ are determined by asymptotics of $\mathrm{b}(s)$ and $\mathrm{c}(s)$ which depend on singularities of $\mathrm{Q}(s)$ that we discuss next.

### 3.2 Singularities of the Matrix $Q(s)$

We study here singularities of the matrix $\mathrm{Q}(s)$, which play central role in the asymptotic analysis of the depth. We prove the following lemma that characterizes the location of singularities of $\mathrm{Q}(s)$.

Lemma 3 Let $\mathrm{Q}(s)=\mathrm{I}-\mathrm{P}(s)$ and $\mathrm{P}(s)=\left\{p_{i j}^{-s}\right\}_{i, j \in \mathcal{A}}$. Let $s_{l}$ denote singularites of $\mathrm{Q}(s)$, where $l \in \mathbb{Z}$ is an integer. Then:
(i) Matrix $Q(s)$ is nonsingular for $\Re(s)<-1$, and $s_{0}=-1$ is a simple pole.
(ii) If and only if

$$
\begin{equation*}
\frac{\ln p_{i j}+\ln p_{1 i}-\ln p_{1 j}}{\ln p_{11}} \in \mathbb{Q}, \quad i, j \in \mathcal{A} \tag{66}
\end{equation*}
$$

where $\mathbb{Q}$ is the set of rational numbers, matrix $Q(s)$ has simple poles on the line $\Re(s)=-1$ which can be written as

$$
s_{l}=-1+l \theta \mathbf{i}
$$

where $\mathbf{i}=\sqrt{-1}$ and

$$
\theta=\frac{n_{1}}{n_{2}}\left|\frac{2 \pi}{\ln p_{11}}\right|,
$$

where integers $n_{1}, n_{2}$ are such that $\left\{\left|\frac{n_{1}}{n_{2} \ln p_{11}}\left(\ln p_{i j}-\ln p_{1 i}+\ln p_{1 j}\right)\right|\right\}_{i j=1}^{V}$ is a set of relative primes.
(iii) Finally,

$$
\mathrm{Q}(-1+l \theta \mathrm{i})=\mathrm{E}^{-l} \mathrm{Q}(-1) \mathrm{E}^{l}
$$

where $\mathrm{E}=\operatorname{diag}\left(1, e^{\theta_{12} \mathrm{i}}, \ldots, e^{\theta_{1 V \mathrm{i}}}\right)$ is the diagonal matrix with $\theta_{i k}=-\theta \ln p_{i \mathrm{k}}$.
Proof. Observe that for $\Re(s)<-1$,

$$
\begin{equation*}
\left|1-{p_{i i}^{-s}}^{-s}\right| \geq 1-\left|p_{i i}^{-s}\right|>1-p_{i i}=\sum_{j \neq i} p_{i j} \geq \sum_{j \neq i}\left|p_{i j}^{-s}\right| \tag{67}
\end{equation*}
$$

hence $Q(s)$ is a strictly diagonal dominant matrix, and therefore nonsingular.
Now, we proceed with the proof of part (ii) of the lemma. For $b \neq 0$ such that $Q(-1+b i)$ is singular, let $\mathbf{x}=\left[x_{1}, x_{2}, \ldots, x_{V}\right]^{T} \neq 0$ be a solution of $\mathrm{Q}(-1+b \mathbf{i}) \mathbf{x}=0$, where

$$
\mathrm{Q}(-1+b \mathrm{i})=\left[\begin{array}{cccc}
1-p_{11} e^{\xi_{11} \mathrm{i}} & -p_{12} e^{\xi_{12} \mathrm{i}} & \ldots & -p_{1 V} e^{\xi_{1 V} \mathbf{i}} \\
-p_{21} e^{\xi_{21} \mathbf{i}} & 1-p_{22} e^{\xi_{22} \mathrm{i}} & \ldots & -p_{2 V} e^{\xi_{2 V} \mathbf{i}} \\
\vdots & \vdots & \ddots & \vdots \\
-p_{i 1} e^{\xi_{11} \mathrm{i}} & -p_{i 2} e^{\xi_{i 2} \mathrm{i}} & \ldots & -p_{i V} e^{\xi_{i V} \mathbf{i}} \\
\vdots & \vdots & \ddots & \vdots \\
-p_{V 1} e^{\xi_{V 1} \mathrm{i}} & -p_{V 2} e^{\xi_{V 2} \mathrm{i}} & \ldots & 1-p_{V V} e^{\xi_{V V} \mathrm{i}}
\end{array}\right]
$$

with $\xi_{i k}=-b \ln p_{i k}$. Without loss of generality, suppose $\left|x_{1}\right|=\max \left\{\left|x_{1}\right|,\left|x_{2}\right|, \ldots,\left|x_{V}\right|\right\} \neq 0$ (since $Q(-1+b \mathbf{i})$ is singular). Then

$$
\left(1-p_{11} e^{\epsilon_{11} \mathrm{i}}\right) x_{1}-p_{12} e^{\xi_{12} \mathbf{i}} x_{2}-\ldots-p_{1 V} e^{\xi_{1 V} \mathbf{i}} x_{V}=0
$$

implies

$$
1-p_{11} e^{\xi_{11} \mathbf{i}}=p_{12} e^{\xi_{12} \mathbf{i}} x_{2} / x_{1}+\ldots+p_{1 V} e^{\xi_{1 V} \mathbf{i}} x_{V} / x_{1}
$$

But as in (67)

$$
\left|1-p_{11} e^{\xi_{11} i}\right| \geq 1-p_{11}
$$

and

$$
\left|p_{12} e^{\xi_{12} \mathbf{i}} x_{2} / x_{1}+\ldots+p_{1 V} e^{\xi_{1 V} \mathrm{i}} x_{V} / x_{1}\right| \leq p_{12}+\ldots+p_{1 V}=1-p_{11} .
$$

Thus

$$
\begin{gathered}
1-p_{11} e^{\xi_{11} \mathbf{i}}=1-p_{11} \\
p_{12} e^{\xi_{12} \mathbf{i}} x_{2} / x_{1}+\ldots+p_{1 V} e^{\xi_{1 V} \mathbf{i}} x_{V} / x_{1}=p_{12}+\ldots+p_{1 V} .
\end{gathered}
$$

This implies

$$
e^{\xi_{11} \mathrm{i}}=e^{\xi_{1} \mathrm{i}} x_{i} / x_{1}=1
$$

and $\left|x_{i}\right|=\left|x_{j}\right|$ for any $i, j=1,2, \ldots, V$, so that $e^{\xi_{i i i}}=1$ for all $i$. Define now $\xi_{i}$ such that $x_{i} / x_{1}=e^{-\xi_{1 i} \mathrm{i}}=e^{\xi_{i} \mathrm{i}}$. Then,
for any $1 \leq j \leq V$. Note that since

$$
-p_{j 1}-p_{j 2}-\ldots-p_{j(j-1)}+1-p_{j j}-\ldots,-p_{j v}=0
$$

we must have $e^{\xi_{j i}{ }^{\mathrm{i}} e^{\xi_{i}} e^{-\xi_{j}}=1, \quad \text { and thus }}$

$$
e^{\xi_{j i}}=e^{\left(\xi_{j}-\xi_{i}\right) i}
$$

Hence $-b\left(\ln p_{j i}+\ln p_{1 j}-\ln p_{1 i}\right)=2 \pi n_{j i}$ for some integer $n_{j i}$, and as a conscquence $\left(\ln p_{i j}+\right.$ $\left.\ln p_{1 i}-\ln p_{1 j}\right) / \ln p_{11}$ is rational for any $i, j=1,2, \ldots, V$.

To prove the inverse part of (ii), suppose $b$ is such that $\left[\left.\frac{b}{2 \pi}\left(\ln p_{j i}+\ln p_{1 j}-\ln p_{1 i}\right) \right\rvert\,\right.$ are integers for any $i, j=1,2, \ldots, V$. Then

$$
\begin{aligned}
Q(-1+b i) & =\left[\begin{array}{cccc}
1-p_{11} e^{\xi_{11} \mathrm{i}} & -p_{12} e^{\xi_{12 \mathrm{i}}} & \ldots & -p_{1 V} e^{\xi_{1 V \mathrm{i}}} \\
-p_{21} e^{\xi_{21} \mathrm{i}} & 1-p_{22} e^{\xi_{22} \mathrm{i}} & \ldots & -p_{2 V} e^{\xi_{2 V} \mathrm{i}} \\
\vdots & \vdots & \ddots & \vdots \\
-p_{i 1} e^{\xi_{i 1} \mathrm{i}} & -p_{i 2} e^{\xi_{i 2} \mathrm{i}} & \ldots & -p_{i V} e^{\xi_{i V} \mathrm{i}} \\
\vdots & \vdots & \ddots & \vdots \\
-p_{V 1} e^{\xi_{V 1} \mathrm{i}} & -p_{V 2} e^{\xi_{V 2} \mathrm{i}} & \ldots & 1-p_{V V} e^{\xi_{V V} \mathrm{i}}
\end{array}\right] \\
& =\left[\begin{array}{ccccc}
1-p_{11} & -p_{12} e^{\left(\xi_{1}-\xi_{2}\right) \mathrm{i}} & \ldots & -p_{1 V} e^{\left(\xi_{1}-\xi_{V}\right) \mathrm{i}} \\
-p_{21} e^{\left(\xi_{2}-\xi_{1}\right) \mathrm{i}} & 1-p_{22} & \ldots & -p_{2 V} e^{\left(\xi_{2}-\xi_{V}\right) \mathrm{i}} \\
\vdots & \vdots & \ddots & \vdots \\
-p_{i 1} e^{\left(\xi_{i}-\xi_{1}\right) \mathrm{i}} & -p_{i 2} e^{\left(\xi_{i}-\xi_{2}\right) \mathrm{i}} & \ldots & -p_{i V} e^{\left(\xi_{i}-\xi_{V}\right) \mathrm{i}} \\
\vdots & \vdots & \ddots & \vdots \\
-p_{V 1} e^{\left(\xi_{V}-\xi_{1}\right) \mathrm{i}} & -p_{V 2} e^{\left(\xi_{V}-\xi_{2}\right) \mathrm{i}} & \ldots & 1-p_{V V}
\end{array}\right]
\end{aligned}
$$

$$
=\left[\operatorname{diag}\left(1, e^{-\xi_{2}}, e^{-\xi_{3}}, \ldots, e^{-\xi_{V}}\right)\right]^{-1} \mathrm{Q}(-1) \operatorname{diag}\left(1, e^{-\xi_{2}}, e^{-\xi_{3}}, \ldots, e^{-\xi_{V}}\right)
$$

Since $Q(-1)$ is singular, so $Q(-1+b \mathrm{i})$ is. Hence $s=-1+b \mathbf{i}$ is a pole of $\mathrm{Q}(s)$ if and only if $\left|\frac{b}{2 \pi}\left(\ln p_{j i}+\ln p_{1 j}-\ln p_{1 i}\right)\right|$ are integers for any $i, j=1,2, \ldots, V$. Since $\left\{\left\lvert\, \frac{\theta}{2 \pi}\left(\ln p_{i j}+\ln p_{1 i}-\right.\right.\right.$ $\left.\ln p_{1 j}\right) \|_{i j=1}^{V}$ is a set of relative primes, hence $b=l \theta$ for some integer $l$. Part (ii) is proved. Part (iii) can be inferred from the above proof.

Observe that for the memorylcss case, that is, when $p_{j i}=\pi_{i}$, condition (66) becomes $\frac{\mathrm{ln} \pi_{i}}{\ln \pi_{j}} \in \mathbb{Q}$ for all $i, j$. This agrees with previous known results (cf. [7]).

Finally, as a simple consequence of the above, we prove the convergence of the infinite product that appears (62).

Lemma 4 The product

$$
\prod_{i=0}^{\infty} Q^{-1}(s-i)
$$

converges for $\Re(s)<-1$, and it can be differentiated with respect to $s$ term by term.
Proof. For $\Re(s)<-1$, every factor of the above infinite product is non-singular, and $\|\mathrm{P}(s)\| \leq V p^{-s}$, where $p=\max _{i, j}\left\{p_{i j}\right\}<1$. For $k$ large enough such that $V p^{k}<\frac{1}{2}$, we have $\|\mathrm{Q}(s-k)\| \leq 1+2 V p^{-s+k}$. Since $\sum_{i=k}^{\infty} p^{-s+i}<\infty$, hence $\left|\prod_{i=0}^{\infty} \mathrm{Q}^{-1}(s-i)\right| \leq$ $\prod_{i=0}^{\infty}\left\|Q^{-1}(s-i)\right\|<\infty$.

### 3.3 Asymptotic Expansions for the Moments in the Poisson Model

As outlined above, we seek the asymptotics of $B_{u}(z, 1)$ and $B_{u u}(z, 1)$ for large $z$, which further will lead through depoissonization to asymptotics of the first two moments of the depth. We derive asymptotic expansions of the moments in the Poisson model by applying the inverse Mellin transform. In particular,

$$
\begin{aligned}
\tilde{B}_{u}(z, 1) & =\frac{1}{2 \pi \mathrm{i}} \int_{-\frac{3}{2}-\mathrm{i} \infty}^{-\frac{3}{2}+\mathrm{i} \infty} B^{*}(s) z^{-s} d s \\
\tilde{B}_{u u}(z, 1) & =\frac{1}{2 \pi \mathrm{i}} \int_{-\frac{3}{2}-\mathrm{i} \infty}^{-\frac{3}{2}+\mathrm{i} \infty} C^{*}(s) z^{-s} d s
\end{aligned}
$$

The evaluation of the above integrals is quite standard (e.g., see $[10,14]$ ): We extend the line of integration to a big rectangle right to the integration line, and observe that bottom and top lines contribute negligibly because the gamma function decreases exponentially with the increase in the magnitude of the imaginary part. The right side positioned at, say $d$,
contributes $|z|^{-d}$ for $d \rightarrow \infty$. Thus, the integral is asymptotically equal to minus the sum of residues positioned right to the line of the integration, that is, $\left(-\frac{3}{2}-i \infty,-\frac{3}{2}+i \infty\right)$. But, the residues of the above depend on the singularities of just studied $Q(s)$ and gamma function. To estimate them, we expand the function under the integral around these singularities.

Let us start with the dominant singularity at $s_{0}=-1$, and derive the Laurent expansion of $\mathbf{x}(s)$ and $\mathbf{v}(s)$. By Lemma 3, we can write

$$
\mathrm{Q}^{-1}(s)=\frac{1}{s+1} \mathrm{Q}_{1}+\mathrm{Q}_{2}+O(s+1)
$$

where $\mathrm{Q}_{1}, \mathrm{Q}_{2}$ are $V \times V$ matrices. Since

$$
\begin{aligned}
\mathbf{x}(s-1) & \left.=\psi+\dot{\mathbf{x}}(-2)(s+1)+O(s+1)^{2}\right) \\
\Gamma(s) & =\frac{-1}{s+1}+\gamma-1+O(s+1)
\end{aligned}
$$

we obtain from (54), (62) and (56), (63)

$$
\begin{aligned}
\mathbf{b}(s) & =\Gamma(s) \mathrm{Q}^{-1}(s) \mathbf{x}(s-1)=\frac{1}{(s+1)^{2}} \mathbf{a}_{1}+\frac{1}{s+1} \mathbf{a}_{2}+O(1) \\
\mathbf{c}(s) & =2 \Gamma(s) \mathrm{Q}^{-2}(s) \mathrm{P}(s) \mathbf{x}(s-1)+\Gamma(s) \mathrm{Q}^{-1}(s) \mathbf{v}(s-1) \\
& =\frac{1}{(s+1)^{3}} \mathbf{f}_{\mathrm{I}}+\frac{1}{(s+1)^{2}} \mathrm{f}_{2}+O\left(\frac{1}{s+1}\right)
\end{aligned}
$$

where $a_{1}, a_{2}, f_{1}$ and $f_{2}$ are vectors of constants for which explicit formulæ are presented below the next lemma. In addition, by (64), (65) and the fact that $x(s-1)=1+O(s+1)$, $v(s-1)=O(s+1)$, we have

$$
\begin{align*}
B^{*}(s) & =\frac{1}{(s+1)^{2}} \pi \mathbf{a}_{1}+\frac{1}{s+1}\left(\pi \mathbf{a}_{2}+\dot{\mathbf{p}}(-1) \mathbf{a}_{1}-1\right)+O(1)  \tag{68}\\
C^{*}(s) & =\frac{1}{(s+1)^{3}} \pi \mathbf{f}_{1}+\frac{1}{(s+1)^{2}}\left(\dot{\mathbf{p}}(-1) \mathbf{f}_{1}+\pi \mathrm{f}_{2}+2 \pi \mathrm{a}_{1}\right)+O\left(\frac{1}{s+1}\right) \tag{69}
\end{align*}
$$

where $\dot{\mathbf{p}}(-1)=\left.\frac{d}{d s}\right|_{s=-1} \mathbf{p}(s)=\left(-\pi_{1} \ln \pi_{1}, \ldots,-\pi_{V} \ln \pi_{V}\right)$.
To derive explicit expressions for the vectors $a_{1}, a_{2}, f_{1}$ and $f_{2}$ we need the following lemma which proof is standard and is omitted (detailed proof can be found in [22]).

Lemma 5 Let us define

$$
\Pi=\left[\begin{array}{cccc}
\pi_{1} & \pi_{2} & \ldots & \pi_{V} \\
\pi_{1} & \pi_{2} & \ldots & \pi_{V} \\
\vdots & \vdots & \ddots & \vdots \\
\pi_{1} & \pi_{2} & \ldots & \pi_{V}
\end{array}\right]=\psi \otimes \pi
$$

and let $\mathrm{Q}^{*}=\left\{q_{j i}^{*}\right\}_{i, j=1}^{V}$ be the adjoint matrix of $\left.\mathrm{Q}(s)\right|_{s=-1}$. Then

$$
\begin{align*}
& \pi \Pi=\pi, \quad \Pi^{2}=\Pi, \quad \Pi \psi=\psi  \tag{70}\\
& \left.\frac{d}{d s} \operatorname{det} \mathrm{Q}(s)\right|_{s=-1}=\left.\frac{d}{d s} \operatorname{det} \mathrm{Q}(s)\right|_{s=-1+j b}=-\omega h, \quad q_{j i}^{*}=\omega \pi_{i}, \quad \mathrm{Q}^{*}=\omega \Pi,  \tag{71}\\
& \mathrm{Q}_{1}=-\frac{1}{h} \Pi, \quad \mathrm{Q}_{2}=-\frac{\dot{\mathrm{Q}}^{*}}{\omega h}-\frac{\beta}{2 \omega h^{2}} \Pi,\left.\quad \mathrm{Q}_{1}\right|_{s=-1+j b}=-\frac{1}{h} \mathrm{E}^{-1} \Pi \mathrm{E} .  \tag{72}\\
& \mathrm{Q}_{1}^{2} \dot{\mathrm{P}}(-1) \psi=\frac{1}{h} \psi, \quad \mathrm{Q}_{1}^{2} \dot{\mathrm{P}}(-1)=\frac{1}{h} \Pi, \tag{73}
\end{align*}
$$

where $s=-1+b \mathrm{i}$ is a pole of $\mathrm{Q}^{-1}(s)$.
Using the above, we finally obtain after some tedious algebra

$$
\begin{aligned}
\mathrm{a}_{1} & =-\mathrm{Q}_{1} \psi=\frac{1}{h} \psi, \\
\mathrm{a}_{2} & =(\gamma-1) \mathrm{Q}_{1} \psi-\mathrm{Q}_{2} \psi-\mathrm{Q}_{1} \dot{\mathbf{x}}(-2), \\
& =-\frac{1}{h}(\gamma-1) \psi+\frac{1}{\omega h} \dot{Q}^{*} \psi+\frac{\beta}{2 \omega h^{2}} \psi+\frac{1}{h} \pi \dot{\mathrm{x}}(-2) \psi, \\
\mathrm{f}_{1} & =-2 \mathrm{Q}_{1}^{2} \psi=\frac{-2}{h^{2}} \psi, \\
\mathbf{f}_{2} & =2\left((\gamma-1) \mathrm{Q}_{1}^{2} \psi-\left(\mathrm{Q}_{1} \mathrm{Q}_{2}+\mathrm{Q}_{2} \mathrm{Q}_{1}\right) \psi-\mathrm{Q}_{1}^{2}(\dot{\mathrm{P}} \psi+\mathrm{P} \dot{\mathbf{x}}(-2))\right) \\
& =2\left(\frac{\gamma-1}{h^{2}} \psi-\frac{1}{\omega h^{2}}\left(\Pi \dot{Q}^{*} \psi+\dot{\mathrm{Q}}^{*} \psi\right)-\frac{\beta}{\omega h^{3}} \psi-\mathrm{Q}_{1}^{2}(\dot{\mathrm{P}} \psi+\mathrm{P} \dot{\mathbf{x}}(-2))\right) \\
& =2\left(\frac{\gamma-1}{h^{2}} \psi-\frac{1}{\omega h^{2}}\left(\Pi \dot{Q}^{*} \psi+\dot{\mathrm{Q}}^{*} \psi\right)-\frac{\beta}{\omega h^{3}} \psi-\frac{1}{h} \psi-\frac{1}{h^{2}} \Pi \dot{\mathbf{x}}(-2)\right) \\
& =2\left(\frac{\gamma-1}{h^{2}}-\frac{\beta}{\omega h^{3}}-\frac{1}{h}-\frac{1}{h^{2}} \pi \dot{\mathbf{x}}(-2)\right) \psi-\frac{2}{\omega h^{2}}\left(\Pi \dot{Q}^{*} \psi+\dot{\mathrm{Q}}^{*} \psi\right) .
\end{aligned}
$$

In summary, using (68) we obtain the following expansion of $B^{*}(s)$ around the dominant pole at $s_{0}=-1$

$$
B^{*}(s)=\frac{1}{(s+1)^{2}} \frac{1}{h}+\frac{1}{s+1}\left(-\frac{1}{h}(\gamma-1)+\frac{1}{\omega h} \pi \dot{Q}^{*} \psi+\frac{\beta}{2 \omega h^{2}}+\frac{1}{h} \pi \dot{\mathbf{x}}(-2)+\frac{h_{\pi}}{h}-1\right)+O(1)
$$

while (69) becomes

$$
\begin{aligned}
C^{*}(s) & =\frac{-2}{h^{2}(s+1)^{3}}+\frac{2}{(s+1)^{2}}\left(\frac{-h_{\pi}}{h^{2}}+\frac{\gamma-1}{h^{2}}-\frac{\beta}{\omega h^{3}}\right. \\
& \left.-\frac{1}{h^{2}} \pi \dot{\mathbf{x}}(-2)-\frac{2}{\omega h^{2}} \pi \dot{\mathrm{Q}}^{*} \psi\right)+O\left(\frac{1}{s+1}\right)
\end{aligned}
$$

Now, we deal with the asymptotics related to the nondominant poles $s_{l}=-1+l \theta \mathrm{i}$ for $l \neq 0$. By Lemma 3, we have

$$
\mathrm{Q}(s)=\frac{-1}{h} \frac{1}{s+1-l \overline{\mathrm{i}}} \times \mathrm{E}^{-1} \Pi \mathrm{E}+O(1)
$$

Therefore,

$$
\begin{aligned}
\mathrm{b}(s) & =\Gamma(s) \mathrm{Q}^{-1}(s) \mathbf{x}(s-1) \\
& =\Gamma(-1+l \theta \mathrm{i})\left(-\frac{1}{h} \mathrm{E}^{-l} \Pi \mathrm{E}^{l}\right) \times(-2+l \theta \mathrm{i}) \frac{1}{s+1-l \theta \mathrm{i}}+O(1), \\
& =-\frac{1}{h} \Gamma(-1+l \theta \mathrm{i})\left(\pi \mathrm{E}^{l} \mathbf{x}(-2+l \theta \mathrm{i})\right) \mathrm{E}^{-l} \psi \frac{1}{s+1-l \theta \mathbf{i}}+O(1), \\
& =-\frac{1}{h} \rho_{l} \psi(l) \frac{1}{s+1-l \theta \mathbf{i}}+O(1), \\
\mathrm{c}(s) & =2 \Gamma(s) \mathrm{Q}^{-2}(s) \mathrm{P}(s) \mathbf{x}(s-1)+\Gamma(s) \mathrm{Q}^{-1}(s) \mathbf{v}(s-1), \\
& =2 \frac{1}{h^{2}} \Gamma(-1+l \theta \mathrm{i})\left(\pi \mathrm{E}^{l} \times(-2+l \theta \mathbf{i})\right) \mathrm{E}^{-l} \psi \frac{1}{(s+1-l \theta \mathbf{i})^{2}}+O\left(\frac{1}{s+1-l \theta \mathbf{i}}\right), \\
& =\frac{2}{h^{2}} \rho_{l} \psi(l) \frac{1}{(s+1-l \theta \mathbf{i})^{2}}+O\left(\frac{1}{s+1-l \theta \mathbf{i}}\right) .
\end{aligned}
$$

where $\rho_{l}=\Gamma(-1+l \theta \mathbf{i})\left(\pi \mathrm{E}^{l} \times(-2+l \theta \mathrm{i})\right)$ and $\psi(l)=\mathrm{E}^{-l} \psi$. In summary, by (64) and (65) at $s=-1+l \theta \mathrm{i}$ we obtain

$$
\begin{aligned}
\mathcal{B}^{*}(s) & =-\frac{1}{h} \rho_{l} \mathrm{p}(-1+l \theta \mathrm{i}) \psi(l) \frac{1}{s+1-l \theta \mathbf{i}}+O(1), \\
C^{*}(s) & =\frac{2}{h^{2}} \rho_{l} \mathrm{p}(-1+l \theta \mathbf{i}) \psi(l) \frac{1}{(s+1-l \theta \mathbf{i})^{2}}+O\left(\frac{1}{s+1-l \theta \mathrm{i}}\right)
\end{aligned}
$$

Finally, we handle singularities in the half plane $\Re(s)>-1$. We consider two cases: $-1<\Re(s) \leq 0$ and $\Re(s)>0$. Let $\mathcal{Z}_{*}$ be the set of singularities $s^{*}$ of $Q(s)$ lying in the strip $-1<\Re\left(s^{*}\right) \leq 0$, while $\mathcal{Z}_{+}$be the set of singularities in $\Re(s)>0$. At the pole $s^{*} \in \mathcal{Z}_{*}$

$$
B^{*}(s)=\frac{1}{s-s^{*}} \pi\left(s^{*}\right) \Gamma\left(s^{*}\right) R\left(s^{*}\right) \mathbf{x}\left(s^{*}-1\right)=\frac{1}{s-s^{*}} r\left(s^{*}\right)
$$

where $R\left(s^{*}\right)$ is the residue matrix of $Q^{-1}(s)$ at $s^{*}$. Note that $s=0$ is the double pole. An application of the inverse Mellin transform gives for $z \rightarrow \infty$,

$$
\begin{align*}
\widetilde{B}_{u}(z, 1) & =\frac{1}{h} z \ln z+\frac{1}{h}\left(\gamma-1-\frac{\beta}{2 \omega h}-\frac{1}{\omega} \pi \dot{Q}^{*} \psi-\pi \dot{\mathrm{x}}(-2)+h-h_{\pi}\right) z  \tag{74}\\
& -\frac{1}{h} \sum_{l=0} \rho_{l} \pi \psi(l) z^{1-l \theta \mathbf{i}}+\sum_{s^{*} \in \mathcal{Z}_{-}} r\left(s^{*}\right) z^{-s^{*}}+r(0)+\sum_{s^{*} \in \mathcal{Z}_{+}} r\left(s^{*}\right) z^{-s^{*}} \\
& =\frac{1}{h} z \ln z+\frac{1}{h}\left(\gamma-1-\frac{\beta}{2 \omega h}-\frac{1}{\omega} \pi \dot{Q}^{*} \psi-\pi \dot{\mathbf{x}}(-2)+h-h \pi\right) z+\delta_{1}(z)+O(\ln z)
\end{align*}
$$

Observe also that $r(0)+\sum_{s^{*} \in \mathcal{Z}_{+}} \tau\left(s^{-}\right) z^{-s^{*}}=O(\ln z)$. In a similar manner, we obtain

$$
\begin{align*}
\widetilde{B}_{u u}(z, 1) & =\frac{1}{h^{2}} z \ln ^{2} z+\frac{2}{h^{2}}\left(\gamma-1-\frac{\beta}{\omega h}-\frac{2}{\omega} \pi \dot{Q}^{*} \psi-h \pi-\pi \dot{\mathrm{x}}(-2)\right) z \ln z \\
& +\frac{2}{h^{2}} \sum_{l=0} \rho_{l} \pi(1-l \theta \mathrm{i}) \psi(l) z^{1-l \theta \mathrm{i}} \ln z+O(z) \tag{76}
\end{align*}
$$

as $z \rightarrow \infty$ in a cone around the real axis.

### 3.4 Analytic Depoissonization

The above asymptotic formulæ concern the behavior of the Poisson mean and the second factorial moment as $z \rightarrow \infty$. More precisely, we had to restrict the growth of $z$ to a linear cone $S_{\theta}=\{z:|\arg (z)| \leq \theta\}$ for some $|\theta|<\pi / 2$. But, our original goal was to derive asymptotics of the mean $\mathrm{E}\left[D_{m}\right]$ and the variance $\operatorname{Var}\left[D_{m}\right]$ in the MI model. To infer such a behavior from its Poisson model asymptotics, we must apply the so called depoissonization lemma. This lemma basically says that $m \mathbf{E}\left[D_{m}\right] \sim \widetilde{B}\left(u(m, 1)\right.$ and $m \mathbf{E}\left[D_{m}\left(D_{m}-1\right)\right] \sim \widetilde{B}_{u u}(m, 1)$ under some weak conditions that will be easy to verify in our case. The reader is referred to $[7,8,9]$ for more details about depoissonization lemma. For completeness, however, we review some depoissonization results that are uscful for our problem.

Let us consider a general problem: For a random variable $X_{n}$ define $g_{n}$ as a functional of the distribution of $X_{n}$ (e.g., $g_{n}=\mathrm{E}\left[X_{n}\right]$ or $\left.g_{n}=\mathrm{E}\left[X_{n}^{2}\right]\right)$, or, in general, assume $g_{n}$ is a sequence of $n$. In some situations (e.g., for limiting distributions we need to consider the generating function $G_{n}(u)=\mathbf{E}\left[u^{X_{n}}\right]$ (for a random variable $X_{n}$ ) for a complex $u$ which can be viewed as such a $g_{n}$ (with a parameter $u$ belonging to a compact set). Define the Poisson transform of $g_{n}$ as $\tilde{G}(z)=\sum_{n=0}^{\infty} g_{n} \frac{z^{n}}{n!} e^{-z}$ (or more generally: $\tilde{G}(z, u)=\sum_{n=0}^{\infty} G_{n}(u) \frac{z^{n}}{n!} e^{-z}$ for $u$ in a compact set). Assume that we know the asymptotics of $\tilde{G}(z)$ for $z$ large and belonging to a cone $S_{\theta}=\{z:|\arg (z)| \leq \theta\}$ for some $|\theta|<\pi / 2$. How can we infer asymptotics of $g_{\pi}$ from $\tilde{G}(z)$ ? An answer is given in the depoissonization lemma below (cf. $[7,8,9]$ ):

## Lemma 6 (Depoissonization Lemma)

(i) Let $\tilde{G}(z)$ be the Poisson transform of a sequence $g_{n}$ that is assumed to be an entire function of $z$. We postulate that for $0<|\theta|<\pi / 2$ the following two conditions simultaneously hold for some numbers $A, B, \xi>0, \beta$, and $\alpha<1$ :
(I) For $z \in S_{0}$

$$
\begin{equation*}
|z|>\xi \Rightarrow|\tilde{G}(z)| \leq B|z|^{\beta} \phi(|z|) \tag{77}
\end{equation*}
$$

where $\phi(z)$ is a slowly varying function (e.g., $\phi(z)=\log ^{d} z$ for some $d>0$ ),
(O) For $z \notin S_{0}$

$$
\begin{equation*}
|z|>\xi \quad \Rightarrow \quad\left|\tilde{G}(z) e^{z}\right| \leq A \exp (\alpha|z|) \tag{78}
\end{equation*}
$$

Then, for large $n$,

$$
\begin{equation*}
g_{n}=\tilde{G}(n)+O\left(n^{\beta-1} \phi(n)\right) \tag{79}
\end{equation*}
$$

or more precisely:

$$
g_{n}=\tilde{G}(n)-\frac{1}{2} \tilde{G}^{\prime \prime}(n)+O\left(n^{\beta-2} \phi(n)\right)
$$

(ii) If the above two conditions, namely (I) and (O), hold for $\tilde{G}(z, u)$ for $u$ belonging to a compact set $\mathcal{U}$, then

$$
\begin{equation*}
G_{n}(u)=\widetilde{G}(n, u)+O\left(n^{\beta-1} \phi(n)\right) \tag{80}
\end{equation*}
$$

for latge $n$ and uniformly in $u \in \mathcal{U}$.
(iii) Let $g(z)$ be an analytic continuation of a sequence $g_{n}$ whose Poisson transform is $\tilde{G}(z)$, and such that $g(z)=O\left(z^{\beta}\right)$ in a linear cone. Then, for some $\theta_{0}$ and for all linear cones $S_{\theta}$ $\left(\theta<\theta_{0}\right)$, there exists $\alpha<1$ and $A>0$ such that

$$
z \notin S_{0} \quad \Rightarrow \quad\left|\widetilde{G}(z) e^{z}\right| \leq A e^{\alpha|z|} .
$$

In summary, when $g(z)$ has a polynomial growth, then conditions (I) and (O) above are automatically satisfied and (79) holds.

Now, we are equipped with the tool to depoissonize $\widetilde{B}_{u}(z, 1)$ and $\widetilde{B}_{u u}(z, 1)$, and obtain asymptotics for the mean $\mathbf{E}\left[D_{m}\right]$ and the variance $\operatorname{Var}\left[D_{m}\right]$. Observe that $\mathbf{E}\left[D_{m}\right]=$ $O(m \ln m)$ and $\operatorname{Var}\left[D_{m}\right]=O\left(m \log ^{2} m\right)$, hence by Lemma 6 we can depoissonize the Poisson estimates. We obtain

$$
\begin{align*}
\mathrm{E}\left[D_{m}\right] & \left.=\frac{1}{h} \ln m+\frac{1}{h}\left(\gamma-1+h-h_{\pi}-\frac{\beta}{2 \omega h}-\frac{1}{\omega} \pi \dot{\mathrm{Q}}^{*} \psi-\pi \dot{\mathrm{x}}(-2)\right)\right)  \tag{81}\\
& +\delta_{1}(m)+O\left(\frac{\ln m}{m}\right) .
\end{align*}
$$

To derive the variance, we observe that $\sum_{s^{*} \in \mathcal{Z} .} r\left(s^{*}\right) m^{-s^{*}}=O\left(m^{-\delta}\right)$ for some $\delta>0$, thus such terms will not appear explicitly in the following formula where only $\Omega(\ln m)$ terms are considered. Again, by Lemma 6 we arrive at

$$
\begin{aligned}
\operatorname{Var}\left[D_{m}\right] & =\frac{1}{h^{2}} \ln ^{2} m+\frac{2}{h^{2}}\left(\gamma-1-\frac{\beta}{\omega h}-\frac{2}{\omega} \pi \dot{\mathrm{Q}}^{*} \psi-h_{\pi}-\pi \dot{\mathbf{x}}(-2)\right) \ln m \\
& +\mathrm{E}\left[D_{m}\right]-\mathbf{E}\left[D_{m}^{2}\right]+O(1) \\
& =\frac{1}{h^{2}} \ln ^{2} m+\frac{2}{h^{2}}\left(\gamma-1-\frac{\beta}{\omega h}-\frac{2}{\omega} \pi \dot{\mathrm{Q}}^{*} \psi-h_{\pi}-\pi \dot{\mathrm{x}}(-2)\right) \ln m+\frac{1}{h} \ln m
\end{aligned}
$$

$$
\begin{align*}
& -\frac{1}{h^{2}} \ln ^{2} m-\frac{2}{h^{2}}\left(\gamma-1+h-h_{\pi}-\frac{\beta}{2 \omega h}-\frac{1}{\omega} \pi \dot{Q}^{*} \psi-\pi \dot{\mathrm{x}}(-2)\right) \ln m+O(1) \\
& =\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{Q}^{*} \psi-h^{2}\right) \ln m+O(1) \tag{82}
\end{align*}
$$

Thus, (7) and (8) of Theorem 1 are proved.

### 3.5 Limiting Distribution

Finally, we prove here the limiting distribution of the depth $D_{m}$, just finishing the proof of Theorem 1. We repeat here the system of functional equations ( 51 ), that is,

$$
\begin{array}{ccc}
B_{z}^{1}(z, u)+B^{1}(z, u) & = & u\left[B^{1}\left(p_{11} z, u\right)+\cdots+B^{V}\left(p_{1 V} z, u\right)\right]+1 \\
\cdots & = & \cdots \\
B_{z}^{V}(z, u)+B^{V}(z, u) & = & u\left[B^{1}\left(p_{V 1} z, u\right)+\cdots+B^{V}\left(p_{V V} z, u\right)\right]+1
\end{array}
$$

Obscrve that $B^{i}(z, 1)-z=0, B(z, 1)-z=0, B^{i}(z, u)-z=(u-1) A_{i}(u, z)$, and $B(z, u)-z=$ ( $u-1$ ) $A(u, z)$, wherc $A_{i}(u, z)$ is a power series of $u$ and thus analytic function of $z$. Let

$$
\begin{aligned}
& Z_{i}^{*}(u, s)=\mathcal{M}\left[B^{i}(z, u)-z ; s\right]=\Gamma(s) \xi_{i}(u, s)=(u-1) A_{i}^{*}(u, s), \quad i \in \mathcal{A} \\
& Z^{*}(u, s)=\mathcal{M}[B(z, u)-z ; s]=\Gamma(s) \xi(u, s)=(u-1) A^{*}(u, s)
\end{aligned}
$$

be the Mellin transforms, where $\xi_{i}(u, s)$ and $\xi(u, s)$ are unknown functions.
Lemma 7 The Mellin transforms $Z_{i}^{*}(u, s), Z^{*}(u, s), A_{i}^{*}(u, s)$ and $A^{*}(u, s)$ exist for $\Re(s) \in$ $(-2,-1)$. In addition, $Z_{i}^{*}(u,-2)=u-1, A_{i}^{*}(u,-2)=1, Z^{*}(u,-2)=u-1, A^{*}(u,-2)=1$.

Proof. By the same argument as in Lemma 2 of [13].
We proceed along the same lines as before, thus leaving out detailed explanations. After applying the Mellin transform to the above system of functional equations, we find

$$
\begin{array}{rlc}
Z^{*}(u, s)-(s-1) Z^{*}(u, s-1) & =u\left[Z_{1}^{*}(u, s) \pi_{1}^{-s}+\ldots+Z_{V}^{*}(u, s) \pi_{V}^{-s}\right] \\
Z_{1}^{*}(u, s)-(s-1) Z_{1}^{*}(u, s-1) & = & u\left[Z_{1}^{*}(u, s) p_{11}^{-s}+\cdots+Z_{V}^{*}(u, s) p_{1 V}^{-s}\right] \\
\cdots & = & \cdots \\
Z_{V}^{*}(u, s)-(s-1) Z_{V}^{*}(u, s-1) & = & u\left[Z_{1}^{*}(u, s) p_{V 1}^{-s}+\cdots+Z_{V}^{*}(u, s) p_{V V}^{-s}\right] .
\end{array}
$$

Let

$$
\Gamma(s) \xi(u, s)=\left[\begin{array}{c}
Z_{1}^{*}(u, s) \\
Z_{2}^{*}(u, s) \\
\vdots \\
Z_{V}^{*}(u, s)
\end{array}\right]=(u-1)\left[\begin{array}{c}
A_{1}^{*}(u, s) \\
A_{2}^{*}(u, s) \\
\vdots \\
A_{V}^{*}(u, s)
\end{array}\right]=(u-1) \mathbf{a}(u, s) .
$$

Then, the above becomes

$$
\boldsymbol{\xi}(u, s)-\boldsymbol{\xi}(u, s-1)=u \mathrm{P}(s) \boldsymbol{\xi}(u, s),
$$

that is,

$$
\boldsymbol{\xi}(u, s)=[1-u \mathrm{P}(s)]^{-1} \xi(u, s-1) .
$$

This leads to

$$
Z^{*}(u, s)=u \mathbf{p}(s) \Gamma(s)[l-u \mathrm{P}(s)]^{-1} \xi(u, s-1)+\Gamma(s) \xi(u, s-1) .
$$

Let now set $u=e^{t}$ for complex $t \rightarrow 0$ so that $u$ is in the vicinity of $u=1$. We denote by $s_{k}(t), k=0, \pm 1, \pm 2, \ldots$ singularities of $\mathrm{Q}^{-1}(t, s)=\left(\mathrm{I}-e^{t} \mathrm{P}(s)\right)^{-1}$. Then, at $s=s_{k}(t)$,

$$
\begin{equation*}
Z^{\mathbf{2}}\left(e^{t}, s_{k}(t)\right)=e^{t} \pi\left(s_{k}(t)\right) \Gamma\left(s_{k}(t)\right) \mathrm{R}_{k} \xi\left(e^{t}, s_{k}(t)-1\right) \frac{1}{s-s_{k}(t)}+O(1) \tag{83}
\end{equation*}
$$

where $\mathrm{R}_{k}$ is the residue matrix of $\mathrm{Q}^{-1}(u, s)=[1-u \mathrm{P}(s)]^{-1}$ at $s=s_{k}(t)$. In addition, one must consider two poles of the gamma function $\Gamma(s)$ at $s_{-1}=-1$ and $s_{0}=0$. The latter pole contribute $O(1)$ while the former $-z \xi(u,-1)$. But, by Lemma 7 we know that $\xi(u,-1)=1$, thus the total contribution of these two poles is $-z+O(1)$. Therefore, by the inverse Mellin transform,

$$
\tilde{B}\left(e^{t}, z\right)=e^{t} \sum_{k=-\infty}^{\infty} \pi\left(s_{k}(t)\right) \Gamma\left(s_{k}(t)\right) \mathrm{R}_{k} \xi\left(e^{t}, s_{k}(t)-1\right) z^{-s_{k}(t)}+O(1)
$$

as $z \rightarrow \infty$ in a conc. As before, the leading contribution to the asymptotics comes from the pole $s_{0}(t)$.

To obtain an asymptotic expansion for the original generating function $B_{m}\left(e^{t}\right)$ we apply the depoissonization lemma Lemma 6 (ii). Since $\tilde{B}\left(z, e^{\ell}\right)=O(z \log z)$, we conclude that $B_{m}\left(e^{l}\right)=\tilde{B}\left(m, e^{l}\right)+O(\log m)$ where

$$
\begin{aligned}
B\left(m, e^{t}\right) & =e^{t} p\left(s_{0}(t)\right) \Gamma\left(s_{0}(t)\right) \mathrm{R}_{0} \xi\left(e^{t}, s_{0}(t)-1\right) m^{-s_{0}(t)} \\
& +e^{t} \sum_{k \neq 0} \mathrm{p}\left(s_{k}(t)\right) \Gamma\left(s_{k}(t)\right) \mathrm{R}_{k} \xi\left(u, s_{k}(t)-1\right) m^{-s_{k}(t)}+O(1)
\end{aligned}
$$

Let now as in (8) and (20) (see also Appendix)

$$
v=\frac{\ddot{\lambda}(-1)-\dot{\lambda}^{2}(-1)}{\dot{\lambda}^{3}(-1)}=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{Q}^{*} \psi-h^{2}\right)
$$

which is the coefficient at $\ln n$ of the $\operatorname{Var}\left[D_{m}\right]$. As in [6] we obtain, after some algebra,

$$
\begin{aligned}
s_{0}(t) & =-1-\frac{t}{h}-\frac{v^{2} t^{2}}{2}+O\left(t^{3}\right), \\
\mathrm{R}_{0} & =-\frac{1}{h} \Pi+O(t), \\
\Gamma\left(s_{0}(t)\right) & =-\frac{h}{t}+O(1), \\
\boldsymbol{\xi}\left(s_{0}(t)-1\right) & =t \psi+O\left(t^{2}\right), \\
\mathrm{p}\left(s_{0}(t)\right) & =\pi+O(t) .
\end{aligned}
$$

To derive the above, we just observe that the expansion of $s_{0}(t)$ is obtained via the Lagrange inversion of $1-e^{t} \lambda(s)$, or better, of function $t+\log \lambda(s)$, at $s=-1$ which results in $t+(s+$ 1) $\dot{\lambda}(-1)+(s+1)^{2}\left(\frac{\ddot{\lambda}(-1)-(\dot{\lambda}(-1))^{2}}{2}\right)+O(s+1)^{3}$. We again identify $\dot{\lambda}(-1)=h$. The residue $\mathrm{R}_{0}$ is computed by using the fact that $\mathrm{Q}^{-1}\left(e^{t}, s\right)=\left(1-e^{\ell} \lambda(s)\right)^{-1} \psi(s) \otimes \pi(s)+O(1)$. Observe also that

$$
\lim _{t \rightarrow 0} \mathrm{p}\left(s_{0}(t)\right) \Gamma\left(s_{0}(t)\right) \mathrm{R}_{0} \xi\left(e^{l}, s_{0}(t)-1\right)=\pi \Pi \psi=1
$$

We now set $t=\frac{\tau}{\sigma_{m}}=O(1 / \sqrt{\ln m})$ for some fixed $\tau$ and $\sigma_{m}=\sqrt{\operatorname{Var} D_{m}}$. Then $m^{-1-s_{0}(t)}=e^{\tau \mu_{m} / \sigma_{m}+\frac{r^{2}}{2}}(1+O(t))$ and $D_{m}\left(e^{t}\right)=B\left(e^{t}\right) / m$ leading to

$$
\begin{aligned}
& e^{-\tau \mu_{m} / \sigma_{m}} D_{m}\left(e^{\tau / \sigma_{m}}, m\right)=e^{-\tau \mu_{m} / \sigma_{\mathrm{m}}}\left(e^{\tau \mu_{\mathrm{m}} / \sigma_{m}+\frac{\tau^{2}}{2}}(1+O(t))\right. \\
+ & \left.e^{-t} m^{-1-s_{0}(t)} \sum_{k \neq 0}(u-1)\left(s_{k}(t)-1\right) \mathbf{p}\left(s_{k}(t)\right) \mathrm{R}_{k} \mathrm{a}\left(s_{k}\left(e^{t}\right)-1, u\right) m^{s_{0}(t)-s_{k}(t)}+O\left(\frac{\log m}{m}\right)\right) \\
= & e^{\frac{r^{2}}{2}}\left(1+t O\left(\sum_{k \neq 0}\left(s_{k}(t)-1\right) \mathbf{p}\left(s_{k}(t)\right) \mathrm{R}_{k} \mathrm{a}\left(s_{k}\left(e^{t}\right)-1, u\right)\right)\right) .
\end{aligned}
$$

In the above, we use the fact that $\Re\left(s_{0}(t)\right) \leq \Re\left(s_{k}(t)\right)$ proved in [6] which allows to bound $\left|m^{s_{0}(t)-s_{k}(t)}\right| \leq 1$. To complete the proof, it suffices to show that the sum appearing above is $O(1)$. Let $s_{k}(t)=x_{k}(t)+y_{k}(t) \mathbf{i}$,

$$
\begin{aligned}
\left|\sum_{k \neq 0}\left(s_{k}(t)-1\right) \mathbf{p}\left(s_{k}(t)\right) \mathrm{R}_{k} \mathbf{a}\left(s_{k}\left(e^{t}\right)-1\right)\right| & \leq \sum_{k \neq 0}\left|\left(s_{k}(t)-1\right)\right|\left\|\mathbf{p}\left(s_{k}(t)\right)\right\|\left\|\mathrm{R}_{k}\right\|\left\|\mathrm{a}\left(s_{k}\left(e^{t}\right)-1\right)\right\| \\
& \leq \sum_{k \neq 0} \frac{1}{\left|y_{k}\right|^{M}}=O(1)
\end{aligned}
$$

Here, we use the fact that $A_{i}(u, z)$ is infinitely differentiable, thus its Mellin transform satisfies $\lim _{y \rightarrow \infty}|y|^{M} A_{i}^{*}(u, x+y \mathrm{i})=0$ for any $M>0$.

In summary, we just shown that

$$
e^{-\tau \mu_{m} / \sigma_{m}} D_{m}\left(e^{\tau / \sigma_{m}}\right)=e^{\frac{\tau^{2}}{2}}\left(1+O\left(\frac{1}{\sqrt{\ln m}}\right)\right)
$$

which completes the proof of Theorem 1.

### 3.6 Non-Stationary MI Model

Here, we show how to adopt the above derivations to the non-stationary model in which the initial distribution is $\mathbf{p}$ instead of $\pi$. First of all, observe that $\mathbf{p}$ appears in equation (4) while the conditional generating functions $B_{m}^{i}(u)$ still satisfy (5). Thus, in (50) we must replace $\pi_{i}$ by $p_{i}$, but again (51) stays unchanged. This leads to the following Mellin transforms of $B_{u}(z, 1)$ and $B_{u u}(z, 1)$ in the non-stationary case

$$
\begin{align*}
B^{*}(s) & =\mathbf{p}(s) \mathbf{b}(s)+\Gamma(s) x(s-1)  \tag{84}\\
C^{*}(s) & =2 \mathbf{p}(s) \mathbf{b}(s)+\mathbf{p}(s) \mathbf{c}(s)+\Gamma(s) v(s-1) \tag{85}
\end{align*}
$$

where $\mathbf{p}(s)=\left(p_{1}^{-s}, \ldots, p_{V}^{-s}\right)$. Observe, however, that $\mathbf{b}(s)$ and $\mathbf{c}(s)$ are exactly the same as in the stationary MI model. Since, as we discussed before, the asymptotics of the mean and the variance in the Poisson model depend mostly on the asymptotics of $b(s)$ and $c(s)$, we may expect similar asymptotics results for the non-stationary model. Indeed, we obtain the following expansions of $B^{*}(s)$ and $C^{*}(s)$ around the dominant pole $s_{0}=-1$

$$
\begin{align*}
B^{\prime \prime}(s) & =\frac{1}{(s+1)^{2}} \mathbf{p} \mathbf{a}_{1}+\frac{1}{s+1}\left(\mathbf{p a}_{2}+\mathbf{p}^{\prime}(-1) \mathbf{a}_{1}-1\right)+O(1)  \tag{86}\\
C^{\prime \prime}(s) & =\frac{1}{(s+1)^{3}} \mathbf{p} \mathbf{f}_{1}+\frac{1}{(s+1)^{2}}\left(\mathbf{p}^{\prime}(-1) \mathbf{f}_{1}+\mathbf{p f}_{2}+2 \pi \mathbf{a}_{1}\right)+O\left(\frac{1}{s+1}\right) \tag{87}
\end{align*}
$$

In view of the above, we conclude that the only term effected by the non-stationarity assumption is related to $\mathbf{p}^{\prime}(-1)$ (and also the fluctuating function) which is responsible for replacing $h_{\pi}$ by $h_{\mathbf{p}}$ in the final results. Similar conclusions hold for the limiting distribution. This proves Corollary 1.

## 4 Analysis of Fixed Number of Phrases (GK) Model

In this section we prove Theorem 2 using a combination of probabilistic and analytic techniques. We start our discussion with the introduction of the so called tree-path that plays a crucial role in the analysis. We study its property in Section 4.1, and in Section 4.2 we make a connection between the tree-path and the depth (i.e., phrase). Finally, in Section 4.3 we obtain the limiting distribution for the phrase while in Section 4.4 we establish the existence of the moments, thus proving Theorem 2.

### 4.1 Tree-Path in Digital Search Trees

We consider a DST tree $\mathcal{T}_{m}$ built over $m$ strings regardless of the model of strings generation (e.g., MI, GK, or hybrids). For $k \leq m$ we denote by $I_{k}\left(\mathcal{T}_{m}\right)$ the depth of insertion of the $k$ th phrase in tree $\mathcal{T}_{m}$. (Observe that $I_{k}\left(\mathcal{T}_{k}\right)=I_{k}\left(\mathcal{T}_{m}\right)$ ). If the tree $\mathcal{T}_{m}$ is known from the context, we often simplify the notation and write $I_{k}$.

We introduce now the so called tree-path. Let $w=x_{1} x_{2} \cdots x_{k}$ be a finite string whose length we also denote as $|w|=k$. We write $(w)_{i}$ for the prefix of $w$ of length $i$, that is, $(w)_{i}=x_{1} x_{2} \cdots x_{i}$. Assume now that $\mathcal{T}_{m}$ is given. The trec-path $C_{m}(w)$ associated with $w$ is a "trace" (path) in $\mathcal{T}_{m}$ when one follows symbols of $w$ along a path in the tree $\mathcal{T}_{m}$ until the paths split. More precisely:

Definition 2 The tree-path $C_{m}(w)$ associated with a given string $w$ in $\mathcal{T}_{m}$ is the largest integer $\ell \leq|w|$ such that there exist $k \leq m$ which satisfies (i) $(w)_{\ell}$ is the prefix of phrase $k$, and (ii) $I_{k}\left(\mathcal{T}_{m}\right)=\ell$. In other words, it is the number of nodes minus one contained in this path.

We now outline some properties of the distribution of the tree-path when DST is random. The next lemma shows that the tree-path distribution satisfies a simple recurrence.

Lemma 8 (i) Consider any model of phrase generations. Then, for all integers $m>1$

$$
\begin{align*}
\operatorname{Pr}\left\{C_{m}(w) \geq k\right\} & =\operatorname{Pr}\left\{C_{m-1}(w) \geq k\right\}+  \tag{88}\\
& +\operatorname{Pr}\left\{C_{m-1}(w)=k-1 \&(w)_{k} \text { is prefix of } m \text { th phrase }\right\}
\end{align*}
$$

for all $k \geq 0$.
(ii) If the strings are generated according to the MI model, then (88) becomes

$$
\begin{equation*}
\operatorname{Pr}\left\{C_{m}^{M I}(w) \geq k\right\}=\operatorname{Pr}\left\{C_{m-1}^{M I}(w) \geq k\right\}+\operatorname{Pr}\left\{C_{m-1}^{M I}(w)=k-1\right\} p_{x_{1}} p_{x_{1} x_{2}} \cdots p_{x_{k-1} x_{k}} \tag{89}
\end{equation*}
$$

where $\mathbf{p}=\left(p_{1}, \ldots, p_{V}\right)$ is the initial probability of generating the first symbol of the string $w=x_{1} \cdots x_{|w|}$.

Proof. To prove (88) we observe that the trec-path in $\mathcal{T}_{m}$ is greater equal than $k$ if and only if either it is greater equal than $k$ in $\mathcal{T}_{m-1}$ (i.e., the $\pi$ th insertion does not follow ( $\left.w\right)_{k}$ ) or the $m$ insertion traces the word $w$ up to $k-1$ and the $k$ th prefix of $w$ is a prefix of the $m$ th phrase.

We need a simple technical lemma whose proof requires pathwise comparison of two stochastic processes (trees).

Lemma 9 Let $w$ be a finite string. Consider two random DST trees $\mathcal{T}_{m_{1}}^{1}$ and $\mathcal{T}_{m_{2}}^{2}$ of respective size $m_{1}$ and $m_{2}$, with tree-paths $C_{m_{1}}^{1}(w)$ and $C_{m_{2}}^{2}(w)$, respectively. We assume that for all $w \in \mathcal{A}^{|w|}$

$$
C_{m_{1}}^{1}(w) \leq_{\mathrm{st}} C_{m_{2}}^{2}(w)
$$

If we insert to both trees the same independent phrase (string), then the corresponding tree paths $C_{m_{1}+1}^{1}(w)$ and $C_{m_{2}+1}^{2}(w)$ still satisfy

$$
C_{m_{1}+1}^{1}(w) \leq_{\mathrm{st}} C_{m_{2}+1}^{2}(w)
$$

for all $w$.
Proof. We remark that we cannot use Lemma 8 since there is no easy way of bounding $\operatorname{Pr}\left\{C_{m}(w)=k-1\right\}$. Thus, we shall rely on another approach, namely stochastic dominance in which the independence assumption plays a central role.

Let us fix a given string $w$. By the pathwise stochastic dominance theorem [19], there exists a probabilistic space on which a pair of DST trecs $\left(\tilde{\mathcal{T}}_{m_{1}}^{1}, \tilde{\mathcal{T}}_{m_{2}}^{2}\right)$ satisfies

- For $i=1,2$ the tree-path distribution of $\tilde{C}_{m_{i}}^{i}(w)$ on $\tilde{\mathcal{T}}_{m_{i}}^{i}$, is the same as the tree-path distribution of $C_{m_{i}}^{i}(w)$ on the original trees $\mathcal{T}_{m_{i}}^{i}$;
- $\tilde{C}_{m_{1}}^{1}(w) \leq \tilde{C}_{m_{2}}^{2}(w)$ for every randorn event.

Now, we insert into both trees $\tilde{\mathcal{T}}_{m_{1}}^{1}$ and $\tilde{\mathcal{T}}_{m_{2}}^{2}$ the same independent random phrase. The path distribution after insertion becomes, respectively, $\tilde{C}_{m_{1}+1}^{1}(w)$ and $\tilde{C}_{m_{2}+1}^{2}(w)$. It is easy to check via Lemma 8 that the distribution of $\bar{C}_{m_{i}+1}^{i}(w)$ will be the same as the distribution of $C_{m_{i}+1}^{i}(w)$. We consider two cases, namely: either $\widetilde{C}_{m_{1}}^{1}(w) \leq \widetilde{C}_{m_{2}}^{2}(w)-1$ or $\tilde{C}_{m_{1}}^{1}(w)=\tilde{C}_{m_{2}}^{2}(w)$ for every $w$. In the first case we must have $\widetilde{C}_{m_{1}+1}^{1}(w) \leq \widetilde{C}_{m_{2}+1}^{2}(w)$ after the insertion since the insertion of the new phrase can only increment by one unit the trec-path. In the second case, we also have $\widetilde{C}_{m_{1}+1}^{1}(w)=\tilde{C}_{m_{2}+1}^{2}(w)=k$ since the insertion of the new phrase can either increment by one unit the tree-pathes of $w$ on both trees or change nothing on both tree-paths, depending whether $(w)_{k}$ is the $k$ th length prefix of the new phrase.

In a typical application of this lemma, we shall assume that for any word $w$ and sizes $m_{1}$ and $m_{2}$ the following

$$
C_{m_{1}}^{G K}(w) \leq_{\mathrm{st}} C_{m_{2}}^{M I}
$$

implies

$$
C_{m_{1}+1}^{G K+M I}(w) \leq_{\mathrm{st}} C_{m 2+1}^{M I}
$$

where $C_{m_{1}+1}^{G K+M I}$ denotes the tree path in the GK model in which a new independent phrase is inserted.

Now, we are in a position to establish main results of this subsection, namely lower and upper bounds on the tree path. Let $C_{m}^{G K}(a w)$ and $C_{m}^{M I}(a w)$ denote the tree-paths in the GK and MI models, respectively, when the associated words $a w$ starts with a given symbol, say $a$. The following lemma presents an upper bound on $C_{m}^{G K}(a w)$ with respect to $C_{m 1}^{M I}(a w)$.

Lemma 10 The tree path $C_{m}^{G K}(a w)$ in the $G K$ model is stochastically bounded from the above by the tree path $C_{m}^{M I}(a w)$ in the $M I$ model in which all the $m$ phrases starts with symbol a (i.e. $\mathbf{p}=\mathbf{p}_{a}$ ), that is,

$$
\begin{equation*}
C_{m}^{G K}(a w) \leq_{\mathrm{st}} C_{m}^{M I}(a w) \tag{90}
\end{equation*}
$$

for all $w \in \mathcal{A}^{|w|}$ and $a \in \mathcal{A}$.
Proof. The proof is by induction on $m$. The property is true for $m=1$. We now suppose it is true for $m-1$. Let us consider the path $C_{m}^{G K}(a w)$ in the GK model. We obtain by Lemma 8

$$
\begin{gathered}
\operatorname{Pr}\left\{C_{m}^{G K}(a w) \geq k+1\right\}=\operatorname{Pr}\left\{C_{m-1}^{G K}(a w) \geq k+1\right\} \\
+\sum_{b=1}^{V} \operatorname{Pr}\left\{C_{m-1}^{G K}(a w)=k \&(m-1) \text { th phrase ends with } b\right\} p_{b a} p_{a x_{1}} p_{x_{1} x_{2}} \cdots p_{x_{k-1} x_{k}}
\end{gathered}
$$

Since $p_{b a} \leq 1$, and

$$
\sum_{b=1}^{b=V} \operatorname{Pr}\left\{C_{m-1}(a w)=k-1 \&(m-1) \text { th phrase ends with } b\right\}=\operatorname{Pr}\left\{C_{m-1}(a w)=k-1\right\}
$$

we obtain

$$
\begin{aligned}
\operatorname{Pr}\left\{C_{m}^{G K}(a w) \geq k+1\right\} & \leq \operatorname{Pr}\left\{C_{m-1}^{G K}(a w) \geq k+1\right\}+ \\
& +\operatorname{Pr}\left\{C_{m-1}^{G K}(a w)=k\right\} p_{a x_{1}} p_{x_{1} x_{2}} \cdots p_{x_{k-1} x_{k}} \\
& =\operatorname{Pr}\left\{C_{m}^{G K+M I}(a w) \geq k+1\right\}
\end{aligned}
$$

The last equality directly follows from Lemma 8 with $p_{a}=1$. Therefore $C_{m}^{G K}(a w) \leq_{s t}$ $C_{m}^{G K+M I}(a w)$. To complete the proof, we use the fact that

$$
\begin{equation*}
C_{m}^{G K+M I}(a w) \leq_{\mathrm{st}} C_{m}^{M I}(a w) \tag{91}
\end{equation*}
$$

which is a consequence of the induction hypothesis, $C_{m-1}^{G K}(a w) \leq_{s t} C_{m-1}^{M I}(a w)$ and Lemma 9. Indeed in both models, $G K+M I$ and $M I$, the last phrase is statistically independent of the $m-1$ first phrases and therefore meets the conditions of Lemma 9.

Finally, we derive a lower bound on the tree path in the GK model. Below, we shall write $r(a)=\min _{i}\left\{p_{i a}\right\}$ and $r=\sum_{a \in \mathcal{A}} \tau(a)$. We denote by $C_{m}^{M I B(r)}(w)$ the path length in the MI model with binomially $(m, r)$ distributed number of phrases. We denote $r$ the probability vector consisting of $\frac{r(a)}{r}$ for all $a \in \mathcal{A}$.

Lemma 11 The tree path $C_{m}^{G K}(w)$ in the GK model is stochastically bounded from the below by the tree path $C_{m-1}^{M I B(r)}(a w)$ in the MI model in which the first symbol of all phrases is distributed according to $\mathbf{r}$, and the number of phrases (strings) are generated according to the binomial distribution with parameters $m$ and $r<1$, that is,

$$
\begin{equation*}
C_{m-1}^{M I B(r)}(w) \leq_{\mathrm{st}} C_{m}^{G K}(w) \tag{92}
\end{equation*}
$$

Proof. The proof is by induction, and we shall imitate our proof of Lemma 10 with a few changes. The property is true for $m=2$, i.e., the second plrase starts with symbol $a$ with a probability smaller than $\tau(a)$ regardless of the actual value of the first phrase. We now suppose the property is truc for $m-1$ and let us take an arbitrary symbol $a \in \mathcal{A}$. We have

$$
\begin{aligned}
\operatorname{Pr}\left\{C_{m}^{G K}(a w)\right. & \geq k+1\}=\operatorname{Pr}\left\{C_{m-1}^{G K}(a w) \geq k+1\right\}+ \\
& +\sum_{b=1}^{v} \operatorname{Pr}\left\{C_{m-1}^{G K}(a w)=k-1 \&(m-1) \text { th phrase ends with } b\right\} \times \\
& \times p_{b a} p_{a x_{1}} p_{x_{1} x_{2}} \cdots p_{x_{k-1} x_{k}} \\
& \geq \operatorname{Pr}\left\{C_{m-1}^{G K}(a w) \geq k\right\}+\operatorname{Pr}\left\{C_{m-1}^{G K}(a w)=k-1\right\} \tau \times \frac{r(a)}{\tau} p_{x_{1} x_{2}} \cdots p_{x_{k-1} x_{k}} \\
& ={ }^{(A)} \operatorname{Pr}\left\{C_{m}^{G K+M I B(r)}(a w) \geq k+1\right\} \\
& \geq{ }^{(B)} \operatorname{Pr}\left\{C_{m-1}^{M I B(r)}(a w) \geq k+1\right\} .
\end{aligned}
$$

Equation (A) follows from Lemma 8 after noticing that the line above could be interpreted as the MI model in which the $m$ phrase is inserted with probability $r$ and the initial symbol of every phrase ahs distribution $\tau(a) / r$. The inequality (B) is a consequence of the induction assumption and Lemma 9. Observe that we omit the first phrase (so we have ( $m-1$ ) in the last line of the above) since it does not fall under our assumptions, i.e., its first symbol is not distributed according to r .

### 4.2 Bounds on the Phrase Length and Depth of Insertion

In this subsection, we translate the bounds on the tree path $C_{m}(w)$ into bounds of the depth of insertion $I_{m}$ in the GK model. We start with a simple observation which relates the depth
of insertion with the tree-path. We have

$$
\begin{aligned}
\operatorname{Pr}\left\{I_{m}=|w| \quad\right. & \& \quad w \text { is a prefix of the } m \text { th phrase }\} \\
& =\operatorname{Pr}\left\{C_{m-1}(w)=|w|-1 \& w \text { is a prefix of the } m \text { th phrase }\right\}
\end{aligned}
$$

which further implies

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{m} \geq k\right\}=\sum_{|w|=k} \operatorname{Pr}\left\{C_{m-1}(w) \geq k-1 \& w \text { is a prefix of the } m \text { th phrase }\right\} \tag{93}
\end{equation*}
$$

This and Lemma 9 lead immediately to the following useful result.

Lemma 12 Consider two random DST trees $\mathcal{T}_{m_{1}}^{1}$ and $\mathcal{T}_{m_{2}}^{2}$, of respective size $m_{1}$ and $m_{2}$, with tree-paths $C_{m_{1}}^{1}(w)$ and $C_{m_{2}}^{2}(w)$, and depths of insertion $I_{m_{1}}^{1}$ and $I_{r_{2}}^{2}$, respectively. If for all $w$

$$
C_{m_{1}}^{1}(w) \leq_{s t} C_{m_{2}}^{2}(w),
$$

then an independent phrase inserted into both trees leads to the following inequality

$$
I_{m_{1}+1}^{1} \leq \leq_{\mathrm{st}} I_{m_{2}+1}^{2} .
$$

on the depths of insertion.

Before we proceed with a formal derivation of the bounds on $I_{m}$, we present here a "guided tour" through the proof. The first step of establishing a bound for $I_{m}^{G K}$ in the GK model is to break a strong dependency between phrases so that the precise results of the MI model can be applied. We accomplish it by deleting the last $K$ phrases before inserting a new phrasc. We denote by $I_{m, K}^{G K}$ the depth of insertion in the GK model when $K$ last phrases are deleted. In order to make this idea useful, we need an inequality relating the depth $I_{m}^{G K}$ and the depth $I_{m, K}^{G K}$. But, in (37)) of Section 2 we prove such a result which we repeat here for reader's convenience

$$
\begin{equation*}
I_{m+1, K}^{G K} \leq I_{m+1}^{G K} \leq I_{m+1, K}^{G K}+K \tag{94}
\end{equation*}
$$

Unfortunately, we could not establish an easy bound on $I_{m, K}^{G K}$. But, in the previous section we proved a lower bound and an upper bound on the tree path that through Lemma 12 will lead to bound on $I_{m-K}^{G K+M I}$ where $I_{m-K}^{G K+M I}$ denotes the depth of insertion in the GK model when one inserts an independent phrase. The last step is to show that distributions of $I_{m, K}^{G K}$ and $I_{m-K}^{G K+M I}$ are within distance $\varepsilon_{m} \rightarrow 0$.

We start the analysis by showing that $I_{m, K}^{G K}$ is within distance $\varepsilon_{m} \rightarrow 0$ from $I_{m-K}^{G K+M I}$ which is crucial to our analysis.

Lemma 13 The random variable $I_{m, K}^{G K}$ is within distance $\varepsilon_{m}=O\left(m^{K \log \rho}\right)$ from $I_{m-K}^{G K+M I}$, where $\rho<1$ is the mixing coefficient of the underlying Markov process. (We shall use a short-hand notation $I_{m, K}^{G K} \stackrel{d}{=} I_{m-K}^{G K+M I}+O\left(\varepsilon_{m}\right)$ in such a situation.)

Proof. We shall use the fact that a Markov process over a finite space is a $\phi$-mixing process with exponentially decreasing mixing coefficient (cf. [1]). More precisely, let for some $d$ and $\ell$ two events, say $A$ and $B$, be defined on the sigma-algebras $\mathcal{F}_{-\infty}^{d}$ and $\mathcal{F}_{d+\ell}^{\infty}$, respectively (i.e., there is a gap of $\ell$ symbols between the events). Then, if the underlying process is Markov over a finite space, then there exists $\rho<1$ such that

$$
|\operatorname{Pr}\{A \& B\}-\operatorname{Pr}\{A\} \operatorname{Pr}\{B\}| \leq \rho^{\ell} \operatorname{Pr}\{A\} \operatorname{Pr}\{B\}
$$

We now associate the event $A$ with first $m-K-1$ phrases and the event $B$ with the $m$ th phrase. Actually, we consider $I_{m, K}^{G K}$ which can be viewed as event $A \& B$ while $I_{m-K}^{G K+M I}$ is composed of two independent events, $A$ and $B$. That is, if $\mathcal{E}_{\ell}$ denotes the event that $K$ last phrases are of length at least $\ell$ symbols, then for any set $D$ of integers

$$
\left|\operatorname{Pr}\left\{I_{m, K}^{G K} \in D \mid \mathcal{E}_{\ell}\right\}-\operatorname{Pr}\left\{I_{m-K}^{G K+M I} \in D \mid \mathcal{E}_{\ell}\right\}\right| \leq \rho^{\ell} \operatorname{Pr}\left\{I_{m-K}^{G K+M I} \in D \mid \mathcal{E}_{\ell}\right\}
$$

Let us now fix $\ell=K \beta \log m$ for some $\beta>0$. In Lemma 14 below we prove that $\operatorname{Pr}\left\{\right.$ not $\left.\mathcal{E}_{\ell}\right\}<$ $K \exp \left(-A m^{\alpha}\right)$ for $\alpha>0$, hence

$$
\left|\operatorname{Pr}\left\{I_{m, K}^{G K} \in D\right\}-\operatorname{Pr}\left\{I_{m-K}^{G K+M I I} \in D\right\}\right| \leq \varepsilon_{m}
$$

with $\varepsilon_{m}=\rho^{K \beta \log m}+K \exp \left(-A m^{\alpha}\right)=O\left(m^{\beta^{\prime} K \log \rho}\right)$ where $\beta^{\prime}>0$.
Lemma 14 There exist $A>0, \alpha>0$ and $\beta>0$ such that for all $m>0, \operatorname{Pr}\left\{I_{m}^{G K} \leq\right.$ $\beta \log m\} \leq \exp \left(-A m^{\alpha}\right)$

Proof. By (93) we have

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{m}^{G K} \geq k\right\} \geq 1-\sum_{|w|=k} \operatorname{Pr}\left\{C_{m-1}(w) \leq k-1\right\} . \tag{95}
\end{equation*}
$$

Thus, we need an estimate for $\operatorname{Pr}\left\{C_{m-1}(w) \leq k-1\right\}$. Observe that by Lemma 8

$$
\begin{aligned}
\operatorname{Pr}\left\{C_{m}(w)=k \mid C_{m-1}(w)=k-1\right\} & =\sum_{a \in \mathcal{A}} \operatorname{Pr}\{\text { last phrase ends with a }\} P\left(a(w)_{k}\right) \\
\operatorname{Pr}\left\{C_{m}(w)=k-1 \mid C_{m-1}(w)=k-1\right\} & =\sum_{a \in \mathcal{A}} \operatorname{Pr}\{\text { last phrase ends with a }\}\left(1-P\left(a(w)_{k-1}\right)\right)
\end{aligned}
$$

where $P(a w)$ denotes the probability of the string $a w$ induced by the underlying probabilistic model. Let now $\mu=\min _{a, b \in \mathcal{A}}\left\{p_{a b}\right\}>0$. Then, the above implies

$$
\begin{aligned}
\operatorname{Pr}\left\{C_{m}(w)=k \mid C_{m-1}(w)=k-1\right\} & \leq P\left(a(w)_{k}\right) \leq 1 \\
\operatorname{Pr}\left\{C_{m}(w)=k-1 \mid C_{m-1}(w)=k-1\right\} & \leq 1-\mu^{k+1}
\end{aligned}
$$

But then $\operatorname{Pr}\left\{C_{m}(w)=k\right) \leq\binom{ m}{k}\left(1-\mu^{k+1}\right)^{m-k}$, and hence

$$
\operatorname{Pr}\left\{C_{m}(w) \leq k\right\} \leq k\binom{m}{k}\left(1-\mu^{k+1}\right)^{m-k} \leq k\binom{m}{k} \exp \left(-\mu^{k}(m-k)\right)
$$

Set now $k=\left\lceil-\frac{\log m}{2 \log \mu}\right\rceil$. Since $\binom{m}{k} \leq \frac{m^{k}}{k!}$, the above becomes

$$
\operatorname{Pr}\left\{C_{m}(w) \leq k\right\} \leq k\binom{m}{k} \exp \left(-\mu^{k}(m-k)\right)=\exp (-\eta \sqrt{m})
$$

where $\eta>0$ is a constant. Finally, returning to (95) with $k=\left\lceil-\frac{\log m}{2 \log \mu}\right\rceil$ and noticing that in this case $\sum_{|w|=k} 1 \leq m^{B}$ for some $B>0$, we obtain

$$
\operatorname{Pr}\left\{I_{m}^{G K} \geq k\right\} \geq 1-m^{B} \exp (-\eta \sqrt{m})
$$

which completes the proof.
Finally, we are in a position to establish an upper bound (cf. Theorem 4) and a lower bound (cf. Theorem 5) on the depth of insertion $I_{m}^{G K}$.

Theorem 4 Let $I_{m-K}^{G K}(a)$ be the depth of insertion in the GK model when the mth phrase starts with symbol a, and $I_{m-K}^{M I}\left(\mathbf{p}_{a}\right)$ be the depth of insertion in the MI model with the initial probability vector $\mathbf{p}_{a}=(0, \ldots, 1, \ldots, 0)$ where 1 is at position $a \in \mathcal{A}$ (i.e., all strings start with symbol $a)$. Then, for any $\beta>0$, there exists $K$ such that $I_{m}^{G K}(a)$ is stochastically dominated by a random variable that is within distance $O\left(n^{-\beta}\right)$ from $I_{m-K}^{M I}\left(\mathrm{p}_{a}\right)+K$

Proof. Let $K$ be a fixed integer. We have from (94)

$$
I_{m}^{G K}(a) \leq I_{m, K}^{G K}(a)+K
$$

We also have

$$
I_{m, K}^{G K}(a) \stackrel{d}{=} I_{m-K}^{G I+M I}(a)+O\left(\varepsilon_{m}\right)
$$

as a consequence of Lemma 13. Lemma 10 implies

$$
I_{m-K}^{G I+N I}(a) \leq_{\mathrm{st}} I_{m-K}^{M I}\left(\mathbf{p}_{a}\right)
$$

which completes the proof.
The proof of a lower bound on $I_{m}^{G K}$ follows the same footsteps as above, so we only sketch it here. As before, we shall write $I_{m}^{M I B(r)}(\mathrm{r})$ for the depth of inscrtion in the MI model in which first symbol in each phrase distributes according to vector $r$ and the number of phrases is distributed according the the binomial $(m, r)$ for some $r<1$ The probability $r$ and the probability vector $\mathbf{r}$ are defined above Lemma 11.

Theorem 5 For any $\beta>0$, there exists $K$ such that $I_{m}^{G K}(a)$ stochastically dominates a random variable that is within distance $O\left(n^{-\beta}\right)$ from $I_{m-K}^{M I B(r)}(r)$ for some $r<1$.

Proof. We have the following chain of inequalities

$$
I_{m}^{G K}(a) \geq I_{m, K}^{G K}(a) \stackrel{d}{=} I_{m-K}^{G I+M I}(a)+O\left(\varepsilon_{m}\right) \succeq_{\mathrm{st}} I_{m-K}^{M I B(r)}(\mathbf{r})
$$

which completes the proof.

### 4.3 Establishing the Limiting Distribution

We prove now that appropriately normalized $I_{m}^{G K}$ converges in distribution to the standard normal distribution. Similar conclusion about the typical depth $D_{m}^{G K}$ will follow directly via the Cesaro limit.

To simplify notation, let $L_{m}=\frac{\log m}{h}$ and $V_{m}=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{Q}^{*} \psi-h^{2}\right) \ln m$. Wc will prove that for all $x=O(1)$

$$
\lim _{m \rightarrow \infty} \operatorname{Pr}\left\{\frac{I_{m}^{G K}-L_{m}}{\sqrt{V_{m}}} \geq x\right\}=\frac{1}{\sqrt{2 \pi}} \int_{x}^{\infty} e^{-t^{2} / 2} d t
$$

By Theorem 4, there exist $\beta>0$ and $K$ such that the following upper bound holds for all $k$ and $m$ :

$$
\begin{equation*}
\operatorname{Pr}\left\{I_{m}^{G K} \geq k \mid \text { last phrase starts with } a\right\} \leq \operatorname{Pr}\left\{I_{m-K}^{M I}\left(\mathbf{p}_{a}\right) \geq k-K\right\}+O\left(n^{-\beta}\right) \tag{96}
\end{equation*}
$$

Thus,

$$
\begin{aligned}
\operatorname{Pr}\left\{I_{m}^{G K} \geq k\right\} & =\sum_{a \in \mathcal{A}} \operatorname{Pr}\left\{I_{m}^{G K} \geq k \mid \text { last } G K \text { phrase starts with } a\right\} \\
& \times \operatorname{Pr}\{\text { last } G K \text { phrase starts with } a\} \\
& \leq \sum_{a \in \mathcal{A}} \operatorname{Pr}\left\{I_{m-K}^{M I}\left(\mathbf{p}_{a}\right) \geq k-K\right\} \operatorname{Pr}\{\text { last } G K \text { phrase starts with } a\}+O\left(n^{-\beta}\right)
\end{aligned}
$$

By Corollary 1 we know that

$$
\lim _{m \rightarrow \infty} \operatorname{Pr}\left\{\frac{I_{m}^{M I}\left(\mathbf{p}_{a}\right)-L_{m}}{\sqrt{V_{m}}} \geq x\right\}=\frac{1}{\sqrt{2 \pi}} \int_{x}^{\infty} e^{-t^{2} / 2} d t
$$

Now, observe that $L_{m-K}=L_{m}+O(1 / m)$ and $V_{m-K}=V_{m}+O(1 / m)$. Also, obviously $\sum_{a \in \mathcal{A}} \operatorname{Pr}\{$ last $G K$ phrase starts with $a\}=1$, which finally leads to

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \operatorname{Pr}\left\{\frac{I_{m}^{G K}-L_{m}}{\sqrt{V_{m}}} \geq x\right\} \leq \lim _{m \rightarrow \infty} \frac{1}{\sqrt{2 \pi}} \int_{x-O(1 / m)}^{\infty} e^{-t^{2} / 2} d t=\frac{1}{\sqrt{2 \pi}} \int_{x}^{\infty} e^{-t^{2} / 2} d t \tag{97}
\end{equation*}
$$

A similar argument works for the lower bound, however, this time we shall use Theorem 5 and Corollary 2 and we do not need to split over symbols $a$. Certainly,

$$
\operatorname{Pr}\left\{I_{m}^{G K} \geq k\right\} \geq \operatorname{Pr}\left\{I_{m-K}^{M I B(r)}(\mathbf{r}) \geq k\right\}+O\left(n^{-\beta}\right)
$$

But, by Corollary 2, $\left(I_{m}^{M I B(r)}\left(\mathbf{p}_{a}\right)-L_{m}\right) / V_{m} \xrightarrow{d} N(0,1)$, hence by a similar line of reasoning as above we conclude that

$$
\liminf _{m \rightarrow \infty} \operatorname{Pr}\left\{\frac{I_{m}^{G K}-L_{m}}{\sqrt{V_{m}}} \geq x\right\} \geq \frac{1}{\sqrt{2 \pi}} \int_{x}^{\infty} e^{-t^{2} / 2} d t
$$

which completes the proof for the limiting distribution of $I_{m}^{G K}$.

### 4.4 Establishing the Convergence of Moments

Finally, we prove the existence and convergence of moments of $\left(I_{m}^{G K}-L_{m}\right) / \sqrt{V_{n}}$ where, as before, $L_{r n}=\frac{\log m}{h}$ and $V_{m}=\frac{1}{h^{3}}\left(-\frac{\beta}{\omega}-\frac{2}{\omega} \pi \dot{Q}^{*} \psi-h^{2}\right) \ln m$. We accomplish this by showing that there exist constants $A_{1}$ and $\alpha_{1}<1$ such that uniformly for all integers $\ell$

$$
\begin{equation*}
\operatorname{Pr}\left\{\left|\frac{I_{m}^{G K}-L_{m}}{\sqrt{V_{m}^{-}}}\right| \geq \ell\right\} \leq A_{1} \alpha_{1}^{\sqrt{\ell}} \tag{98}
\end{equation*}
$$

Indeed, above will prove the existence of the moments exist and by uniform dominated theorem their convergence to the moments of the normal distribution. Notice that in any model $I_{m}$ cannot be greater than $m$ and therefore there is no need to check the incquality for values of $\ell$ beyond $m$.

Below, we present details of the derivations only for the casc $\operatorname{Pr}\left\{I_{m}^{G K}-L_{m} \geq \ell \sqrt{V_{m}}\right\}$ since the case $\operatorname{Pr}\left\{I_{m}^{G K}-L_{m} \leq-\ell \sqrt{V_{m}}\right\}$ can be handled in a similar manner.

By (94) we know that $I_{m}^{G K} \leq I_{m, K}^{G K}+K$ for a fixed $K$. But, Lemma 13 asserts that $I_{m, K}^{G K}$ is within distance $\varepsilon_{m}=O\left(m^{K \log \rho}\right)$, where $\rho<1$, from $I_{m-K}^{G K+M I}$. More precisely, for any set of integers $B$

$$
\operatorname{Pr}\left\{I_{m, K}^{G K} \in B\right\} \leq\left(1+\varepsilon_{m}\right) \operatorname{Pr}\left\{I_{m-K}^{G K+M I} \in B\right\}+O\left(e^{-\eta \sqrt{m}}\right)
$$

for $\eta>0$. From Theorem 4 we know also that

$$
I_{m-K}^{G K+M I}(a) \leq_{\mathrm{s} t} I_{m-K}^{M I}\left(\mathbf{p}_{a}\right)
$$

where above we indicated that phrases starts with symbol a. Finally, Corollary 1 implies that there are constants $A$ and $\alpha<1$ such that

$$
\operatorname{Pr}\left\{\left|\frac{\left.I_{m}^{M I}\left(\mathbf{p}_{a}\right)-L_{m}\right]}{\sqrt{V_{m}}}\right| \geq \ell\right\} \leq A \alpha^{\ell}
$$

Putting everything together, we obtain

$$
\begin{aligned}
\operatorname{Pr}\left\{I_{m}^{G K} \geq L_{m}+\ell \sqrt{V_{m}}\right\} & \leq\left(1+\varepsilon_{m}\right) \sum_{a \in \mathcal{A}} \operatorname{Pr}\left\{I_{m-K}^{M I}\left(\mathbf{p}_{a}\right) \geq k-K\right\} \\
& \times \operatorname{Pr}\{\text { last } G K \text { phrase starts with } a\}+O\left(e^{-\eta \sqrt{m}}\right) \\
& \leq A\left(1+\varepsilon_{m}\right) \alpha^{\ell}+O\left(e^{-\eta \sqrt{m}}\right) \leq A_{1} \alpha_{1}^{\sqrt{\ell}}
\end{aligned}
$$

since $\ell$ cannot be greater than $m$ and therefore $O\left(e^{-\eta \sqrt{m}}\right)$ can be dominated by $A_{1} \alpha_{1}^{\sqrt{\ell}}$ term. This prove the existence and convergence of moments which completes the proof of Theorem 2.

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## Appendix A: Alternative Representation of Theorem 1 Results

In this appendix, we show how to prove our alternative representations (19)-(20) for the mean $\mathbf{E}\left[D_{m}\right]$ and $\operatorname{Var}\left[D_{m}\right]$. Instead of presenting a detailed derivations, as in Section 3, we rather sketch here the proof.

We concentrate on evaluating the mean. The starting point is (62), that is,

$$
\mathbf{x}(s)=\mathrm{Q}^{-1} \mathbf{x}(s-1)=\sum_{k=0}^{\infty} \mathrm{P}^{k}(s) \mathbf{x}(s-1)
$$

Before we apply the spectral representation to $\mathrm{P}^{k}(s)$, we need some notation. Let us denote by $\lambda(s), \mu_{2}(s), \ldots, \mu_{V}(s)$ the eigenvalues of $\mathrm{P}(s)$ with $|\lambda(s)|>\left|\mu_{1}(s)\right| \geq \cdots \geq\left|\mu_{V}(s)\right|$. The correspondinging left eigenvectors are $\pi(s), \pi_{2}(s), \ldots, \pi_{V}(s)$ while the right eigenvectors are $\psi(s), \psi_{2}(s), \ldots, \psi_{V}(s)$. As in [6], we adopt an optional notation for the scalar product of vectors, namely, we either write as before $x y$ for product of vectors $x$ and $y$ or $\langle x, y\rangle$. The latter notation is convenient when scalar products are often used, as in this appendix.

By spectral representation (cf. [16]), matrix $\mathrm{P}(s)$ can be represented as

$$
\mathrm{P}^{k}(s) \mathrm{x}(s-1)=\lambda^{k}(s)\langle\pi(s), \mathrm{x}(s-1)\rangle \psi(s)+\sum_{i=2}^{V} \mu_{i}^{k}(s)\left\langle\pi_{i}(s), \mathrm{x}(s-1)\right\rangle \psi_{i}(s) .
$$

Thus $\mathbf{b}(s)=\Gamma(s) \mathbf{x}(s)$ becomes

$$
\begin{equation*}
\mathrm{b}(s)=\frac{\Gamma(s)\langle\pi(s), \mathbf{x}(s-1)\rangle \psi(s)}{1-\lambda(s)}+\sum_{i=2}^{V} \frac{\Gamma(s)\left\langle\pi_{i}(s), \mathbf{x}(s-1)\right\rangle \psi_{i}(s)}{1-\mu_{i}(s)} \tag{99}
\end{equation*}
$$

In order to obtain leading asymptotics of $B^{*}(s)=\mathbf{p}(s) \mathbf{b}(s)+\Gamma(s) x(s-1)$ (cf. (64)), we need Laurent's expansion of the above around the roots of $\lambda(s)=-1$. Observe that the second term of (99) contributed $o(m)$ since $\lambda(s)$ is the largest eigenvalue (cf. [6]), hence we further ignore this negligible term in our derivations. To simplify the presentation, we only deal here with the root $s_{0}=-1$. We use our previous expansions for $\mathrm{x}(s-1)$ and $\Gamma(s)$ together with

$$
\begin{aligned}
\frac{1}{1-\lambda(s)} & =\frac{-1}{\dot{\lambda}(-1)} \frac{1}{s+1}+\frac{\ddot{\lambda}(-1)}{2 \dot{\lambda}^{2}(-1)}+O(s+1) \\
\psi(s) & =\psi+\dot{\psi}(-1)(s+1)+O\left((s+1)^{2}\right)
\end{aligned}
$$

This finally leads to

$$
\begin{aligned}
B^{\prime}(s) & =\frac{-1}{\dot{\lambda}(-1)} \frac{1}{(s+1)^{2}} \\
& +\frac{1}{s+1}\left(\frac{\langle\pi, \dot{\mathbf{x}}(-2)\rangle}{\dot{\lambda}(-1)}-\frac{\gamma-1}{\dot{\lambda}(-1)}+\frac{\langle\mathrm{p}(-1), \dot{\psi}(-1)\rangle}{\dot{\lambda}(-1)}+\frac{\ddot{\lambda}(-1)}{2 \dot{\lambda}^{2}(-1)}-1\right)+O(1) .
\end{aligned}
$$

After finding the inverse Mellin transform of the above and depoissonizing, we prove the alternative representation (19).

Finally, we turn our attention to the second factorial moment and the variance. We need to study $\mathbf{c}(s)=\Gamma(s) \mathbf{v}(s)$ where $\mathbf{v}(s)=2 \mathrm{Q}^{-1}(s) \mathrm{P}(s) \mathbf{x}(s)+\mathrm{Q}^{-1}(s) \mathbf{v}(s-1)$. As before, we
obtain

$$
\mathbf{c}(s)=\frac{2 \Gamma(s)\langle\pi(s), \mathbf{x}(s-1)\rangle(\pi(s), \mathrm{P}(s) \psi(s)\rangle \psi(s)}{(1-\lambda(s))^{2}}+O\left((1-\lambda(s))^{-1}\right) .
$$

Similar algebra as above leads to

$$
\begin{aligned}
\mathbf{c}(s) & =\frac{-2}{\dot{\lambda}^{2}(-1)} \frac{1}{(s+1)^{3}} \\
& +\frac{1}{(s+1)^{2}}\left(\frac{\ddot{\lambda}(-1)}{2 \dot{\lambda}^{3}(-1)}+2 \frac{\gamma-1-\langle\pi, \dot{\mathrm{x}}(-2)\rangle-\langle\mathrm{p}(-1), \dot{\psi}(-1)\rangle-\dot{\lambda}(-1)}{\dot{\lambda}^{2}(-1)}\right) \\
& +O\left(\frac{1}{s+1}\right) .
\end{aligned}
$$

This is sufficient to prove (20), after some tedious algebra that was helped by Maple.


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