# A Non-Adaptive Algorithm for the Quantitative Group Testing Problem 

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#### Abstract

Consider an $n$ dimensional binary feature vector with $k$ non-zero entries. The vector can be interpreted as the incident vector corresponding to $n$ items out of which $k$ items are defective. The quantitative group testing (QGT) problem aims at learning this binary feature vector by queries on subsets of the items that return the total number of defective items. We consider this problem under the non-adaptive scenario where the queries on subsets are designed collectively and can be executed in parallel. Most of the existing efficient non-adaptive algorithms for the sublinear regime where $k=n^{\alpha}$ with $0<\alpha<1$ fall short of the information-theoretic lower bound, with a multiplicative gap of $\log k$. Recently, Hahn-Klimroth and Müller (2022) closed this gap by providing a non-adaptive algorithm with decoding complexity of $O\left(n^{3}\right)$.

In this work, we present a concatenated construction method yielding a non-adaptive algorithm with the decoding complexity of $O\left(n^{2 \alpha}+n \log ^{2} n\right)$. The probability of decoding failure is analyzed by establishing a connection between the QGT problem and the so-called balls into bins problem. Our algorithm reduces the gap between the information-theoretic and computational bound for the number of required queries/tests from $\log k$ to $\log \log k$. This narrows the gap in number of tests for non-adaptive algorithms within the class of algorithms with $o\left(n^{2}\right)$ decoding complexity. Moreover, although our algorithm exhibits a $\log \log k$ gap in terms of the number of tests, it is surpassed by the existing asymptotically optimal construction only in scenarios where $k$ is exceptionally large for moderate values of $\alpha$, such as $k>10^{27}$ for $\alpha=0.7$, thereby highlighting the practical applicability of our proposed concatenated construction.


Keywords: Statistical inference, Quantitative group testing, Urn models, Compressed sensing

## 1. Introduction

Quantitative Group Testing (QGT) is the problem of detecting $k$ defective items within a collection of $n$ items through a series of tests conducted on $m$ distinct pools. The result of each individual test yields the number of defective items present in the corresponding pool. This problem finds applications across a diverse spectrum of domains, including identifying rare variant carriers in genome sequencing (Cao et al., 2014), network traffic control (Wang et al., 2015), resource allocation in random access channels (De Marco et al., 2021), and signal recovery (Matsumoto et al., 2023; Mazumdar and Pal, 2021).

The QGT problem is often studied within two primary statistical models that define the arrangement of underlying items. In the probabilistic model, each item has a probability $p$ of being defective, which is independent of other items. In the combinatorial model, the number of defective items, denoted as $k$, is known in advance. Strategies for addressing the QGT problem fall into
two distinct categories: adaptive and non-adaptive. In the non-adaptive approach, all tests must be predetermined and can be executed in parallel. In contrast, in the adaptive approach, one can observe the results of previous tests and utilize this information to design subsequent tests . In particular, Bshouty (2009) provides an order-optimal adaptive algorithm that is within a factor of 2 of the information-theoretic lower bound. Recently, Soleymani and Javidi (2024) introduced an adaptive algorithm with tunable adaptation, allowing for a trade-off between the number of tests and the number of adaptive stages. In this paper, we focus on the non-adaptive QGT problem within the combinatorial model.

Developing an algorithm for non-adaptive QGT can be considered as the construction of a binary matrix, where the rows represent the tests and the columns correspond to the items. An item is included in a pool test if the corresponding entry in the associated row is equal to 1 . This setup bears a significant resemblance to the compressive sensing (CS) (Candès et al., 2006; Donoho, 2006) and sparse recovery problems that are studied extensively in the literature. In particular, the QGT problem can be regraded as a binary CS problem, with the additional constraint that the measurement matrix must also be binary.

In this paper, we investigate the QGT problem within the sublinear regime, where $k=n^{\alpha}$ for $0<\alpha<1$. The information-theoretic lower bound on the number of tests required in a non-adaptive scheme for the QGT problem in this regime is (Djackov, 1975)

$$
\begin{equation*}
m_{0} \stackrel{\text { def }}{=} \frac{2 k}{\log k} \log \frac{n}{k}, \tag{1}
\end{equation*}
$$

Since this problem can be regarded as a special case of the CS problem, utilizing the linear programming methods proposed recovers the underlying binary vector w.h.p. using $\Theta(k \ln n)$. Several studies in the literature addressed the construction of techniques tailored to the QGT problem (Karimi et al., 2019b,a; Coja-Oghlan et al., 2020; Feige and Lellouche, 2020; Gebhard et al., 2022; Soleymani et al., 2023; Tan et al., 2023). All such algorithms contribute to a reduction in the required number of tests compared to linear programming techniques, primarily focusing on improving multiplicative constant factors and the gap to the information theoretic gaps remains at $O(\log k)$.

Recently, Hahn-Klimroth and Müller (2022); Hahn-Klimroth et al. (2023) successfully bridged this gap by introducing a non-adaptive algorithm that infers all $k$ defective items w.h.p. with the number of tests approaching $\frac{1+\sqrt{\alpha}}{1-\sqrt{\alpha}} m_{0}$ as $n$ grows large and a decoding complexity of $O\left(n^{3}\right)$. Notably, this marks the first polynomial-time algorithm achieving an orderwise optimal number of tests for recovering defective items in the QGT problem. Our main focus in this study is to continue the effort of reducing the decoding complexity of non-adaptive algorithms for the QGT problem by addressing the fundamental question of whether it is possible to narrow the mentioned $O(\log k)$ gap while achieving decoding complexity superior to $O\left(n^{3}\right)$, ideally approaching near-linear complexity. To this end, we provide a construction method that reduces the $O(\log k)$ gap to $O(\log \log k)$ in the number of required tests accompanied by a computational complexity of $O\left(n^{2 \alpha}+n \log n\right)$. Our approach involves the introduction of a concatenation method that can be utilized to tackle the QGT problem through a divide-and-conquer strategy. The theoretical guarantee on the probability of successful detection of the defective items is established by introducing a certain probabilistic urn model (Sprott, 1978) and making a connection to the well-known balls into bins problem (Park, 1980). Our contribution can be summarized in the following:

- Decoding Complexity: In the sublinear regime, for $0<\alpha<\frac{1}{2}$, the decoding complexity of our algorithm is $O(n \log n)$, achieving an almost linear decoding complexity while narrowing
the computational-semantic gap from $\log k$ to $\log \log k$ in number of tests. For $\frac{1}{2}<\alpha<1$, the complexity is $o\left(n^{2 \alpha}\right)$, improving upon the state-of-the-art complexity of $O\left(n^{3}\right)$.
- Non-asymptotic behavior: Our proposed algorithm outperforms the algorithm proposed by Hahn-Klimroth and Müller (2022) in the non-asymptotic regime in terms of the number of tests for moderate values of $\alpha$. In other words, in order for the mentioned algorithm to surpass ours in terms of the number of tests, the parameter $n$ must be significantly large.

The rest of the paper is organized as follows. Section 2 introduces the system model and presents a concatenation approach that is utilized in our proposed test matrix design. Section 3 provides an overview of a result regarding the balls into bins problem and introduces a novel probabilistic urn model tailored to the QGT setting. We establish a connection between these two models, allowing us to utilize the probability bounds derived from the former for the latter. In Section 4, we discuss two fundamental components that can serve as the building blocks for our concatenated scheme, and provide our construction and analyse its performance. Finally, the paper is concluded in Section 5.

## 2. Problem Setting and a Concatenated Construction

### 2.1. Notation

The vectors and matrices are represented by bold lower and upper case characters, respectively. The component $i$ of a vector $\boldsymbol{a}$ is represented by $a_{i}$. The number of non-zero elements in $\boldsymbol{a}$ is referred to as the Hamming weight of $\boldsymbol{a}$ and is denoted by $|\boldsymbol{a}|$. The Hamming distance between two vectors $\boldsymbol{a}$ and $\boldsymbol{b}$ is $|\boldsymbol{a}-\boldsymbol{b}|$. Throughout this paper, $\log (\cdot)$ denotes the base 2 logarithm, while $\ln (\cdot)$ represents the natural logarithm with base $e$. The set of all non-negative integers is denoted by $\mathbb{N}_{0}$.

### 2.2. Problem Setting

The problem of quantitative group testing (QGT) is the following: A set of $n$ items is given out of which $k$ items are defective. The incident vector corresponding to these items is a binary vector $x \in\{0,1\}^{n}$ such that $x_{i}=1$ if the item is defective and $\boldsymbol{x}_{i}=0$, otherwise. Let $\mathscr{B}_{n, k}$ denote the set of all binary vectors of size $n$ with $k$ non-zero elements. In this paper, we consider the combinatorial model for the QGT problem where $x$ is picked from $\mathscr{B}_{n, k}$ uniformly at random. This problem is reffered to as $(n, k)$-QGT problem. The goal is to identify all defective items, i.e., recover $x$, by performing tests/measurements over as few subsets of the items as possible. The outcome of a test is the number of defective items belonging to the underlying subset. This problem is also referred to as the coin weighing problem in the literature where one attempts to detect $k$ counterfeit coins from a collection of $n$ coins using a spring scale, given that all counterfeit coins weight the same, which is different from the weight of genuine coins. Designing a non-adaptive strategy for the $(n, k)$-QGT problem is equivalent to constructing a binary matrix with $n$ columns such that the sum of any subset of size $k$ of the columns is distinct. This matrix is referred to as a search matrix, which is defined as follows.

Definition 1 A binary matrix $\boldsymbol{A} \in\{0,1\}^{m \times n}$ is called an $(n, k)$-search matrix if $\boldsymbol{A x} \neq \boldsymbol{A} \boldsymbol{x}^{\prime}$ for all $x \neq x^{\prime} \in \mathscr{B}_{n, k}$.

Next, Definition 1 is extended to to the case where the entries of $x$ are non-negative integers.

Definition 2 A binary matrix $\boldsymbol{D} \in\{0,1\}^{m \times n}$ is called $a\left(d_{1}, \cdots, d_{n}\right)$-detecting matrix if $\boldsymbol{D} \boldsymbol{x} \neq \boldsymbol{D} \boldsymbol{x}^{\prime}$ for all $\boldsymbol{x} \neq \boldsymbol{x}^{\prime} \in\left\{\boldsymbol{a}: 0 \leqslant \boldsymbol{a}_{i}<d_{i}, \forall i \in[n]\right\}$. In particular, if $d_{i}=d$ for all $i \in[n]$, the matrix is referred to as an $(n, d)$-detecting matrix.

We leverage an off-the-shelf construction for the $(n, d)$-detecting matrix, in particular the one introduced by Bshouty (2009). This serves as a building block in constructing measurement matrices for the QGT problem with balanced incident vectors. The notion of balanced vector is defined next.

Definition 3 A binary vector $x$ of length $n$ is called an $(m, t)$-balanced vector for some integer $m$ that divides $n$ if

$$
\max \left(\sum_{i=1}^{\frac{n}{m}} x_{i}, \sum_{i=\frac{n}{m}+1}^{2 \frac{n}{m}} x_{i}, \cdots, \sum_{i=n-\frac{n}{m}+1}^{n} x_{i}\right) \leqslant t
$$

In other words, in an ( $m, t$ )-balanced vector $\boldsymbol{x}$, the number of ones inside each block of length $\frac{n}{m}$ does not exceed $t$.

### 2.3. Concatenated Construction

One of the main ideas of the construction provided in this paper involves reducing a QGT problem into carefully selected smaller QGT subproblems. These subproblems are then tackled using a test matrix tailored for solving the smaller instances, in conjunction with a detection matrix with precisely crafted parameters. To establish this, the following theorem demonstrates that the Kronecker product of an $\left(n_{1}, d\right)$-detecting matrix with an $\left(n_{2}, d\right)$-search matrix can uniquely recover any vector $x$ of length $n_{1} n_{2}$ that is $\left(n_{1}, d\right)$-balanced.

Theorem 4 Let $\boldsymbol{A} \in\{0,1\}^{m_{1} \times n_{1}}$ and $\boldsymbol{B} \in\{0,1\}^{m_{2} \times n_{2}}$ be an $\left(n_{1}, d\right)$-detecting matrix and an $\left(n_{2}, d\right)$ search matrix, respectively. Let

$$
C \xlongequal[=]{\text { def }} \boldsymbol{A} \otimes B .
$$

Then, one can uniquely recover the binary vector $\boldsymbol{x}$ of length $n_{1} n_{2}$ from $C x$ if $\boldsymbol{x}$ is $\left(n_{1}, d\right)$-balanced.
Moreover, if there exist algorithms $\mathscr{A}$ and $\mathscr{B}$ that retrieve $\boldsymbol{a}$ and $\boldsymbol{a}^{\prime}$ from $\boldsymbol{A} a$ and $\boldsymbol{B a}$ within computation complexities of $O_{\mathscr{A}}$ and $O_{\mathscr{B}}$, respectively, then one can recover $\boldsymbol{x}$ from Cx with $a$ computation complexity of $m_{2} O_{\mathscr{A}}+n_{1} O_{\mathscr{B}}$.

Proof: Suppose $x$ is a binary vector of length $n_{1} n_{2}$ that is $\left(n_{1}, d\right)$-balanced. That is, we have $\sum_{i=(j-1) n_{2}+1}^{j n_{2}} x_{i} \leqslant d$ for all $j \in\left[n_{1}\right]$. We prove the statement of the theorem by proposing a decoder that recovers $\boldsymbol{x}$ from $\boldsymbol{C} \boldsymbol{x}$. Suppose that there exist algorithms $\mathscr{A}: \mathbb{N}_{0}^{m_{1}} \rightarrow\{0,1, \cdots, d\}^{n_{1}}$ and $\mathscr{B}: \mathbb{N}_{0}^{m_{2}} \rightarrow\{0,1\}^{n_{2}}$ that recover $\boldsymbol{a}$ and $\boldsymbol{a}^{\prime}$ from $\boldsymbol{A} \boldsymbol{a}$ and $\boldsymbol{B} \boldsymbol{a}^{\prime}$ within time complexities of $O_{1}$ and $O_{2}$, respectively. Let $\boldsymbol{b} \xlongequal{\text { def }} \boldsymbol{C} \boldsymbol{x}$ denote the corresponding measurement vector. We use the so-called mixed Kronecker matrix-vector product property as follows:

$$
\begin{equation*}
\operatorname{vec}\left(\boldsymbol{B} \boldsymbol{X} \boldsymbol{A}^{T}\right)=(\boldsymbol{A} \otimes \boldsymbol{B}) \operatorname{vec}(\boldsymbol{X}) \tag{2}
\end{equation*}
$$

for any $n_{2} \times n_{1}$ matrix $\boldsymbol{X}$, where the $\operatorname{vec}(\cdot)$ operator constructs a column vector from the input matrix $\boldsymbol{T}=\left[\boldsymbol{t}_{1}|\cdots| \boldsymbol{t}_{m_{1}}\right]$ by vertically stacking the column vectors of $\boldsymbol{T}$ beneath each other, i.e.,
$\operatorname{vec}(\boldsymbol{T}) \stackrel{\text { def }}{=}\left[\boldsymbol{t}_{1}^{t}, \cdots, \boldsymbol{t}_{m_{1}}^{t}\right]^{t}$. We partition $\boldsymbol{x}$ into $n_{1}$ column vectors of length $n_{2}$ and arrange them as the columns of a new $n_{2} \times n_{1}$ matrix referred to as $\boldsymbol{X}$, i.e., $\boldsymbol{X} \xlongequal{\text { def }}\left[x_{1}|\cdots| x_{n_{1}}\right]$. Following this, it can be verified that $\operatorname{vec}(\boldsymbol{X})=\boldsymbol{x}$. Then, one can write

$$
\begin{equation*}
\boldsymbol{F}_{m_{2} \times n_{1}} \stackrel{\text { def }}{=} \boldsymbol{B}_{m_{2} \times n_{2}} \boldsymbol{X}_{n_{2} \times n_{1}}=\left[\boldsymbol{B} \boldsymbol{x}_{1}\left|\boldsymbol{B} \boldsymbol{x}_{2}\right| \cdots \mid \boldsymbol{B} \boldsymbol{x}_{n_{1}}\right] . \tag{3}
\end{equation*}
$$

Let $\boldsymbol{F}_{n_{1} \times m_{2}}^{T}=\left[\tilde{f}_{1}\left|\tilde{f}_{2}\right| \cdots \mid \tilde{f}_{m_{2}}\right]$. Note that $\tilde{f}_{i}$ 's represent the columns of $\boldsymbol{F}^{T}$, hence they are column vectors of length $n_{1}$. Then, we have

$$
\begin{align*}
\operatorname{vec}\left(\boldsymbol{B} \boldsymbol{X} \boldsymbol{A}^{T}\right) & =\operatorname{vec}\left(\boldsymbol{F} \boldsymbol{A}^{T}\right)=\operatorname{vec}\left(\left(\boldsymbol{A} \boldsymbol{F}^{T}\right)^{T}\right)=\operatorname{vec}\left(\left[\boldsymbol{A} \tilde{\boldsymbol{f}}_{1}\left|\boldsymbol{A} \tilde{\boldsymbol{f}}_{2}\right| \cdots \mid \boldsymbol{A} \tilde{\boldsymbol{f}}_{m_{2}}\right]^{T}\right)  \tag{4}\\
& =\operatorname{vec}\left(\left[\boldsymbol{y}_{1}|\cdots| \boldsymbol{y}_{m_{1}}\right]^{T}\right)=\left[y_{11}, \cdots, y_{m_{2} 1}, y_{12}, \cdots, y_{m_{2} 2}, \cdots, y_{1 m_{1}}, \cdots, y_{m_{2} m_{1}}\right]^{T} \tag{5}
\end{align*}
$$

where $\boldsymbol{y}_{i}=\left[y_{i 1}, y_{i 2}, \cdots, y_{i m_{1}}\right]^{T} \stackrel{\text { def }}{=} \boldsymbol{A} \tilde{\boldsymbol{f}}_{i}$, for all $i \in\left[m_{2}\right]$. Combining (2) together with (5) and recalling that $\boldsymbol{A} \otimes \boldsymbol{B}=\boldsymbol{C}$ and $\operatorname{vec}(\boldsymbol{X})=\boldsymbol{x}$ result in

$$
\begin{equation*}
\boldsymbol{b}=\left[y_{11}, \cdots, y_{m_{2} 1}, y_{12}, \cdots, y_{m_{2} 2}, \cdots, y_{1 m_{1}}, \cdots, y_{m_{2} m_{1}}\right]^{T} \tag{6}
\end{equation*}
$$

Hence, all $\boldsymbol{y}_{i}$ values for each $i \in\left[m_{2}\right]$ can be deduced from the measurement vector $\boldsymbol{b}$ using (6); specifically, $y_{i j}=b_{(j-1) m_{2}+i}$ for all $i \in\left[m_{2}\right]$ and $j \in\left[m_{1}\right]$. Then, one can arrange the following equations:

$$
\begin{equation*}
A \tilde{\boldsymbol{f}}_{i}=\boldsymbol{y}_{i}, \quad \forall i \in\left[m_{2}\right] \tag{7}
\end{equation*}
$$

Note that all $\tilde{f}_{i}$ 's then can be determined by utilizing the decoding algorithm corresponding to the matrix $A$ which is algorithm $\mathscr{A}$, i.e., $\tilde{f}_{i}=\mathscr{A}\left(\boldsymbol{y}_{i}\right)$, provided that the entries of $\tilde{f}_{i}$ 's are not larger than $d$. This requirement is necessary because matrix $\mathbf{A}$ is an $\left(n_{1}, d\right)$-detecting matrix. Hence the correct recovery of $\tilde{f}_{i}$ is only guaranteed if all the entries of $\tilde{f}_{i}$ do not exceed $d$. As $\tilde{f}_{i}$ 's are the columns of $\boldsymbol{F}^{T}$, it is sufficient to show that the entries of $\boldsymbol{F}$ do not exceed $d$. To establish this, note that all entries of $\boldsymbol{B} \boldsymbol{x}_{j}$ for all $j \in\left[n_{1}\right]$ are not greater than $d$. This follows because each entry is the sum of at most $n_{2}$ binary variables (recall that $\boldsymbol{B}$ is a binary matrix with $n_{2}$ columns), at most $d$ of which are 1 . The latter holds since $\boldsymbol{x}$ is $\left(n_{1}, d\right)$-balanced, leading to the constraint that the total number of 1's in $\boldsymbol{x}_{j}$ does not exceed $d$, according to Definition 3. Therefore, the entries of $\mathbf{F}$ satisfy the same condition due to (3). This confirms that none of the entries in $\tilde{f}_{i}$ exceed $d$. As a result, one can reconstruct the matrix $\boldsymbol{F}$ by invoking the algorithm $\mathscr{A}$ for $m_{2}$ times and stacking the results as the columns of $\boldsymbol{F}^{T}$. Next, let $f_{j}$ for all $j \in\left[n_{1}\right]$ denote the columns of $\boldsymbol{F}$, i.e., $\boldsymbol{F}_{m_{2} \times n_{1}}=\left[f_{1}\left|f_{2}\right| \cdots \mid f_{n_{1}}\right]$. This together with (3) leads to

$$
\begin{equation*}
\boldsymbol{B} \boldsymbol{x}_{j}=f_{j}, \quad \forall j \in\left[n_{1}\right] \tag{8}
\end{equation*}
$$

Similarly, each of the individual equations can be solved by invoking the decoding algorithm corresponding to the test matrix $B$, namely, $\mathscr{B}$. In other words, we have $x_{j}=\mathscr{B}\left(f_{j}\right)$ for all $j \in\left[n_{1}\right]$. Therefore, the incident vector $x$ can be recovered by invoking the algorithms $\mathscr{A}$ and $\mathscr{B}$ for $m_{2}$ and $n_{1}$ times, respectively. This results in $m_{2} O_{1}+n_{1} O_{2}$ combined computation complexity to decode $x$ from $C x$, which completes the proof.

The result of Theorem 4 implies that when dealing with a binary vector $x$ of length $n_{1} n_{2}$ that is $\left(n_{1}, d\right)$-balanced, the Kronecker product of an $\left(n_{1}, d\right)$-detection matrix with an $\left(n_{2}, d\right)$-search matrix yields a test matrix that can uniquely recover $\boldsymbol{x}$. In order to gain an insight into the implication of the result, one can regard the problem as a coin weighing problem with $n_{1} n_{2}$ coins. The $\left(n_{1}, d\right)-$ balanced assumption implies that one can partition all coins into $n_{1}$ subgroups consisting of $n_{2}$
coins, such that the total number of counterfeit coins (corresponding to 1's in the underlying incident vector $\boldsymbol{x}$ ) inside each subgroup does not exceed $d$. Suppose that the matrix $\boldsymbol{B} \in\{0,1\}^{m_{2} \times n_{2}}$ is an $\left(n_{2}, d\right)$-search matrix which solves this sub-problem. A naive approach then would be to utilize $\boldsymbol{B}$ for each subgroup independently in order to identify the location of all counterfeit coins. The overall test matrix then can be characterized as $\boldsymbol{I}_{n_{1}} \otimes \boldsymbol{B}$, which results in a total number of $m_{2} m_{1}$ tests. In comparison, the test matrix constructed in Theorem 4, which is the concatenation of $\boldsymbol{B}$ with an $\left(n_{1}, d\right)$-detecting matrix, requires only $m_{2} m_{1}$ tests where $m_{1}=o\left(n_{1}\right)$ is achievable. Therefore, the result of Theorem 4 introduces an approach for addressing the QGT problem through a divide-and-conquer strategy that significantly reduces the number of tests. However, this is true only when the underlying binary vector to be reconstructed adheres to a specific condition termed as a balanced vector, as defined in Definition 3. We will show in the next section that this holds with probability approaching 1 as the length of $x$ grows large, if one carefully chooses the size of the underlying subgroups under the combinatorial model for the QGT problem.

## 3. Connection to Urn Models

In this section we first review some results on the so-called balls into bins problem from the literature. We then demonstrate that a randomly selected vector $x$ from the set $\mathcal{B}_{n, k}$ is a balanced vector with certain parameters by making a connection between the underlying probabilistic model for the incident vector $x$ to the balls into bins problem. The description of this problem is as follows: There are $n$ bins and $m$ balls. Each ball is thrown into one of these $n$ bins. In the simplest scenario, the bins are chosen uniformly at random, i.e., each with probability $\frac{1}{n}$, as illustrated in Figure 1. There are various questions that arise concerning the statistics of the ball distributions including the statistics of the number of empty bins, the maximum number of balls inside a single bins, etc. Several variants of this problem are extensively studied in the literature, collectively falling under the broad category referred to as urn models (Sprott, 1978; Park, 1980). Let $M$ denote the maximum number of balls inside the individual bins in the balls into bins problem. We will use the following result on the probability of $M$ exceeding a certain threshold proposed by Raab and Steger (1998).

Theorem 5 ((Raab and Steger, 1998), Theorem 1) Let $M$ be the random variable that counts the maximum number of balls in any bin, if we throw $m$ balls independently and uniformly at random into $n$ bins. Then, if $m=c n \ln n$ for some constant $c$, we have

$$
\begin{equation*}
\operatorname{Pr}\left[M>\left(d_{c}-1+\alpha\right) \ln n\right]=o(1), \tag{9}
\end{equation*}
$$

for any $\alpha>1$. Here $d_{c}$ is a solution to

$$
\begin{equation*}
1+x(\ln c-\ln x+1)-c=0 \tag{10}
\end{equation*}
$$

that is larger than $c$.
It's worth noting that within the parameter range considered in Theorem 5, on average, each bin holds around $c \ln n$ balls. In particular, the pigeonhole principle implies that the maximum is always greater than $c \ln n$. The importance of the result in Theorem 5 is that it ensures the maximum number of balls in each bin deviates from its average by a constant that does not depend on $n$.

Next, we explore the connection between the ( $n, k$ )-QGT problem with the balls into bins problem and demonstrate how the result of Theorem 5 can be applied within the context of the QGT


Figure 1: Illustration of balls into bins problem with $m$ balls and $n$ bins.


Figure 2: Sampling without replacement demonstrated as an LDLC model. $n$ items are distributed into $l$ bins, with each selection of a red-colored ball representing a placement into a bin according to the LDLC model.
problem. Recall that $x$ is drawn from $\mathcal{B}_{n, k}$ uniformly at random, i.e., the subset of defective items is uniformly distributed over all $k$-subsets of $[n]$. One possible approach to generate a $k$-subset with incident vector $x$ according to the uniform distribution from $\mathscr{B}_{n, k}$ is to repeat the following for $k$ steps: At each step, one chooses an item among those items in [ $n$ ] that have not been chosen in previous steps uniformly at random and labels it as a defective item. This is often referred to as sampling without replacement. Let the set of all items be partitioned into $l$ distinct subsets, as depicted in Figure 2. Each of these partitions can be viewed as a bin, and the act of selecting an item from a particular partition can be regarded as placing a ball into that bin. Therefore, sampling $k$ items without replacement from the set of $n$ items can can be comprehended within the framework of urn models. One may note that the procedure of placing balls in the mentioned urn model does not align with the procedure in the balls into bins problem, where the result of Theorem 5 remains applicable. To see that, let $X_{j}^{i}$ for $i \in[k]$ and $j \in[l]$ denote random variables representing the number of items that have not been selected yet from bin $j$ (corresponding to uncolored items in Figure 2) before ball $i$ is placed. The main difference of the above urn model with the balls into bins problem is that ball $i$ is placed into bin $j$ with a probability that is proportional to $X_{j}^{i}$, whereas all bins are chosen uniformly at random in all steps in the balls into bins problem as depicted in Figure 1. Thus, the urn model depicted in Figure 2 differs from the balls into bins problem in two major aspects. First, the capacity of each urn is limited, meaning there is a constraint on the number of balls that can be placed into each bin. In particular, for the scenario considered above where $n$ items are divided into $l$ partitions, the number of defective items in each bin cannot exceed $\frac{n}{l}$. Second, when regarded as an urn model, sampling without replacement intrinsically exhibits a form of negative reinforcement, where a new ball is less likely to be placed in a bin that already contains more balls, compared to other bins. We refer to this urn model as linearly de-preferential with limited capacity (LDLC) urn model. In the next lemma, we show that the probability of the maximum number of balls inside bins in the LDLC urn model is always lower bounded by that of the balls into bins problem.

Lemma 6 Let $M_{n, m}^{\mathscr{R}}$ denote the maximum number of balls in the balls into bins problem with $n$ bins and $m$ balls. Let also $M_{n, m}^{\mathscr{L}}$ denote the maximum number of balls in the LDLC urn model with $n$ bins and $m$ balls. Then, for any positive integer $t$, we have

$$
\begin{equation*}
\operatorname{Pr}\left[M_{n, m}^{\mathscr{B}} \leqslant t\right] \leqslant \operatorname{Pr}\left[M_{n, m}^{\mathscr{L}} \leqslant t\right] . \tag{11}
\end{equation*}
$$

The proof of Lemma 6 is provided in Appendix A.
The result of Lemma 6 can be intuitively justified by comparing the bin selection process in each step. Note that in the balls into bins problem, a ball is placed into a randomly chosen bin, regardless
of the arrangement of previously placed balls. In contrast, in the LDLC model, the probability of selecting a bin with a higher number of balls already in it is reduced, creating a disincentive to choose bins with more balls. This negative feedback mechanism makes it less likely for the bin containing the maximum number of balls to be chosen for the next ball placement. In other words, the random variable $M_{n, m}^{\mathscr{L}}$ is less likely to increase compared to $M_{n, m}^{\mathscr{B}}$ at each step.

Building upon the results of Theorem 5 and Lemma 6, it is shown in the following corollary that the incident vector $\boldsymbol{x}$ drawn uniformly from $\mathscr{B}_{n, k}$ is an $\left(l, c^{\prime} \log k\right)$-balanced vector with high probability for any constant $c^{\prime}>2$ and a carefully chosen parameter $l$.

Corollary 7 Let $x$ be a vector that is drawn from $\mathscr{B}_{n, k}$ according to a uniform distribution. Then, $x$ is $\left(\frac{k}{\ln k},(e+\gamma) \ln k\right)$-balanced with probability $1-o(1)$, for any constant $\gamma>0$.

Proof: Suppose $x$ is partitioned into $l$ distinct subsets as shown in Figure 2 where $l=\frac{k}{\ln k}$. Note that this resembles an LDLC urn model with $m=k$ balls and $n=\frac{k}{\ln k}$ bins. Then,

$$
\begin{align*}
\operatorname{Pr}\left[M_{\substack{k \\
\ln , k}}^{\mathscr{L}} \leqslant(e+\gamma) \ln k\right] & \leqslant \operatorname{Pr}\left[M_{\substack{k \\
k \\
\mathscr{R}}} \leqslant(e+\gamma) \ln k\right]  \tag{12}\\
& =1-o(1), \tag{13}
\end{align*}
$$

where (12) follows by the result of Lemma 6 and (13) is established by recognizing that the number of balls and bins satisfy the condition in Theorem 5. In other words, we have $n \ln n=\frac{k}{\ln k} \log \left(\frac{k}{\ln k}\right)=$ $k(1+o(1))=m$, corresponding to the choice of $c=1$ in Theorem 5 . Consequently, $d_{c}$ represents the larger solution to the equation (10) with $c=1$, which results in $d_{c}=e$. Lastly, note that $\alpha-1$ for any $\alpha>1$ can be replaced with $\gamma$ for any $\gamma>0$.

The implication of the result established in Corollary 7 bears considerable significance for the $(n, k)$-QGT problem. In essence, as $n$ grows large, this implies that the problem can be decomposed into $\frac{k}{\ln k}$ instances of $\left(\frac{n \ln k}{k}, e \ln k\right)$-QGT problems. This decomposition enables us to tackle the problem by using the concatenated approach detailed in Section 2. Roughly speaking, this suggests that a typical incident vector $\boldsymbol{x}$, selected uniformly at random from $\mathscr{B}_{n, k}$, can be divided into roughly $\frac{k}{\ln k}$ distinct subsets, with the total number of defective items across all these subsets not exceeding $e \ln k$. Recall that the pigeonhole principle implies that the maximum is at least $\ln k$. This shows the bound derived through observing the connection to the balls into bins problem, as characterized in Lemma 6, is essentially optimal up to a constant factor less than $e$. In the next section, we provide specific details about the constructions employed for certain detecting and search matrices that are utilized in the concatenation method proposed in Section 2 to tackle the $(n, k)$-QGT problem.

## 4. Proposed Construction

In this section, we briefly overview specific constructions for $(n, d)$-detecting matrices and $(n, k)$ search matrices employed in our construction. In particular, we show that the decoding algorithm of the ( $n, d$ )-detecting matrices constructed by Bshouty (2009) can be modified to detect decoding failure when at least one of the entries of the underlying integer vector exceeds $d$. Then, building upon the results provided in Section 2 and Section 3, we propose our concatenated construction.

### 4.1. An Optimal Construction for Detecting Matrices

We provide a concise summary of the construction method proposed in Bshouty (2009) for detecting matrices. This construction is designed to generate optimal $\left(d_{1}, \cdots, d_{n}\right)$-detecting matrices. In
particular, our specific focus is on detecting matrices where $d_{i}=d$ for all $i \in[n]$. Therefore, we will explore the construction tailored to this specific case.

The primary idea behind constructing the $(n, d)$-detection matrix in Bshouty (2009) involves carefully crafting a certain collection of binary functions defined on the domain $\{-1,+1\}^{2^{v}}$ for a given integer $v$. The specific characteristics of these functions are detailed in Appendix B. These functions are labeled by $\boldsymbol{a} \in\{0,1\}^{v}$ and $i \in\left[l_{a}\right]$, where $l_{a}$ is a non-negative integer that satisfies the condition $d^{l_{a}} \leqslant 2^{|a|-1}<d^{l_{a}+1}$. We represent these functions as $g_{a, l_{a}}:\{-1,+1\}^{v} \rightarrow\{0,1\}$, and the entire set containing such functions is denoted as $\mathcal{G}_{v, d}$. The $(n, d)$-detecting matrix takes the form of a matrix with rows labeled by binary vectors of $\pm 1$ values and columns labeled by functions from the set $\mathcal{G}_{v, d}$. Then, the $(i, j)$ entry of $\boldsymbol{M}$ is equal to the evaluation of the function corresponding to column $j$, evaluated at the evaluation point corresponding to row $i$. In simpler terms, each column of this matrix stores the results of applying a function from $\mathcal{G}_{v, d}$ to all possible inputs from $\{-1,+1\}^{\nu}$. The number of rows in matrix $M$ is at most

$$
\begin{equation*}
\frac{2 n}{\log n} \log d(1+o(1)), \tag{14}
\end{equation*}
$$

which asymptotically matches the information theoretic lower bound for the minimum number of non-adaptive tests. Let $B \stackrel{\text { def }}{=}\left\{\chi_{\boldsymbol{a}}(\boldsymbol{x}) \stackrel{\text { def }}{=} \prod_{a_{i}=1} x_{i} \mid \boldsymbol{a} \in\{0,1\}^{\nu}\right\}$. Since $B$ is an orthogonal set of functions, any function $f(x)$ with domain $\{-1,1\}^{v}$, including the functions in $\mathcal{G}_{v, d}$, can be uniquely represented as

$$
\begin{equation*}
f(x)=\sum_{a \in\{0,1\}^{v}} \hat{f}_{a} \chi_{a}(x), \tag{15}
\end{equation*}
$$

where $\hat{f}_{a}$ is called the Fourier coefficient of $\chi_{\boldsymbol{a}}(x)$ in $f(x)$ and is equal to

$$
\begin{equation*}
\hat{f}_{a}=\frac{1}{2^{v}} \sum_{x \in\{-1,+1\}^{v}} f(x) \chi_{a}(x) . \tag{16}
\end{equation*}
$$

The functions in the set $\mathcal{G}_{v, d}$ have an important property that is utilized in establishing the matrix $M$ as a $d$-detecting matrix. This property also enables the development of an efficient decoder for determining $x$ from $M x$, provided that all entries of $x$ are non-negative integers less than $d$. This property can be described as in the following. Consider a linear combination of such functions as $h_{\lambda}(x)=\sum_{a \in A} \sum_{i \in\left\{0,1, \cdots, l_{a}\right\}} \lambda_{a, i} g_{a, i}(x)$ for some $A \subset\{0,1\}^{v}$. Then, if $\mathbf{b} \in A$ is a maximal element in $A$, the Fourier coefficient of $h(x)$ in $\chi_{\mathbf{b}}$ is equal to $\sum_{j \in\left\{0,1, \cdots, l_{a}\right\}} \lambda_{\mathbf{b}, j} d^{j}$. This property is exploited during the decoding procedure as discussed below.

Suppose we have a vector $\lambda \in\{0,1, \cdots, d-1\}^{n}$. This vector can be regarded as a function within the basis $\mathcal{G}_{v, d}$. Therefore, recovering $\lambda$ is equivalent to retrieving the function $h_{\lambda}(x)=$ $\sum_{a \in\{0,1\}^{v}} \sum_{i \in\left\{0,1, \cdots, l_{a}\right\}} \lambda_{a, i} g_{a, i}(x)$. Note that the measurement vector $\boldsymbol{M} \boldsymbol{\lambda}$ provides all the evaluations of $h_{\lambda}(x)$ over $\{-1,1\}^{\nu}$. Therefore one can determine all the Fourier coefficients of $h_{\lambda}(x)$ at the decoder and search for a maximal $a \in\{0,1\}^{v}$ whose corresponding Fourier coefficient is non-zero. All the entries of $\lambda$ corresponding to $a$ can be recovered from the expansion of $\sum_{j \in 0,1, \ldots, l_{a}} \lambda_{\mathrm{a}, j} d^{j}$ in base $d$. This process can be repeated recursively, replacing $h_{\lambda}(x)$ with $h_{\lambda}(x)-\sum_{j \in 0,1, \cdots, l_{a}} \lambda_{\mathrm{a}, j} g_{a, j}$ until all the entries of $\lambda$ are determined. This demonstrates that the recovery of $\lambda$ can be accomplished in $O\left(n^{2}\right)$ time. Next, we present our observation that extends the decoding algorithm to handle cases where the vector $\boldsymbol{\lambda}$ violates the assumption that all of its non-negative integer entries are less than $d$. In particular, we show that the decoding procedure discussed earlier, with a slight

```
Algorithm 1 Modified decoding algorithm for detecting matrix \(M\)
Input: \(M \lambda\), i.e., \(h_{\lambda}(x)\) for all \(x \in\{-1,1\}^{v}\).
Output: \(\lambda\) or decoding failure.
while \(h_{\lambda}(x) \neq 0\) do
    Compute Fourier coefficients \(\hat{h}_{a}\) of \(\chi_{a}(x)\) in \(h_{\lambda}(x)\), defined in (16), for all \(\boldsymbol{a} \in\{0,1\}^{\nu}\).
    Find a maximal \(\tilde{\boldsymbol{a}} \in\{0,1\}^{v}\) such that \(\hat{h}_{\tilde{a}} \neq 0\).
    Expand \(\hat{h}_{\tilde{a}}\) in base \(d\). The coefficients are \(\lambda_{\boldsymbol{a}, j}\) for \(j \in\left[l_{a}\right]\), where \(l_{\boldsymbol{a}}\) is a non-negative integer that satisfies
    the condition \(d^{l_{a}} \leqslant 2^{|a|-1}<d^{l_{a}+1}\).
    if The expansion surpasses the maximum integer representable with \(l_{a_{i}}+1\) digits in base \(d\) : then
        Break; declare decoding failure.
    end
    Set \(h_{\lambda}(\boldsymbol{x})=h_{\lambda}(\boldsymbol{x})-\sum_{j \in\left\{0,1, \cdots, l_{a}\right\}} \lambda_{\mathbf{a}_{i} j} g_{a_{i} j}(\boldsymbol{x})\).
    Compute the Fourier coefficient \(\hat{h}_{\tilde{a}}\) of \(\chi_{\tilde{a}}(x)\) in \(h_{\lambda}(x)\).
    if \(\hat{h}_{\tilde{a}} \neq 0\) then
        Break; declare decoding failure.
    end
end
Return \(\lambda\).
```

adjustment, can determine whether any entry of $\lambda$ exceeds $d$ and declare it as a decoding failure. This modified algorithm is detailed in Algorithm 1.

Lemma 8 Consider $\boldsymbol{M} \boldsymbol{\lambda}=\boldsymbol{v}$, where $\boldsymbol{\lambda}$ is a non-negative integer vector. By modifying the algorithm described in Bshouty (2009), as provided in Algorithm 1, it can effectively detect a decoding failure when there exists at least one entry in $\boldsymbol{\lambda}$ that is greater than or equal to $d$.

Proof: Let $h_{i}(x)$ represent the function remaining after iteration $i$ during the recursive decoding procedure described above. Suppose that at iteration $i$, we select the maximal binary vector denoted as $\boldsymbol{a}_{i}$. In this case, we can express $h_{i}(x)$ as: $h_{i}(x)=h_{i-1}(x)-\sum_{j \in\left\{0,1, \ldots, l_{a}\right\}} \lambda_{\mathbf{a}_{i} j} g_{a_{i} j}$. This means that all the coefficients corresponding to the functions indexed by $\boldsymbol{a}_{i}$ are effectively subtracted from $h_{i}(x)$, resulting in the removal of all such functions, i.e., $g_{\boldsymbol{a}_{i}, j}$ for all $j \in\left[l_{\boldsymbol{a}_{i}}\right]$. Consequently, the Fourier coefficient of $\chi_{a_{i}}$ in $h(x)$ becomes zero.

In the event that any of the $\lambda_{\mathbf{a}_{i}, j}$ values is greater than or equal to $d$, two scenarios arise. Either the Fourier coefficient of $\chi_{a_{i}}$ in $h_{i}(x)$ surpasses the maximum integer representable with $l_{a_{i}}+1$ digits in base $d$, or at least one of the $\lambda_{a_{i} j}$ values is incorrectly recovered during the decoding process at iteration $i$. In the former case, the decoder effortlessly detects that at least one entry is not an integer less than $d$. In the latter case, where a decoding failure is not immediately apparent from the Fourier coefficient of $\chi_{a_{i}}$ in $h_{i}(x)$, the decoder must additionally verify whether the Fourier coefficient of $\chi_{a_{i}}$ in the function derived from the pruning of $h_{i}(x)$ is zero or not. If this coefficient is indeed zero, it shows that the entries are determined correctly, implying that all such entries are integers less than $d$. Conversely, if the coefficient is non-zero, the decoder realizes that at least one entry violates the initial assumption.

The result of Lemma 8 implies that by incorporating this modification into the decoder for the ( $n, d$ )-detecting matrix constructed by Bshouty (2009), the decoder gains the ability to determine whether the input vector adheres to the assumption that all of its non-negative integer entries are less than $d$ or not. It can also correctly recover these entries if they do satisfy the assumption. Later in this section, we will demonstrate how this modification enables our decoding algorithm for the
$(n, k)$-QGT problem to identify cases where it cannot recover the underlying incident vector. This capability allows the algorithm to declare a decoding failure, affirming that if it does produce a vector, that vector must indeed be the unique incident vector.

## 4.2. $(n, k)$-Search matrix with $k=O(\log n)$

In this subsection, we focus on addressing the $(n, k)$-QGT problem with a specific condition: when $k$ is $O(\log n)$. This scenario represents a much sparser setting compared to the case where $k=n^{\alpha}$. Note that the number of defective items still grows large, although with a speed that is logarithmic in $n$, which means that algorithms designed for the sparse regime where $k$ is constant, may not be optimal. In this regime, for $(n, k)$-search matrices, we utilize the construction of BCH codes from coding theory literature. In particular, the parity-check matrix of BCH codes with length $n=2^{m-1}$ for some positive integer $m$, and parameter $k$, referred to as an $(n, k) \mathrm{BCH}$ code, is an $(n, k)$-search matrix. These codes has been also utilized as a building block to construct test matrices for the QGT problem in Karimi et al. (2019b). For a concise summary on BCH codes, please refer to Appendix C. It suffices to mention that the parity-check matrix of this code has at most $k \log (n+1)$ measurements/rows. Moreover, a binary vector $x$ of Hamming weight at most $k$ can be recovered from $H x$ in $O(k n)$ time complexity by utilizing the well known algorithms for decoding BCH codes, e.g., the Berlekamp-Massey algorithm (Berlekamp, 2015; Massey, 1969) the Euclidean algorithm (Sugiyama et al., 1975) the Berlekamp-Welch algorithm (Welch and Berlekamp, 1986).

### 4.3. Our Construction

Now that all necessary building blocks and theoretical results are established, we are ready to propose our construction for $(n, k)$-search matrices. The main idea behind our construction can be described as in the following: The set of all items are first partitioned into $\frac{k}{\ln k}$ groups. Then, the result of Corollary 7 implies that there are at most $e \ln k$ in each partition with probability $1-o(1)$. The problem then can be solved by using the concatenation scheme outlined in Section 2, where the $(n, d)$-detecting matrix reviewed in Section 4.1 and the parity-check of a BCH code with suitable parameters, as discussed in Section 4.2 are used as building blocks.

Let $\boldsymbol{D}$ denote a $\left(\frac{k}{\ln k},(e+\gamma) \ln k\right)$-detecting matrix constructed according to the scheme overviewed in 4.1 for some $\gamma>0$. Let $\tilde{\boldsymbol{H}}$ denote the parity-check matrix of a $(\tilde{n},(e+\gamma) \ln k) \mathrm{BCH}$ code, where $\tilde{n}$ is the smallest power of 2 such that $\frac{n \ln k}{k} \leqslant \tilde{n}$. Let also $\boldsymbol{H}$ be a submatrix of $\tilde{H}$ consisting of any $\frac{k}{\ln k}$ columns of $\tilde{\boldsymbol{H}}$. One can see that $\boldsymbol{H}$ is a $\left(\frac{k}{\ln k},(e+\gamma) \ln k\right)$-search matrix since $\tilde{\boldsymbol{H}}$ is a $(\tilde{n},(e+\gamma) \ln k)$-search matrix. Then, the following matrix is our proposed test matrix for the $(n, k)$-QGT problem in the sublinear regime where $k=n^{\alpha}$ for some $0<\alpha<1$ :

$$
\begin{equation*}
T=D \otimes H \tag{17}
\end{equation*}
$$

The following theorem establishes the probability of successfully decoding the underlying incident vector when utilizing this test matrix.

Theorem 9 Let $\boldsymbol{x}$ be a binary vector sampled from the uniform distribution over $\mathcal{B}_{n, k}$. Then, $\boldsymbol{x}$ can be recovered from $T \boldsymbol{x}$, where $T$ is characterized in (17), with probability $1-o(1)$ and the decoding complexity of

$$
O\left(k^{2}+n \log ^{2} k\right)
$$

using no more than

$$
(2 e+\gamma) \frac{k}{\log k} \log \left(\frac{n}{k}\right) \log \log k,
$$

tests for any $\gamma>0$. Moreover, in cases where the decoder cannot obtain a unique reconstruction of $\boldsymbol{x}$, it can reliably detect this condition and report a decoding failure.

Proof: The result of Corollary 7 implies that the vector $x$ is $\left(\frac{k}{\ln k},(e+\gamma) \ln k\right)$-balanced with probability $1-o(1)$, for any constant $\gamma>0$. In other words, for any positive constant $\gamma$, it is highly likely that the number of defective items allocated to each of the $\frac{k}{\ln k}$ partitions will not exceed $(e+\gamma) \ln k$. Therefore, the result of Theorem 4 can be applied to this case with $n_{1}=\frac{k}{\ln k}, m_{1}=$ $\frac{2 n_{1}}{\log n_{1}} \log ((e+\gamma) \ln k)$ according to (14), and $n_{2}=\frac{n \ln k}{k}, m_{2}=(e+\gamma) \ln k \log (\tilde{n}+1)$, as described in 4.2, and $d=(e+\gamma) \ln k$. Recall that $\tilde{n}$ is the smallest power of 2 such that $\frac{n \ln k}{k} \leqslant \tilde{n}$. This implies $\tilde{n} \leqslant \frac{2 n \ln k}{k}$, and thus, $m_{2} \leqslant(e+\gamma) \ln k \log \left(\frac{2 n \ln k}{k}+1\right)$. Specifically, the result of Theorem 4 implies that the underlying incident vector can be recovered with complexity $m_{2} O_{\mathscr{A}}+n_{1} O_{\mathscr{B}}$, where $O_{\mathscr{A}}$ and $O_{\mathscr{B}}$ denote the complexities of the underlying decoders associated with $\boldsymbol{D}$ and $\boldsymbol{H}$, respectively. Recall that the complexity of the decoding algorithm described in 4.1 is $O_{\mathscr{A}}=O\left(n_{1}^{2}\right)=O\left(\frac{k^{2}}{\ln ^{2} k}\right)$ and the decoding complexity of BCH codes, described in 4.2 , is $O_{\mathscr{B}}=O\left(m_{2} n_{2}\right)=O\left(\frac{n \ln ^{3} k}{k}\right)$. Therefore, the overall combined computational complexity is $O\left(k^{2}+n \log ^{2} k\right)$. Note that decoding failure only occurs when $x$ is not $\left(\frac{k}{\ln k},(e+\gamma) \ln k\right)$-balanced, in which case the decoder detects it during the decoding procedure (Lemma 8). Also, the number of rows in $\boldsymbol{D}$ is upper bounded by $\frac{2 n_{1} \log d}{\log n_{1}}(1+o(1))=\frac{2 k}{\log k \ln k} \log \log k(1+o(1))$, according to (14), and the number of rows in $\boldsymbol{H}$ is $m_{2} \log (\tilde{n}+1) \leqslant(e+\gamma) \ln k \log \left(\frac{2 n \ln k}{k}\right)$. Overall, the number of rows in $T$ is upper bounded by $(2 e+\gamma) \frac{k}{\log k} \log \left(\frac{n}{k}\right) \log \log k$, for an arbitrarily small $\gamma>0$.

## 5. Conclusion

In this work, we proposed a method to recover a $k$-sparse binary vector from additive tests in the sublinear regime, where $n=k^{\alpha}$ with $0<\alpha<1$. Our construction uses $2 e_{0} \log \log (k)$ tests with decoding complexity of $O\left(n^{2 \alpha}+n \log ^{2} n\right)$ while the asymptotically optimal scheme (HahnKlimroth and Müller, 2022; Hahn-Klimroth et al., 2023) requires $\left(\frac{1+\sqrt{\alpha}}{1-\sqrt{\alpha}}\right) m_{0}$ tests with the decoding complexity of $O\left(n^{3}\right)$, where $m_{0}$ denotes the information-theoretic lower bound. This demonstrates that our construction offers improvement in decoding complexity at the cost of a gap of $\log \log k$ in terms of the number of measurements. In particular, our construction achieves an almost linear decoding complexity for $\alpha<\frac{1}{2}$ and sub-quadratic decoding complexity for $\alpha \geqslant \frac{1}{2}$. Moreover, despite a double logarithmic gap, our scheme numerically outperforms the asymptotically optimal construction when $\alpha$ is in the higher range of the unit interval. Specifically, the crossover value for $k$ is $k^{*}=\exp \left(\exp \left(\frac{1+\sqrt{\alpha}}{e(1-\sqrt{\alpha})}\right)\right)$, growing double-exponentially in $\frac{1}{1-\sqrt{\alpha}}$. Hence, for moderate values of $\alpha$, the crossover becomes significantly large, highlighting the practicality of our scheme in the sublinear regime for higher $\alpha$ values. For example, with $\alpha=0.7$, we find $k^{*} \approx 10^{27}$.

Compared to other works in the literature that provide nearly linear decoding complexity (Gebhard et al., 2022; Karimi et al., 2019b,a), which fall short of the information-theoretic bound by a factor of $\log k$, our construction reduces this gap by a factor of $\frac{\log k}{\log \log k}$ while offering the same decoding complexity for $0<\alpha<\frac{1}{2}$.

Moreover, our decoding algorithm reliably identifies decoding failures when unique recovery is not feasible. As a result, in cases where decoding is successful, the inferred feature vector is error-free, a feature not guaranteed in the aforementioned schemes.

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## Appendix A. Proof of Lemma 6

Proof: The main idea involves creating a combined process that mixes random allocation based on the balls into bins problem with an extra step called the leakage phase, similar to the assignment method in the LDLC urn model. Showing that this extra step does not increase the maximum number of balls leads to the desired outcome. Consider the process shown in Figure 3, which includes $n$ bins and $m$ balls. Each bin has a capacity of $c$, meaning it can hold up to $c$ balls. Initially, each ball, numbered $i \in[m]$, is randomly placed in one of the $n$ bins, namely bin $l$, according to the uniform distribution, following the same method as the original urn model in the balls into bins problem. Let $X_{j}^{i}$ represent the number of empty spaces in bin $j$ before ball $i$ is placed in any bin $j \in[n]$. Upon ball $i$ landing in bin $j$ during the first step, the subsequent phase involves the potential leakage of the same ball to another bin $j \neq l$, with probability.

$$
\begin{equation*}
p_{l j}^{i} \stackrel{\text { def }}{=} \frac{\left(X_{j}^{i}-X_{l}^{i}\right)_{+}}{\sum_{k=1}^{n} X_{k}^{i}}, \quad \forall j \in[n] \backslash\{l\} \tag{18}
\end{equation*}
$$

where $(x)_{+} \stackrel{\text { def }}{=} \max (0, x)$. Also the ball remains in bin $j$ with probability

$$
\begin{equation*}
p_{l l}^{i} \stackrel{\text { def }}{=} 1-\sum_{j=1, j \neq l}^{n} p_{l j}^{i} . \tag{19}
\end{equation*}
$$

To demonstrate that equations (18)-(19) constitute a valid probability distribution, we need to establish that $\sum_{j=1, j \neq l}^{n} p_{l j}^{i} \leqslant 1$. To this end, consider the following:

$$
\sum_{j=1, j \neq l}^{n} p_{l j}^{i}=\sum_{j=1, j \neq l}^{n} \frac{\left(X_{j}^{i}-X_{l}^{i}\right)_{+}}{\sum_{k=1}^{n} X_{k}^{i}} \leqslant \frac{\sum_{j: X_{j}^{i}>X_{l}^{i}} X_{j}^{i}}{\sum_{k=1}^{n} X_{k}^{i}} \leqslant \frac{\sum_{j=1}^{n} X_{j}^{i}}{\sum_{k=1}^{n} X_{k}^{i}}=1,
$$

where the first and the second inequalities arise from the fact that $X_{l}^{i} \geqslant 0$ and $X_{j}^{i} \geqslant 0$, for all $j \in[n]$, respectively.

Next, we determine the probability that ball $i$ ultimately lands in bin $j$ after the second random transition described above. This probability is denoted as $q_{j}^{i}$. Let $A_{l}^{i}$ for all $l \in[n]$ represent the probability of ball $i$ landing in bin $l$ after the first step. It follows that $\operatorname{Pr}\left[A_{l}^{i}\right]=\frac{1}{n}$ since the initial random allocation in the first step is uniform. Similarly, let $B_{j}^{i}$ denote the event that ball $i$ lands in


Figure 3: Demonstration of the combined process equivalent to random allocation within the LDLC urn model is depicted here. In this instance, with a total of $n=4$ bins, two steps are illustrated concerning the placement of ball $i=7$. Each bin has a capacity denoted by $c$. Solid lines denote the outcomes of the random decision made at each step, while dashed lines represent all alternative potential decisions. Transition probabilities are represented for each possibility. Note that during the initial allocation step, all transition probabilities are identical which resembles the allocation procedure in the balls into bins problem. However, during the second allocation, probabilities are contingent upon the remaining space within the respective bins. Specifically, the transition from bin 2 to bin 4 is prohibited ( $p_{24}^{i}=0$ ) because bin 4 possesses only $c-3$ vacant spaces, whereas bin 2 has $c-2$ vacant spaces remaining.
bin $j$ after the second step. Then, we have

$$
\begin{align*}
q_{j}^{i} & \stackrel{\text { def }}{=} \operatorname{Pr}\left[B_{j}^{i}\right]=\sum_{l=1}^{n} \operatorname{Pr}\left[A_{l}^{i}\right] \operatorname{Pr}\left[B_{j}^{i} \mid A_{l}^{i}\right]=\frac{1}{n} \sum_{l=1}^{l=n} \operatorname{Pr}\left[B_{j}^{i} \mid A_{l}^{i}\right]  \tag{20}\\
& =\frac{1}{n} \sum_{l=1}^{l=n} p_{l j}^{i}=\frac{1}{n}\left(1-\sum_{l=1, l \neq j}^{n} p_{j l}^{i}\right)+\frac{1}{n} \sum_{l=1, l \neq j}^{l=n} p_{l j}^{i}  \tag{21}\\
& =\frac{1}{n}\left[1+\sum_{l=1}^{n} \frac{\left(X_{j}^{i}-X_{l}^{i}\right)_{+}-\left(X_{l}^{i}-X_{j}^{i}\right)_{+}}{\sum_{k=1}^{n} X_{k}^{i}}\right]  \tag{22}\\
& =\frac{1}{n}\left[1+\frac{\sum_{l: X_{j}^{i} \geqslant X_{l}^{i}}\left(X_{j}^{i}-X_{l}^{i}\right)-\sum_{l: X_{j}^{i}<X_{l}^{i}}\left(X_{l}^{i}-X_{j}^{i}\right)}{\sum_{k=1}^{n} X_{k}^{i}}\right]  \tag{23}\\
& =\frac{1}{n}\left(1+\frac{n X_{j}^{i}-\sum_{l=1}^{n} X_{l}^{i}}{\sum_{k=1}^{n} X_{k}^{i}}\right)=\frac{X_{j}^{i}}{\sum_{k=1}^{n} X_{k}^{i}} \tag{24}
\end{align*}
$$

where (20) arises from the law of total probability and noting that $\operatorname{Pr}\left[A_{l}^{i}\right]=\frac{1}{n}$ for all $l$, (21) is established by considering $\operatorname{Pr}\left[B_{j}^{i} \mid A_{l}^{i}\right]$ as the transition probability $p_{l j}^{i}$, then delineating the sum for each case, whether $l=j$ or $l \neq j$, (22) is derived from equations (18) and (19), taking into account the condition $\left(X_{j}^{i}-X_{l}^{i}\right)_{+}=\left(X_{l}^{i}-X_{j}^{i}\right)_{+}=0$ when $l=j$, (23) arises from the definition of $(x)_{+}$, and, (24) follows through simplification steps.

The result in (24) demonstrates that the probability of ball $i$ ultimately landing in bin $j$ is proportional to $X_{j}^{i}$, the number of remaining spaces in bin $j$, thereby validating the consistency of the discussed process with the LDLC model. Consequently, the maximum number of balls in each bin is $M_{n, m}^{\mathscr{L}}$. Note also that if we omit the second step for each iteration of placing ball $i$, the bins are chosen uniformly at random, resembling the balls into bins problem. Let $M_{n, m}^{\mathscr{B}}$ represent the maximum number of balls in the bins for this scenario. A crucial observation is that the maximum number of balls never increases during the leakage process (the second step), as the probability of ball $i$ transitioning from bin $l$ to another bin $j$ where $X_{j}^{i}<X_{l}^{i}$ is zero. This implies that ball $i$ might either remain in the same bin or leak to a bin with a lesser number of balls. Let $P^{\mathscr{L}}=\left\{\left(l^{i}, j^{i}\right): i \in[m]\right\}$ denote a sample outcome of the random assignment rule with leakage, and $Q^{\mathscr{B}}=\left\{l^{i}: i \in[m]\right\}$ denote the same path without leakage. The set $Q^{\mathscr{B}}$ is derived from $P^{\mathscr{L}}$ by removing the second elements in all pairs. Note that for all samples $Q^{\mathscr{B}}$ with $M_{n, m}^{\mathscr{B}} \leqslant t$ for some constant $t$ we also have $M_{n, m}^{\mathscr{L}} \leqslant t$ for the corresponding sample(s) $P^{\mathscr{L}}$ since the maximum number of balls never increases during the leakage step. This implies that the event $M_{n, m}^{\mathscr{B}} \leqslant t$ is a subset of the event $M_{n, m}^{\mathscr{L}} \leqslant t$, thus implying

$$
\begin{equation*}
\operatorname{Pr}\left[M_{n, m}^{\mathscr{B}} \leqslant t\right] \leqslant \operatorname{Pr}\left[M_{n, m}^{\mathscr{L}} \leqslant t\right] . \tag{25}
\end{equation*}
$$

## Appendix B. Summary of the Construction for Detecting Matrix in Bshouty (2009)

In this section, we first introduce a construction for an $\left(k_{1}, k_{2}, \cdots, k_{n}\right)$-detection matrix, i.e., a test matrix for the coin weighing problem with constraints that achieves the asymptotically optimal number of tests. This scheme is proposed by Bshouty (2009). The $i$ th component of a vector $\boldsymbol{x}$ is denoted by $x_{i}$. Let $<$ denotes the usual lattice partial ordering over on dimensional vectors. Specifically, for two $\boldsymbol{a}, \boldsymbol{b} \in\{0,1\}^{n}$, we say $\boldsymbol{b}<\boldsymbol{a}$ if and only if the support of $\boldsymbol{b}$ is a subset of $\boldsymbol{a}$. For $\boldsymbol{a}, \boldsymbol{b} \in\{0,1\}^{n}$ where $\boldsymbol{b}<\boldsymbol{a}$, we define

$$
\begin{equation*}
f_{a, b}(x)=\prod_{a_{i}=1} \frac{(-1)^{b_{i} \boldsymbol{x}_{i}+1}}{2} . \tag{26}
\end{equation*}
$$

Algorithm 2 provides the steps of the construction proposed for detection matrices in Bshouty (2009) in detail. The algorithm outputs a binary matrix $\boldsymbol{M}$ whose number of rows are asymptotically optimal as the number of coins, namely $l$, grows large.

```
Algorithm 2 Constructing an \(\left(k_{1}, \cdots, k_{l}\right)\) detecting matrix \(\boldsymbol{M}\).
Input: \(k_{1}, \cdots, k_{l},\left(k_{1} \leqslant k_{2} \cdots \leqslant k_{l}\right), l\).
Find the maximal integer \(v\) such that
\[
\begin{equation*}
(v-2) 2^{v-1} \leqslant \log \left(k_{l-2^{v}}^{2^{v}} \prod_{i=1}^{n-2^{v}} k_{i}\right) . \tag{27}
\end{equation*}
\]
```

Set $M_{2^{v} \times l}=0_{2^{v} \times l}, r=0$ and $s=0$.
for $\forall \boldsymbol{a} \in\{0,1\}^{\nu} \backslash\{0\}$ do
Find $l_{a}$ such that

$$
k_{r+1} k_{r+2} \cdots k_{r+l_{a}} \leqslant 2^{\|a\|_{0}-1}<k_{r+1} k_{r+2} \cdots k_{r+l_{a}} k_{r+l_{a}+1}
$$

Construct $G_{\boldsymbol{a}}=\left\{f_{a, \boldsymbol{b}(x)}: \boldsymbol{b}<\boldsymbol{a},\|\boldsymbol{b}\|_{0} \equiv 0 \bmod 2\right\}\left(\right.$ check $\left.\left\|G_{\boldsymbol{a}}\right\|_{0}=2^{\|\boldsymbol{a}\|_{0}-1}\right)$
Choose any collection of subsets

$$
G_{a, 0}, G_{a, 1}, \cdots, G_{a, l_{a}} \subset G_{a}
$$

such that

$$
\left\|G_{a, 0}\right\|_{0}=1, \quad\left\|G_{a, i}\right\|=k_{r+1} k_{r+2} \cdots k_{r+i} \quad \forall i=1, \cdots l_{a} .
$$

for $i \in\left[l_{a}\right]$ do
Set

$$
h_{a, i}(x) \stackrel{\text { def }}{=} \sum_{g(x) \in G_{a, i}(x)} g(x) .
$$

Set

$$
\boldsymbol{M}[:, r+i]=\left.h_{a, i}(\boldsymbol{x})\right|_{\{-1,+1\}^{v}} .
$$

end
Set

$$
h_{a, l_{a}}(x) \stackrel{\text { def }}{=} \sum_{g(x) \in G_{a, l a}(x)} g(x) .
$$

Set

$$
M\left[:, l-2^{v}+s+1\right]=\left.h_{a, l_{a}}(x)\right|_{\{-1,+1\}^{v}} .
$$

Set $r=r+l_{a}$, and, $s=s+1$.
end
Return M.
Let $m$ denote the number of rows in $\boldsymbol{M}$, and $\left(k_{1}, \cdots, k_{l}\right)$-detection matrix constructed according to Algorithm 2. It is proved in Bshouty (2009) that $m=\frac{2 l \log \frac{k}{L}}{\log l}(1+o(1))$.

## Appendix C. BCH codes

An $[n, k, d]$ binary code $\mathcal{C}$ is a $k$-dimensional linear subspace of the binary vector space of dimension $n$ over $\mathbb{F}_{2}$. The rate of $\mathcal{C}$ is $R \xlongequal{\text { def }} \frac{k}{n}$. The parameter $d$ represents the minimum Hamming distance between the codewords. A generator matrix of $\mathcal{C}$ is a $k \times n$ matrix whose rows span $\mathcal{C}$. A paritycheck matrix of $\mathcal{C}$ is an $(n-k) \times n$ whose rows span the null space of $G$, i.e., $G H^{T}=0$. This implies that $H x=0$ if $x \in \mathcal{C}$. A code with a minimum distance of $d$ can uniquely recover any codeword $x$ in the presence of up to $t=\left\lfloor\frac{d-1}{2}\right\rfloor$ errors. This property implies that all $\mathbf{H e}$ 's are distinct for every $\boldsymbol{e} \in\{0,1\}^{n}$ with Hamming weight at most $t$.

Note that all arithmetic operations are conducted within the field $\mathbb{F}_{2}$, which also implies validity over the real numbers $\mathbb{R}$. Consequently, a parity-check matrix of a $t$-error correcting code satisfies the criteria of an $(n, t)$-search matrix, as outlined in Definition 1.

In particular, we employ binary BCH codes, known for providing an optimal trade-off between minimum distance and code rate for cases where $d \leqslant \frac{n}{\log n}$. Such codes has been also utilized as a building block to construct test matrices for the QGT problem in Karimi et al. (2019b).

Let $n=2^{m}-1$ for for some $m \geqslant 3$ and $t<2^{m-1}$ be an integer. There exists a binary BCH that corrects $t$ errors with $n-k \leqslant m t$ and $d \geqslant 2 t+1$. Let

$$
\tilde{\boldsymbol{H}} \xlongequal{\text { def }}\left[\begin{array}{ccccc}
1 & \gamma & \gamma^{2} & \cdots & \gamma^{n-1}  \tag{28}\\
1 & \gamma^{3} & \left(\gamma^{3}\right)^{2} & \cdots & \left(\gamma^{3}\right)^{n-1} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
1 & \gamma^{2 t-1} & \left(\gamma^{2 t-1}\right)^{2} & \cdots & \left(\gamma^{2 t-1}\right)^{n-1}
\end{array}\right]
$$

where $\gamma$ is a primitive element in $\mathbb{F}_{2^{m}}$. Then, the matrix $\boldsymbol{H}$ that is derived from $\tilde{\boldsymbol{H}}$ by replacing its entries by their corresponding representation in $\mathbb{F}_{2}$ as a column vector is a parity-check matrix of this binary code. Since each element of the the extension field $\mathbb{F}_{2^{m}}$ is represented in $m$ elements in the binary base field, the matrix $\boldsymbol{H}$ has at most $m t=t \log (n+1)$ rows. Equivalently, the matrix $H$ is an $(n, t)$-search matrix with at most $t \log (n+1)$ measurements/rows. Moreover, a binary vector $\boldsymbol{x}$ of Hamming weight at most $t$ can be recover from $H x$ in $O(t n)$ time complexity by utilizing the well known algorithms for decoding BCH codes, e.g., Berlekamp-Massey algorithm Berlekamp (2015); Massey (1969) the Euclidean algorithm Sugiyama et al. (1975) the Berlekamp-Welch algorithm Welch and Berlekamp (1986). We utilize the aforementioned parity-check matrix of the BCH code as a search matrix that solves the QGT problem where the number of defective items is logarithmic in the total number of items $n$.

