# Average-case linear matrix factorization and reconstruction of low width Algebraic Branching Programs 

Neeraj Kayal<br>Microsoft Research India<br>neeraka@microsoft.com

Vineet Nair<br>Indian Institute of Science<br>vineet@iisc.ac.in

Chandan Saha<br>Indian Institute of Science<br>chandan@iisc.ac.in

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#### Abstract

Let us call a matrix $X$ as a linear matrix if its entries are affine forms, i.e. degree one polynomials. What is a minimal-sized representation of a given matrix $F$ as a product of linear matrices? Finding such a minimal representation is closely related to finding an optimal way to compute a given polynomial via an algebraic branching program. Here we devise an efficient algorithm for an average-case version of this problem. Specifically, given $w, d, n \in \mathbb{N}$ and blackbox access to the $w^{2}$ entries of a matrix product $F=X_{1} \cdots X_{d}$, where each $X_{i}$ is a $w \times w$ linear matrix over a given finite field $\mathbb{F}_{q}$, we wish to recover a factorization $F=Y_{1} \cdots Y_{d^{\prime}}$, where every $Y_{i}$ is also a linear matrix over $\mathbb{F}_{q}$ (or a small extension of $\mathbb{F}_{q}$ ). We show that when the input $F$ is sampled from a distribution defined by choosing random linear matrices $X_{1}, \ldots, X_{d}$ over $\mathbb{F}_{q}$ independently and taking their product and $n \geq 4 w^{2}$ and the characteristic of $\mathbb{F}_{q}$ is at least $(n d w)^{\Omega(1)}$ then an equivalent factorization $F=Y_{1} \cdots Y_{d}$ can be recovered in (randomized) time $(w d n \log q)^{O(1)}$. We also show that in this situation, if we are instead given a single entry of $F$ rather than its $w^{2}$ correlated entries then the recovery can be done in (randomized) time $\left(d^{w^{3}} n \log q\right)^{O(1)}$.


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## 1 Introduction

Polynomial matrix factorization. In this paper, we are interested in factorization of a polynomial matrix (that is a matrix with multivariate polynomial entries) into linear matrices, if such a factorization exists. We call this problem linear matrix factorization. It is a natural generalization of the problem of factoring a multivariate polynomial into linear factors for which there is a known efficient randomized algorithm [KT90]. Motivated by applications in control theory, polynomial matrix factorization has been studied in the literature under various restrictions on input and output matrices (see [LMW17] and the references therein). To our knowledge, these restrictions are quite different from the requirement of outputting linear matrix factors of an input polynomial matrix. Our primary motivation for studying this problem stems from the problem of learning or reconstruction of algebraic branching programs (ABPs) - a powerful subclass of arithmetic circuits capturing determinant and iterated matrix multiplication computations (see Definition 1.1).

Reconstruction. Circuit reconstruction is a notable problem in algebraic complexity theory alongside proving lower bounds and polynomial identity testing. Reconstruction of a circuit class $\mathcal{C}$ is the following problem: Given black-box access (i.e. membership query access) to a polynomial function $f$ that is computed by a circuit of size $s$ from $\mathcal{C}$, output a circuit (preferably from $\mathcal{C}$ ) of size not much larger than $s$ (ideally, a polynomial or quasi-polynomial function of $s$ ) computing $f$. As reconstruction of general circuits is believed to be a hard problem, research on reconstruction has focused on interesting restricted models (see the survey [SY10] and the references in [KNST17]), and on the average-case complexity of this problem. In [GKL11] and [GKQ13], average-case reconstruction algorithms were given for multilinear formulas and general formulas respectively under intuitive input distributions. Algebraic branching programs being more powerful than formulas, the problem of efficient average-case reconstruction of ABPs was posed in our earlier work [KNST17] under a natural distribution (see Definition 1.2) ${ }^{1}$. We explain below why the average-case ABP reconstruction problem is interesting to study under this distribution, and in what sense this work is an improvement over [KNST17].

Does lower bound imply reconstruction (even in the average-case)? An intriguing question in circuit complexity is whether or not lower bound implies some kind of learning. More precisely, if there is an explicit function that cannot be computed by circuits in a class $\mathcal{C}$ of size $s$ then can we design an efficient learning algorithm for circuits in $\mathcal{C}$ of size $s$ ? The intuitive reason for expecting a positive answer rests on the high level view that a lower bound proof points to some structural property/weakness of a circuit class and the same property is potentially useful in designing learning algorithms for the class. Indeed, for Boolean circuits, a recent result [CIKK16] has shown that a natural lower bound proof (in the sense of [RR97]) for a circuit class implies quasipolynomial time PAC learning over the uniform distribution for the same class. This generic result is preceded by results (evidences) that hinted at such a connection, like the learning algorithms for $\mathrm{AC}^{0}$ circuits [LMN93] and $\mathrm{AC}^{0}$ circuits with few majority gates [JKS02] ${ }^{2}$. Analogous to Boolean circuits, does a natural lower bound proof (in the sense of [FSV17, GKSS17]) for an algebraic circuit

[^0]class imply efficient reconstruction for the same class? ${ }^{3}$ There are a few results in favor of such a connection, like the reconstruction algorithms for read-once oblivious ABPs, set-multilinear ABPs and non-commutative ABPs [FS13,KS06]. However, there are many other interesting arithmetic circuit classes for which we know of strong lower bounds (that are also algebraically natural), but not efficient reconstruction algorithms. Instances of such classes are homogeneous depth three circuits [NW97], homogeneous depth four circuits [KLSS17, KS17], constant depth multilinear circuits [RY09], multilinear formulas [Raz09], regular formulas [KSS14], and a few other classes [KS16a,KS16b]. Even for more general models like arithmetic formulas and homogeneous ABPs , it makes sense to ask - can we reconstruct sub-quadratic size formulas or sub-quadratic size homogeneous ABPs efficiently? Quadratic lower bounds for formulas and homogeneous ABPs are known [Kal85,Kum17]. In the absence of a generic connection (analogous to [CIKK16]) in the algebraic setting, it would be useful to gather more evidences/examples, perhaps by also moderating the reconstruction setup. Average-case reconstruction is one such natural way to relax the setup ${ }^{4}$. After all, we have little insight into even the 'weaker' question of whether or not lower bound implies average-case reconstruction. However, we should choose an input distribution for an average-case reconstruction problem that is relevant in the context of lower bound.

For the discussion ahead, we denote a $n$-variate, degree- $d$ polynomial as a $(n, d)$-polynomial; a random $(n, d)$-polynomial denotes a $(n, d)$-polynomial with coefficients chosen independently and uniformly at random from $\mathbb{F}$. Assume that $\mathbb{F}$ is a sufficiently large finite field $\mathbb{F}_{q}$, although this requirement is not necessary for the most part of the arguments.

Choosing an input distribution. A lower bound proof for a class $\mathcal{C}$ shows that an explicit $(n, d)$ polynomial is not computable by size-s circuits from $\mathcal{C}$. Such a proof demonstrates some weakness of size-s circuits from $\mathcal{C}$ computing $(n, d)$-polynomials. Typically, the explicit polynomial has degree $d \leq n$ (as in determinant/ permanent [Raz09,RY09, GKKS14] or the Nisan-Wigderson design polynomial [KSS14] or the elementary/power symmetric polynomials [NW97,SW01,Kum17] or a variant of the design polynomial [KST16]), and $s$ is much larger than $n$ (and hence also larger than d). In order to define a corresponding average-case reconstruction problem, we should ideally define an input distribution that is supported on $(n, d)$-polynomials computable by size-s circuits in $\mathcal{C}$; moreover, the distribution should be polynomial-time samplable and reasonably natural. For many circuit classes, defining such a distribution is a bit of a challenge as some of the natural $P$-samplable distributions tend to be primarily supported on $(n, d)$-polynomials, where $d$ is close to the size $s$ of the circuits [GKL11, GKQ13]. But, for some classes, like homogeneous ABPs and homogeneous depth three circuits, these requirements from an input distribution can be mitigated easily. We study the former model in this paper.

## Choosing a distribution on ABPs. A well-known ABP homogenization argument [Nis91] implies

[^1]the following: If a $(n, d)$-polynomial is computable by an ABP A of size ${ }^{5} s$ then it is also computable by an ABP B of width $w \leq s$ and length $d$. If A is a homogeneous $\mathrm{ABP}^{6}$ of size $s$ then B is also a homogeneous ABP of size $s$. In [Kum17], a quadratic lower bound for homogeneous ABP is given by essentially showing that any $(w, d, n)$-ABP computing the power symmetric polynomial $\sum_{i=1}^{n} x_{i}^{d}$ must satisfy $w \geq \frac{n}{2}$, implying that the size of such an ABP is $s \approx w d=\Omega(n d)$. Choosing $d=\Theta(n)$ yields the quadratic bound. A $\Omega\left(n^{1+\epsilon}\right)$ lower bound on $w$, for a constant $\epsilon>0$, would imply a $\Omega\left(n^{1+\epsilon}\right)$ lower bound on the size of general ABPs - such a bound is hitherto unknown. Thus, for the average-case ABP reconstruction problem (see Problem 2), the distribution given in Definition 1.2 is quite appropriate to study as it produces $(n, d)$-polynomials computable by ABPs of size $s \approx w d$ that can potentially be much larger than $n$ and $d$. Further, from the perspective of the quadratic lower bound [Kum17], the average-case ABP reconstruction problem is interesting even for $w=O(n)$. We make progress in this direction by giving a nontrivial ${ }^{7}$ reconstruction algorithm for $w \leq \sqrt{n} / 2$, irrespective of $d$. The algorithm outputs a ( $w, d, n$ )-ABP (with high probability) for the input polynomial chosen according to the distribution.

Comparison to [KNST17]. In [KNST17], we gave a reconstruction algorithm for $w \leq \sqrt{\frac{n}{d}}$. Observe that, under this width constraint, the size $s \approx w d$ of an ABP is upper bounded by max $(n, d)$. Whereas, in this paper we give a reconstruction algorithm for $w \leq \sqrt{n} / 2$ (independent of $d$ ), and hence the size of the ABPs here can be $s=\Theta(\sqrt{n} d)$. To highlight this improvement, if we set $d=\Theta(n)$ (as in several lower bound results [Kum17,KST16,SW01,NW97]) then the width constraint in [KNST17] reduces to $w=O(1)$; moreover, for $d=\Theta(n)$ the size of the ABPs in this work can be $\Theta\left(n^{1.5}\right)$ which is subtantially larger than both $n$ and $d$. On the flip side, the running time of the algorithm in [KNST17] is polynomial in $w, d, n$ and $\log q$, whereas the algorithm here has time complexity $\left(d^{w^{3}} n \log q\right)^{O(1)}$. The exponential dependence on $w^{3}$ comes from a step in our algorithm that solves polynomial equations; all the remaining steps have $(w d n \log q)^{O(1)}$ running time. In fact, the main step (linear matrix factorization) of our algorithm has $(w d n \log q)^{O(1)}$ time complexity (Theorem 1). It may be possible to get around this expensive solvability step and reduce the overall complexity of the algorithm - we leave this as an open question in Section 1.4.

Our proof approach is also quite different from that of [KNST17]. In [KNST17], the Lie algebra of the iterated matrix multiplication polynomial is analyzed to establish a connection between the layer spaces of a full-rank ABP and the irreducible invariant subspaces of the Lie algebra of the polynomial computed by the ABP. This in turn helped reduce the problem to reconstruction of a set-multilinear ABP. We cannot hope to do a similar reduction here as the number of variables is

[^2]much fewer (and independent of $d$ ). Instead, our proof hinges on the following three steps:

1. Showing the uniqueness of the corner spaces when $w \leq \sqrt{n} / 2$, and finding these spaces. This step involves solving polynomial equations.
2. Recovering the intermediate matrices modulo the corner spaces and rearranging them in the correct order. This is the linear matrix factorization step.
3. Completing the affine forms in the intermediate matrices by showing linear independence of the so-called minors of a random $A B P$.

Along the way, we give an efficient equivalence test for the determinant (which is used to get partial access to the intermediate matrices) over finite fields. The details of these steps are given in Section 1.3 and subsequent sections. We think that these steps give us some crucial insights into the structure of a random ABP which may find applications in other similar problems and in resolving some of the questions stated in Section 1.4.

### 1.1 The problems

We study two related problems in this work, average-case matrix factorization and average-case ABP reconstruction. The average-case matrix factorization problem aids us in making progress on average-case ABP reconstruction (see also the remark after Problem 2). The definition of an ABP given below is quite standard and similar to the one stated in [KNST17].

Definition 1.1 (Algebraic branching program). An algebraic branching program (ABP) of width $w$ and length $d$ is a product expression $X_{1} \cdot X_{2} \ldots X_{d}$, where $X_{1}, X_{d}$ are row, column linear matrices over $\mathbb{F}$ of length $w$ respectively, and $X_{i}$ is a $w \times w$ linear matrix over $\mathbb{F}$ for $i \in[2, d-1]$. The polynomial computed by the ABP is the entry of the $1 \times 1$ matrix obtained from the product $\prod_{i=1}^{d} X_{i}$. An ABP of width $w$, length $d$, and in $n$ variables will be called a $(w, d, n)$-ABP over $\mathbb{F}$.

## Remarks:

(a) The iterated matrix multiplication polynomial $\left(\mathrm{IMM}_{w, d}\right)$ is computed by a $(w, d, n)$ - ABP where each entry in $X_{i}$ is a distinct variable, for all $i \in[d]$, and hence $n=w^{2}(d-2)+2 w$.
(b) A polynomial computed by a $(w, d, n)$-ABP can be viewed as an entry of a product of $d, w \times w$ linear matrices $X_{1}, X_{2}, \ldots, X_{d}$. The $w \times w$ matrix $F=X_{1} \cdot X_{2} \ldots X_{d}$ is then called a $(w, d, n)$ matrix product. We note that in the matrix product formulation $X_{1}, X_{d}$ are $w \times w$ linear matrices, while in the ABP formulation $X_{1}, X_{d}$ are row and column linear matrices of length $w$ respectively; hopefully, the context will make the dimensions of these matrices clear.

To study average-case reconstruction for ABP, [KNST17] defined a natural distribution on the polynomials computed by it. The distribution is expressed by a random ( $w, d, n$ )-ABP.

Definition 1.2 (Random ABP and matrix product). A random ( $w, d, n$ )-ABP over $\mathbb{F}$ is a $(w, d, n)$-ABP $X_{1} \cdot X_{2} \ldots X_{d}$ over $\mathbb{F}$, where $X_{i}$ is a random linear matrix chosen independently for every $i \in[d]$. Similarly, a random ( $w, d, n$ )-matrix product over $\mathbb{F}$ is a ( $w, d, n$ )-matrix product $F=X_{1} \cdot X_{2} \ldots X_{d}$ over $\mathbb{F}$, where $X_{i}$ is a random linear matrix chosen independently for every $i \in[d]$.

Having defined the distributions, the two average-case problems can be posed as follows.

Problem 1 (Average-case matrix factorization). Design an algorithm which when given $w, d, n \in \mathbb{N}$, and blackbox access to $w^{2},(n, d)$-polynomials $\left\{f_{s t}\right\}_{s, t \in[w]}$ that constitute the entries of a random $(w, d, n)$-matrix product $F$ over $\mathbb{F}_{q}$, outputs $d, w \times w$ linear matrices $Y_{1}, \ldots, Y_{d}$ over $\mathbb{F}_{q}$ (or a small extension of $\mathbb{F}_{q}$ ) such that $F=Y_{1} \cdot Y_{2} \ldots Y_{d}$, with high probability ${ }^{8}$. The desired running time of the algorithm is $(w d n \log q)^{O(1)}$.

Problem 2 (Average-case ABP reconstruction). Design an algorithm which when given $w, d, n \in$ $\mathbb{N}$, and blackbox access to a $(n, d)$-polynomial $f$ computed by a random $(w, d, n)$-ABP over $\mathbb{F}_{q}$, outputs a $(w, d, n)$-ABP over $\mathbb{F}_{q}$ (or a small extension of $\mathbb{F}_{q}$ ) computing $f$, with high probability. The desired running time of the algorithm is $(w d n \log q)^{O(1)}$.

Remark: In Problem 1 we have blackbox access to $w^{2}$ polynomials constituting the entries of a matrix, whereas in Problem 2 we have blackbox access to a single polynomial. In this sense, Problem 1 is supposedly easier than Problem 2. Still, Problem 1 is of independent interest because if the coefficients of the affine forms are chosen adversarially (instead of randomly) in $X_{1}, X_{2}, \ldots, X_{d}$ then even for $w=3$ the problem becomes as hard as formula reconstruction [BC92].

### 1.2 Our results

Throughout this article, $\mathbb{F}$ will denote $\mathbb{F}_{q}$ with $\operatorname{char}(\mathbb{F}) \geq(w d n)^{7}$, and $\mathbb{L}$ the field $\mathbb{F}_{q^{w}}{ }^{9}$. Also, we will assume $d \geq 5$. Theorem 1 solves Problem 1 for $n \geq 2 w^{2}$.

Theorem 1 (Average-case matrix factorization). For $n \geq 2 w^{2}$, there is a randomized algorithm that takes as input blackbox access to $w^{2},(n, d)$-polynomials $\left\{f_{s t}\right\}_{s, t \in[w]}$ that constitute the entries of a random $(w, d, n)$-matrix product $F=X_{1} \cdot X_{2} \ldots X_{d}$ over $\mathbb{F}$, and with probability $1-(w d n)^{-\Omega(1)}$ returns $w \times w$ linear matrices $Y_{1}, Y_{2}, \ldots, Y_{d}$ over $\mathbb{L}$ satisfying $F=\prod_{i=1}^{d} Y_{i}$. The algorithm runs in $(w d n \log q)^{O(1)}$ time and queries the blackbox at points in $\mathbb{L}^{n}$.

## Remarks:

- The constraint on $\operatorname{char}(\mathbb{F})$ is a bit arbitrary, the results in this paper hold as long as $|\mathbb{F}|$ and $\operatorname{char}(\mathbb{F})$ are sufficiently large polynomial functions in $w, d$ and $n$.
- Uniqueness of factorization: The proof of the theorem shows that there are $C_{i}, D_{i} \in \mathrm{GL}(w, \mathbb{L})$ such that $Y_{i}=C_{i} \cdot X_{i} \cdot D_{i}$, for every $i \in[d]$. Moreover, there are $c_{1}, \ldots, c_{d-1} \in \mathbb{L}^{\times}$satisfying $C_{1}=D_{d}=I_{w}, D_{i} \cdot C_{i+1}=c_{i} I_{w}$ for $i \in[d-1]$, and $\prod_{i=1}^{d-1} c_{i}=1$. At a very high level, it is this uniqueness feature of a random matrix product that guides the algorithm to find a factorization for $F$. In the worst-case, such a factorization need not be unique even if the determinants of the $X_{i}$ 's are coprime irreducible polynomials. For instance ${ }^{10}$,

$$
\left[\begin{array}{ll}
x_{1} & x_{2} \\
x_{3} & x_{4}
\end{array}\right] \cdot\left[\begin{array}{cc}
2 x_{3}-x_{2} & x_{4} \\
x_{1} & x_{3}
\end{array}\right]=\left[\begin{array}{cc}
x_{3} & x_{1} \\
x_{4} & 2 x_{3}-x_{2}
\end{array}\right] \cdot\left[\begin{array}{cc}
x_{1} & x_{2} \\
x_{3} & x_{4}
\end{array}\right]=\left[\begin{array}{cc}
2 x_{1} x_{3} & x_{1} x_{4}+x_{2} x_{3} \\
2 x_{3}^{2}-x_{2} x_{3}+x_{1} x_{4} & 2 x_{3} x_{4}
\end{array}\right]
$$

Using Theorem 1, Theorem 2 addresses Problem 2 for $n \geq 4 w^{2}$.

[^3]Theorem 2 (Average-case ABP reconstruction). For $n \geq 4 w^{2}$, there is a randomized algorithm that takes as input blackbox access to a $(n, d)$-polynomial $f$ computed by a random $(w, d, n)$-ABP over $\mathbb{F}$, and with probability $1-(w d n)^{-\Omega(1)}$ returns a $(w, d, n)$-ABP over $\mathbb{L}$ computing $f$. The algorithm runs in time $\left(d^{w^{3}} n \log q\right)^{O(1)}$ and queries the blackbox at points in $\mathbb{L}^{n}$.

## Remarks:

1. Comparison to [KNST17]: [KNST17] gave an efficient randomized algorithm to solve Problem 2 when $n \geq w^{2} d^{11}$. Theorem 2 improves over [KNST17] by relaxing the constraint on $n$ to $n \geq 4 w^{2}$, but pays in the running time which is exponential in $w^{312}$. Nevertheless, Theorem 2 gives a nontrivial average-case reconstruction algorithm for $w \leq \sqrt{n} / 2$, irrespective of $d$.
2. Time-complexity: There is one step in the algorithm that finds the affine forms in $X_{1}$ and $X_{d}$ by solving systems of polynomial equations over $\mathbb{F}$, and this takes $d^{O\left(w^{3}\right)}$ field operations. Except this step, every other step runs in $(w d n \log q)^{O(1)}$ time. If the complexity of this step is improved then the overall time complexity of the algorithm will also come down.
3. Not pseudorandom: Consider a formal ( $w, d, n$ )-ABP where the coefficients of the affine forms are distinct $\mathbf{y}$-variables, and let $h(\mathbf{x}, \mathbf{y})$ be the polynomial computed by this ABP. Here, $|\mathbf{y}|=$ $(n+1) \cdot\left(w^{2}(d-2)+2 w\right)=m$ (say). If $w=O(\sqrt{n})$, the family $H=\left\{h(\mathbf{x}, \mathbf{b}): \mathbf{b} \in \mathbb{F}^{m}\right\}$ is not pseudorandom under the distribution defined by $\mathbf{b} \in_{r} \mathbb{F}^{m}$. This is because, the $w$ affine forms in $X_{1}$ are linearly independent with high probability. So, the variety of $f=h(\mathbf{x}, \mathbf{b})$ (denoted by $\mathbb{V}(f)$ ) has a subspace of dimension $n-w$ over $\mathbb{F}$; a random polynomial does not have this property with high probability. Using a randomized algorithm (Theorem 2.6 and 3.9 in [HW99]) we can check if $\mathbb{V}(f)$ has a large subspace in $\left(d^{w^{2}} n \log q\right)^{O(1)}$ time. Observe that $\left(d^{w^{2}} n \log q\right)^{O(1)}$ is close to $\exp (n)$ for $w=O(\sqrt{n})$, and so the algorithm does not take time $\gg \exp (n)$ to distinguish $f$ from a random polynomial thereby implying that $H$ is not a pseudorandom family.
4. Comparison to [GKQ13]: [GKQ13] gave an efficient average-case reconstruction algorithm for formulas. Their input is picked from a distribution defined by complete binary trees with alternating layers of + and $\times$ gates and with random affine forms at the leaves. As width- 3 ABPs form a complete model for formulas under p-projections [BC92], Theorem 2 can also be seen as giving another average-case reconstruction algorithm for formulas (when $w=3$ ), albeit with a different input distribution. Our result does not subsume [GKQ13] as the input distributions appear incomparable to us.
The proof of Theorem 1 requires an efficient affine equivalence test for the determinant over finite fields. An $n$-variate polynomial $f(\mathbf{x})$ is affine equivalent to an $m$-variate polynomial $g$, for $n \geq m$, if there is an $A \in \mathbb{F}^{m \times n}$ of rank $m$ and an $\mathbf{a} \in \mathbb{F}^{m}$ such that $f=g(A \cdot \mathbf{x}+\mathbf{a})$. Further, for $m=n$, $f$ is equivalent to $g$ if there is an $A \in \mathrm{GL}(n, \mathbb{F})$ such that $f=g(A \cdot \mathbf{x})$. Given blackbox access to a $(n, w)$-polynomial $f$, where $n \geq w^{2}$, the affine equivalence test problem for the determinant is to check whether $f$ is affine equivalent to $\operatorname{Det}_{w}$, and if yes then output a $B \in \mathbb{F}^{w^{2} \times n}$ of rank $w^{2}$ and $\mathbf{a} \mathbf{b} \in \mathbb{F}^{w^{2}}$ such that $f=\operatorname{Det}_{w}(B \cdot \mathbf{x}+\mathbf{b})$. The algorithm in the theorem below almost solves this problem over finite fields - it returns a $B \in \mathbb{L}^{w^{2} \times n}$ of rank $w^{2}$ and a $\mathbf{b} \in \mathbb{L}^{w w^{2}}$.
[^4]Theorem 3 (Determinant equivalence test). There is a randomized algorithm that takes as input blackbox access to a $(n, w)$-polynomial $f \in \mathbb{F}[\mathbf{x}]$, where $n \geq w^{2}$, and does the following with probability $1-\frac{n^{O(1)}}{q}$ : If $f$ is affine equivalent to $\operatorname{Det}_{w}$ then it outputs $a B \in \mathbb{L}^{w^{2} \times n}$ of rank $w^{2}$ and $a \mathbf{b} \in \mathbb{L}^{w^{2}}$ such that $f=\operatorname{Det}_{w}(B \cdot \mathbf{x}+\mathbf{b})$, else it outputs ' $f$ not affine equivalent to $\operatorname{Det}_{w}$ '. The algorithm runs in $(n \log q)^{O(1)}$ time and queries the blackbox at points in $\mathbb{L}^{n}$.

## Remarks:

1. Comparison to [Kay12]: An efficient equivalence test for the determinant over $\mathbb{C}$ was given in [Kay12]. The computation model in [Kay12] assumes that arithmetic over $\mathbb{C}$ and root finding of univariate polynomials over $\mathbb{C}$ can be done efficiently. While we follow the general strategy of analyzing the Lie algebra of the determinant and reduction to PS-equivalence from [Kay12], our algorithm is somewhat simpler: Unlike [Kay12], our algorithm does not involve the Cartan subalgebras and is almost the same as the simpler equivalence test for the permanent polynomial in [Kay12]. The simplification is achieved by showing that the characteristic polynomial of a random element of the Lie algebra of Det ${ }_{w}$ splits completely over $\mathbb{L}$ with high probability (Lemma 5.2) - this is crucial for Theorem 1 as it allows the algorithm to output a matrix factorization over a fixed low extension of $\mathbb{F}$, namely $\mathbb{L}$.
2. Average-case $A B P$ reconstruction over $\mathbb{Q}$ : In our arguments, Theorem 3 is the only place where we need the underlying field is finite. In other words, the algorithms in Theorems 1 and 2 work over $\mathbb{Q}$ if only there is an efficient equivalence test for $\operatorname{Det}_{w}$ over $\mathbb{Q}$. Also, if there is an affine equivalence test for $\operatorname{Det}_{w}$ that outputs $B, \mathbf{b}$ over the base field $(\mathbb{Q}$ or $\mathbb{F})$ then the algorithm in Theorem 2 would output an ABP over the base field.

### 1.3 Algorithms and their analysis

The algorithms mentioned in Theorem 1 and 2 are given in Algorithm 1 and 2, respectively. In this section, we briefly discuss their correctness and complexity - for the missing details, we allude to the relevant parts of the subsequent sections.

### 1.3.1 Analysis of Algorithm 1

Since $F=X_{1} \cdot X_{2} \ldots X_{d}$ is a random $(w, d, n)$-matrix product, with probability $1-(w d n)^{-\Omega(1)}$, the following property is satisfied: Every $X_{i}$ is a full rank linear matrix (that is the affine forms in $X_{i}$ are $\mathbb{F}$-linearly independent), and $\operatorname{det}\left(X_{1}\right), \operatorname{det}\left(X_{2}\right), \ldots, \operatorname{det}\left(X_{d}\right)^{13}$ are coprime irreducible polynomials (see Claim 2.3). We analyze Algorithm 1 assuming that this property of the input is satisfied. Algorithm 1 has three main stages:

1. Computing the irreducible factors of $\operatorname{det}(F)$ (Steps 2-6): From blackbox access to the entries of $F$, a blackbox access to $\operatorname{det}(F)$ is computed in $(w d n \log q)^{O(1)}$ time using Gaussian elimination. Subsequently, using Kaltofen-Trager's factorization algorithm [KT90], blackbox access to the irreducible factors $g_{1}, g_{2}, \ldots, g_{d}$ of $\operatorname{det}(F)$ are constructed in $(w d n \log q)^{O(1)}$ time (see Lemma 2.1). Since $\operatorname{det}\left(X_{1}\right), \ldots, \operatorname{det}\left(X_{d}\right)$ are coprime irreducible polynomials, there is a permutation $\sigma$ of $[d]$, and $c_{i} \in \mathbb{F}^{\times}$for all $i \in[d]$, such that $c_{i} \cdot \operatorname{det}\left(X_{i}\right)=g_{\sigma(i)}$ and $\prod_{i=1}^{d} c_{i}=1$. For the next
[^5]```
Algorithm 1 Average-case matrix factorization
    INPUT: Blackbox access to \(w^{2},(n, d)\)-polynomials \(\left\{f_{s t}\right\}_{s, t \in[w]}\) that constitute the entries of a ran-
    \(\operatorname{dom}(w, d, n)\) - matrix product \(F=X_{1} \cdot X_{2} \ldots X_{d}\).
    OUTPUT: Linear matrices \(Y_{1}, Y_{2}, \ldots, Y_{d}\) over \(\mathbb{L}\) such that \(F=Y_{1} \cdot Y_{2} \ldots Y_{d}\).
    /* Factorization of the determinant */
    Compute blackbox access to \(\operatorname{det}(F)\).
    Compute blackbox access to the irreducible factors of \(\operatorname{det}(F)\); call them \(g_{1}, g_{2}, \ldots, g_{d}\).
    if the number of irreducible factors is not equal to \(d\) then
        Output 'Failed'.
    end if
    /* Affine equivalence test for determinant */
    Set \(j=1\).
    while \(j \leq d\) do
        Call the algorithm in Theorem 3 with input as blackbox access to \(g_{j}\); let \(B_{j}\) and \(\mathbf{b}_{j}\) be its
        output. Construct the \(w \times w\) full-rank linear matrix \(Z_{j}\) over \(\mathbb{L}\) determined by \(B_{j}\) and \(\mathbf{b}_{j}\).
        if the algorithm outputs ' \(g_{j}\) not affine equivalent to \(\operatorname{Det}_{w}{ }^{\prime}\) ' then
            Output 'Failed'.
        end if
        Set \(j=j+1\).
    end while
    /* Rearrangement of the matrices */
    Call Algorithm 3 on input blackbox access to \(F\) and \(Z_{1}, \ldots, Z_{d}\), and let \(Y_{1}, \ldots, Y_{d}\) be its output.
    if Algorithm 3 outputs 'Rearrangement not possible' then
        Output 'Failed'.
    end if
23.
24. Output \(Y_{1}, Y_{2}, \ldots, Y_{d}\).
```

two stages, assume $w>1$ as the $w=1$ case gets solved readily at this stage.
2. Affine equivalence test (Steps 9-16): Let $j=\sigma(i)$ and $X_{i}^{\prime}$ be the matrix $X_{i}$ with the affine forms in the first row multiplied by $c_{i}$. Then, $g_{j}=\operatorname{det}\left(X_{i}^{\prime}\right)=c_{i} \cdot \operatorname{det}\left(X_{i}\right)$, which is affine equivalent to $\operatorname{Det}_{w}$. At step 11, the algorithm ${ }^{14}$ in Theorem 3 finds a $B_{j} \in \mathbb{L}^{w^{2} \times n}$ of rank $w^{2}$ and $\mathbf{b}_{j} \in \mathbb{L}^{w w^{2}}$ such that $g_{j}=\operatorname{Det}_{w}\left(B_{j} \cdot \mathbf{x}+\mathbf{b}_{j}\right)$, with probability $1-(w d n)^{-\Omega(1)}$. Let $Z_{j}$ be the matrix obtained by appropriately replacing the entries of the $w \times w$ symbolic matrix with the affine forms in $B_{j} \cdot \mathbf{x}+\mathbf{b}_{j}$ such that $\operatorname{det}\left(Z_{j}\right)=g_{j}=\operatorname{det}\left(X_{i}^{\prime}\right)$. This certifies that there are matrices $C_{i}, D_{i} \in \operatorname{SL}(w, \mathbb{L})$ satisfying, $Z_{j}=C_{i} \cdot X_{i}^{\prime} \cdot D_{i}$ or $Z_{j}^{T}=C_{i} \cdot X_{i}^{\prime} \cdot D_{i}$ (see Fact 1 in Section 5.1). Multiplying the first column of $C_{i}$ with $c_{i}$, and calling the resulting matrix $C_{i}$ again, we see that there are matrices $C_{i}, D_{i} \in \mathrm{GL}(w, \mathbb{L})$ satisfying, $Z_{j}=C_{i} \cdot X_{i} \cdot D_{i}$ or $Z_{j}^{T}=C_{i} \cdot X_{i} \cdot D_{i}$. Observe that such $C_{i}, D_{i}$ are unique up to multiplications by elements in $\mathbb{L}^{\times 15}$.
3. Rearrangement of the retrieved matrices (Steps 19-22): At step 19, Algorithm 3 constructs the matrices $Y_{1}, Y_{2}, \ldots, Y_{d}$ by determining the permutation $\sigma$ and whether $Z_{\sigma(i)}=C_{i} \cdot X_{i} \cdot D_{i}$ or $Z_{\sigma(i)}^{T}=C_{i} \cdot X_{i} \cdot D_{i}$. Internally, Algorithm 3 uses Algorithm 4 which when given blackbox access to $F_{d}=F$ and a $Z$ (that is either $Z_{k}$ or $Z_{k}^{T}$ for some $k \in[d]$ ), does the following with probability $1-(w d n)^{-\Omega(1)}$ : If $Z=C_{d} \cdot X_{d} \cdot D_{d}$ then it outputs a $\tilde{D}_{d}=a_{d} D_{d}$ for some $a_{d} \in \mathbb{L}^{\times}$. For all other cases - if $Z=C_{i} \cdot X_{i} \cdot D_{i}$ or $Z^{T}=C_{i} \cdot X_{i} \cdot D_{i}$ for $i \in[d-1]$, or $Z^{T}=C_{d} \cdot X_{d} \cdot D_{d}$ - it outputs 'Failed'. Algorithm 4 uses the critical fact that $F$ is a random matrix product to accomplish the above and locate the unique last matrix. The running time of the algorithm, which is $(w d n \log q)^{O(1)}$, and its proof of correctness ${ }^{16}$ are discussed in Section 3.2. Algorithm 3 calls Algorithm 4 on inputs $F, Z_{k}$ and $F, Z_{k}^{T}$ for all $k \in[d]$. If Algorithm 4 returns a matrix $\tilde{D}_{d}$ for some $k \in[d]$ on either inputs $F, Z_{k}$ or $F, Z_{k}^{T}$ then it sets $M_{d}=Z_{k}$ or $M_{d}=Z_{k}^{T}$ respectively, and $\sigma(d)=k$. Subsequently, Algorithm 3 computes blackbox access to a length $d-1$ matrix product $F_{d-1}=F \cdot \tilde{D}_{d} \cdot M_{d}^{-1}=X_{1} \cdots X_{d-2} \cdot\left(X_{d-1} \cdot a_{d} C_{d}^{-1}\right)$, and repeats the above process to compute $M_{d-1}$ and $\sigma(d-1)$ with the inputs $F_{d-1}$ and $\left\{Z_{1}, \ldots, Z_{d}\right\} \backslash Z_{\sigma(d)}$. Thus, using Algorithm 4 repeatedly, Algorithm 3 iteratively determines $\sigma$ and $M_{d}, M_{d-1}, \ldots, M_{2}$ : At the $(d-t+1)$-th iteration, for $t \in[d-1,2]$, it computes a matrix $\tilde{D}_{t}=a_{t}\left(C_{t+1} \cdot D_{t}\right)$ for some $a_{t} \in \mathbb{L}^{\times}$, sets $M_{t}$ and $\sigma(t)$ accordingly, creates blackbox access to $F_{t-1}=F_{t} \cdot \tilde{D}_{t} \cdot M_{t}^{-1}$ and prepares the list $\left\{Z_{1}, \ldots, Z_{d}\right\} \backslash\left\{Z_{\sigma(d)}, Z_{\sigma(d-1)}, \ldots, Z_{\sigma(t)}\right\}$ for the next iteration. Finally, setting $Y_{1}=F_{1}$ and $Y_{i}=M_{i} \cdot \tilde{D}_{i}^{-1}$, for all $i \in[2, d]$, we have $F=\prod_{i=1}^{d} Y_{i}$.

### 1.3.2 Analysis of Algorithm 2

Let $f$ be the polynomial computed by a $(w, d, n)$-ABP $X_{1} \cdot X_{2} \ldots X_{d}$. We can assume that $f$ is a homogeneous degree- $d$ polynomial and the entries in each $X_{i}$ are linear forms (i.e., affine forms with constant term zero), owing to the following simple homogenization trick.

[^6]Homogenization of $A B P$ : Consider the $(n+1)$-variate homogeneous degree- $d$ polynomial

$$
f_{\mathrm{hom}}=x_{0}^{d} \cdot f\left(\frac{x_{1}}{x_{0}}, \frac{x_{2}}{x_{0}}, \ldots, \frac{x_{n}}{x_{0}}\right) .
$$

The polynomial $f_{\text {hom }}$ is computable by the $(w, d, n)$ - $\mathrm{ABP} X_{1}^{\prime} \cdot X_{2}^{\prime} \ldots X_{d^{\prime}}^{\prime}$, where $X_{i}^{\prime}$ is equal to $X_{i}$ but with the constant term in the affine forms multiplied by $x_{0}$. If we construct an ABP for $f_{\text {hom }}$ then an ABP for $f$ is obtained by setting $x_{0}=1$.

```
Algorithm 2 Average-case ABP reconstruction
    INPUT: Blackbox access to a \((n, d)\)-polynomial \(f\) computed by a random \((w, d, n)\)-ABP.
    OUTPUT: A \((w, d, n)\)-ABP over \(\mathbb{L}\) computing \(f\).
    /* Computing the corner spaces */
    Call Algorithm 5 on \(f\) to compute bases of the two unique \(w\)-dimensional \(\mathbb{F}\)-linear spaces \(\mathcal{X}_{1}\)
    and \(\mathcal{X}_{d}\), spanned by linear forms in \(\mathbb{F}[\mathbf{x}]\), such that \(f\) is zero modulo each of \(\mathcal{X}_{1}\) and \(\mathcal{X}_{d}\).
    if Algorithm 5 outputs 'Failed' then
        Output 'Failed to construct an \(\mathrm{ABP}^{\prime}\).
    end if
    Compute a transformation \(A \in \operatorname{GL}(n, \mathbb{F})\) that maps the bases of \(\mathcal{X}_{1}\) and \(\mathcal{X}_{d}\) to distinct variables
    \(\mathbf{y}=\left\{y_{1}, y_{2}, \ldots, y_{w}\right\}\) and \(\mathbf{z}=\left\{z_{1}, z_{2}, \ldots z_{w}\right\}\) respectively, where \(\mathbf{y}, \mathbf{z} \subseteq \mathbf{x}\). Let \(\mathbf{r}=\mathbf{x} \backslash(\mathbf{y} \uplus \mathbf{z})\),
    \(X_{1}^{\prime}=\left(\begin{array}{llll}y_{1} & y_{2} & \ldots & y_{w}\end{array}\right), X_{d}^{\prime}=\left(\begin{array}{llll}z_{1} & z_{2} & \ldots & z_{w}\end{array}\right)^{T}\) and \(f^{\prime}=f(A \cdot \mathbf{x})\).
```

    7.
    . /* Computing the coefficients of the r variables */
    Construct blackbox access to the \(w^{2}\) polynomials that constitute the entries of the \(w \times w\) matrix
    \(F=\left(\left.\frac{\partial f^{\prime}}{\partial y_{s} z_{t}}\right|_{\mathbf{y}=0, \mathbf{z}=0}\right)_{s, t \in[w]}\).
    Call Algorithm 1 on input \(F\) to compute a factorization of \(F\) as \(S_{2} \cdot S_{3} \ldots S_{d-1}\).
    if Algorithm 1 outputs 'Failed' then
        Output 'Failed to construct an ABP'.
    end if
    /* Computing the coefficients of the \(\mathbf{y}\) and \(\mathbf{z}\) variables */
    Call Algorithm 6 on inputs \(f^{\prime}\) and \(\left\{S_{2}, S_{3}, \ldots, S_{d-1}\right\}\) to compute matrices \(T_{2}, T_{3}, \ldots, T_{d-1}\) such
    that \(f^{\prime}\) is computed by the ABP \(X_{1}^{\prime} \cdot T_{2} \cdots T_{d-1} \cdot X_{d}^{\prime}\).
    if Algorithm 6 outputs 'Failed' then
        Output 'Failed to construct an ABP'.
    end if
    Apply the transformation \(A^{-1}\) on the \(\mathbf{x}\) variables in the matrices \(X_{1}^{\prime}, X_{d}^{\prime}\), and \(T_{k}\) for \(k \in[2, d-1]\).
    Call the resulting matrices \(Y_{1}, Y_{d}\), and \(Y_{k}\) for \(k \in[2, d-1]\) respectively.
    21. Output $Y_{1} \cdot Y_{2} \ldots Y_{d}$ as the ABP computing $f$.

We give an overview of the three main stages in Algorithm 2. As in Algorithm 1, the matrices $X_{1}, X_{2}, \ldots, X_{d}$ are assumed to be full rank linear matrices and further, for a similar reason, the $2 w$ linear forms in $X_{1}$ and $X_{d}$ are assumed to be $\mathbb{F}$-linearly independent.

1. Computing the corner spaces (Steps 2-6): Polynomial $f$ is zero modulo each of the two wdimensional $\mathbb{F}$-linear spaces $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ spanned by the linear forms in $X_{1}$ and $X_{d}$ respec-
tively ${ }^{17}$. We show in Lemma 4.1, if $n \geq 4 w^{2}$ then with probability $1-(w d n)^{-\Omega(1)}$ the following holds: Let $\mathbb{K} \supseteq \mathbb{F}$ be any field. If $f=0 \bmod \left\langle l_{1}, \ldots, l_{w}\right\rangle$, where $l_{i}$ 's are linear forms in $\mathbb{K}[\mathbf{x}]$, then the $l_{i}$ 's either belong to the $\mathbb{K}$-span of the linear forms in $X_{1}$ or belong to the $\mathbb{K}$ span of the linear forms in $X_{d}$. In this sense, the spaces $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ are unique. The algorithm invokes Algorithm 5 which computes bases of $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ by solving $O(n)$ systems of polynomial equations over $\mathbb{F}$. Such a system has $d^{O\left(w^{2}\right)}$ equations in $m=O\left(w^{3}\right)$ variables and the degree of the polynomials in the system is at most $d$; we intend to find all the solutions in $\mathbb{F}^{m}$. It turns out that owing to the uniqueness of $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$, the variety over $\overline{\mathbb{F}}{ }^{18}$ defined by such a system has exactly two points and these points lie in $\mathbb{F}^{m}$. From the two solutions, bases of $\mathcal{X}_{1}$ and of $\mathcal{X}_{d}$ can be derived. The two solutions of the system are computed by a randomized algorithm running in $\left(d^{w^{3}} \log q\right)^{O(1)}$ time ([Ier89, HW99], see Lemma 2.2) - the algorithm exploits the fact that the variety over $\overline{\mathbb{F}}$ is zero-dimensional. Thus, at step 2 , the two spaces are either equal to $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ or $\mathcal{X}_{d}$ and $\mathcal{X}_{1}$ respectively. Without loss of generality, we assume the former. Once bases of the corner spaces $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ are computed, an invertible transformation $A$ maps the linear forms in the bases to distinct variables (as the linear forms in $X_{1}$ and $X_{d}$ are $\mathbb{F}$-linearly independent).
2. Computing the coefficients of the $\mathbf{r}$ variables (Steps 9-13): There is an ABP $X_{1}^{\prime} \cdot X_{2}^{\prime} \ldots X_{d}^{\prime}$ computing $f^{\prime}=f(A \cdot \mathbf{x})$, where $X_{1}^{\prime}$ and $X_{d}^{\prime}$ are equal to $\left(y_{1} y_{2} \ldots y_{w}\right)$ and $\left(z_{1} z_{2} \ldots z_{w}\right)^{T}$ respectively. For $k \in[2, d-1]$, let $R_{k}=\left(X_{k}^{\prime}\right)_{\mathbf{y}=0, \mathbf{z}=0}{ }^{19}$ and $F=R_{2} \cdot R_{3} \ldots R_{d-1}$. As $X_{1} \cdot X_{2} \ldots X_{d}$ is a random $(w, d, n)$ - $\mathrm{ABP}, R_{2} \cdot R_{3} \ldots R_{d-1}$ is a random $(w, d-2, n-2 w)$-matrix product over $\mathbb{F}$. The $(s, t)$-th entry of $F$ is equal to $\left(\frac{\partial f^{\prime}}{\partial y_{s} z_{t}}\right)_{\mathbf{y}=0, \mathbf{z}=0}$, for $s, t \in[w]$. Blackbox access to each of the $w^{2}$ entries of $F$ are constructed in $(w d n \log q)^{O(1)}$ time using Claim 2.1. From $F$, Algorithm 1 computes linear matrices $S_{2}, \ldots, S_{d-1}$ over $\mathbb{L}$ in $\mathbf{r}=\mathbf{x} \backslash(\mathbf{y} \uplus \mathbf{z})$ variables such that $F=S_{2} \cdot S_{3} \ldots S_{d-1}$. Moreover, the uniqueness of factorization implies there are linear matrices $T_{2}, \ldots, T_{d-1}$ over $\mathbb{L}$ in the $\mathbf{x}$-variables, satisfying $\left(T_{k}\right)_{\mathbf{y}=0, \mathbf{z}=0}=S_{k}$, such that $f^{\prime}$ is computed by the ABP $X_{1}^{\prime} \cdot T_{2} \cdots T_{d-1} \cdot X_{d}^{\prime}$.
3. Computing the coefficients of $\mathbf{y}$ and $\mathbf{z}$ variables in $T_{k}$ (Steps 16-20): Algorithm 6 finds the coefficients of the $\mathbf{y}$ and $\mathbf{z}$ variables in the linear forms present in $T_{2}, \ldots, T_{d-1}$ in $(w d n \log q)^{O(1)}$ time. We present the idea here; the detail proof of correctness is given in Section 4.2. In the following discussion, $M(i, j)$ denotes the $(i, j)$-th entry, $M(i, *)$ the $i$-th row, and $M(*, j)$ the $j$-th column of a linear matrix $M$. Let us focus on finding the coefficients of $y_{1}$ in the linear forms present in $T_{2}(1, *), T_{3}, \ldots, T_{d-2}, T_{d-1}(*, 1)$. There are $w^{2}(d-4)+2 w$ linear forms in these matrices and these would be indexed by $\left[w^{2}(d-4)+2 w\right]$. Let $c_{e}$ be the coefficient of $y_{1}$ in the $e$-th linear form $l_{e}$ for $e \in\left[w^{2}(d-4)+2 w\right]$. We associate a polynomial $h_{e}(\mathbf{r})$ in $\mathbf{r}$ variables with $l_{e}$ as follows: If $l_{e}$ is the $(i, j)$-th entry of $T_{k}$ then $h_{e} \stackrel{\text { def }}{=}$ $\left[S_{2}(1, *) \cdot S_{3} \cdots S_{k-2} \cdot S_{k-1}(*, i)\right] \cdot\left[S_{k+1}(j, *) \cdot S_{k+2} \cdots S_{d-2} \cdot S_{d-1}(*, 1)\right]{ }^{20}$. Observe that if $f^{\prime}$ is treated as a polynomial in $\mathbf{y}$ and $\mathbf{z}$ variables with coefficients in $\mathbb{L}(\mathbf{r})$ then the coefficient of $y_{1}^{2} z_{1}$ is exactly $\sum_{e \in\left[w^{2}(d-4)+2 w\right]} c_{e} \cdot h_{e}(\mathbf{r})$. On the other hand, this coefficient is $\left(\frac{\partial f^{\prime}}{\partial y_{1}^{2} z_{1}}\right)_{\mathbf{y}=0, \mathbf{z}=0^{\prime}}$

[^7]for which we can obtain blackbox access using Claim 2.1. This allows us to write the equation,
\[

$$
\begin{equation*}
\sum_{e=1}^{w^{2}(d-4)+2 w} c_{e} \cdot h_{e}(\mathbf{r})=\left(\frac{\partial f^{\prime}}{\partial y_{1}^{2} z_{1}}\right)_{\mathbf{y}=0, \mathbf{z}=0} . \tag{1}
\end{equation*}
$$

\]

We show in Lemma 4.2 and Corollary 4.1 that the polynomials $h_{e}$, for $e \in\left[w^{2}(d-4)+\right.$ $2 w]$, are $\mathbb{L}$-linearly independent with probability ${ }^{21} 1-(w d n)^{-\Omega(1)}$. By substituting random values to the $\mathbf{r}$ variables in the above equation, we can set up a system of $w^{2}(d-4)+2 w$ linear equations in the $c_{e}$ 's. The linear independence of the $h_{e}{ }^{\prime}$ s ensures that we can solve for $c_{e}$ (by Claim 2.2).

### 1.3.3 Proof strategy for Theorem 3

The algorithm in Theorem 3 has three stages:

1. Reduction to equivalence testing: Applying known techniques - 'variable reduction' (Claim 5.1) and 'translation equivalence' (Claim 5.2) - the affine equivalence testing problem is efficiently reduced to equivalence testing for Det $_{w}$ with high probability. An equivalence test takes blackbox access to a $w^{2}$-variate polynomial $g(\mathbf{y})$ as input and does the following with high probability: If $g$ is equivalent to Det $_{w}$ then it outputs a $Q \in \operatorname{GL}\left(w^{2}, \mathbb{L}\right)$ such that $g=\operatorname{Det}_{w}(Q \cdot \mathbf{y})$ else it outputs ' $g$ not equivalent to $\operatorname{Det}_{w}$ '.
2. Reduction to PS-equivalence: The reduction is given in Algorithm 7. The algorithm proceeds by computing an $\mathbb{F}$-basis of the Lie algebra of the group of symmetries of $g$ (denoted as $\mathfrak{g}_{g}$, see Claim 5.3). It then picks an element $F$ uniformly at random from $\mathfrak{g}_{g}$ and computes its characteristic polynomial $h(x)$. Since $F \in \mathfrak{g}_{g}$, it is similar to a $L \in \mathfrak{g}_{\text {Det }_{w v}}$ (see Fact 3 in Section 5.1), implying that their characteristic polynomials are equal. As $F$ is a random element of $\mathfrak{g}_{g}, L$ is also a random element of $\mathfrak{g}_{\operatorname{Det}_{w}}$. In Lemma 5.2, we show that the characteristic polynomial $h$ of a $L \in_{r} \mathfrak{g}_{\text {Det }_{w}}$ is square-free and splits completely over $\mathbb{L}$, with high probability ${ }^{22}$. The roots of $h$ are computed in randomized $(w \log q)^{O(1)}$ time ( [CZ81], see also [vzGG03]). From the roots, a $D \in \mathrm{GL}\left(w^{2}, \mathbb{L}\right)$ can be computed such that $D^{-1} F D$ is diagonal ${ }^{23}$. Thereafter, the structure of the group of symmetries of $\operatorname{Det}_{w}$ and its Lie algebra helps argue, in Section 5.2, that $f(D \cdot \mathbf{x})$ is PS-equivalent to $\operatorname{Det}_{w}$.
3. Doing the PS-equivalence: This step follows directly from [Kay12] (see Lemma 5.1).

### 1.4 Few questions

The following questions are immediate from the above discussions:
(a) Can we compute the corner spaces in $(w d \log q)^{O(1)}$ time? If so then the overall complexity of the algorithm would come down to $(w d \log q)^{O(1)}$.

[^8](b) In the equivalence test for the determinant, can we output a linear matrix over the base field $\mathbb{F}$ instead of a matrix over the extension $\mathbb{L}$ ?
(c) Is it possible to do nontrivial reconstruction in the average-case when $w$ is significantly larger than $\sqrt{n}$, say for $w=\frac{n}{2}$ ?
(d) For $w$ significantly larger than $\sqrt{n}$, say $w=n^{2}$, can we show that linear factorization of a random ( $w, d, n$ )-matrix product is unique (in the sense as in the second remark after Theorem 1)?

## 2 Preliminaries

### 2.1 Notations

$\mathrm{GL}(w, \mathbb{F})$ is the set of $w \times w$ invertible matrices over $\mathbb{F}$, and $\operatorname{SL}(w, \mathbb{F})$ the set of $w \times w$ matrices over $\mathbb{F}$ with determinant one. Bold letters $\mathbf{x}, \mathbf{y}, \mathbf{z}, \mathbf{u}, \mathbf{v}, \mathbf{w}$ are used to represent either column vectors (or sets) of variables or column vectors of field elements, calligraphic letters like $\mathcal{X}$ to represent vector spaces, capital letters like $A, B, C, S, T$ for matrices or sets - the context of a usage of any of these symbols would hopefully make its purpose clear. The derivative of a polynomial $f$ with respect to a monomial $\mu$ is denoted as $\frac{\partial f}{\partial \mu}$ or $\partial_{\mu} f$.

### 2.2 Algorithmic preliminaries

The following result on blackbox polynomial factorization is proved in [KT90].
Lemma 2.1 ( [KT90]). There is a randomized algorithm that takes as input blackbox access to a ( $n, d$ )polynomial $f$ over $\mathbb{F}$, and constructs blackbox access to the irreducible factors of $f$ over $\mathbb{F}$ in $(n d \log q)^{O(1)}$ time with success probability $1-\frac{(n d)^{O(1)}}{q}$.

Let $I$ be an ideal of $\mathbb{F}[\mathbf{x}]$ generated by $(n, d)$-polynomials $g_{1}, \ldots, g_{m}$, and $\mathbb{V}_{\overline{\mathbb{F}}}(I)$ the variety or the algebraic set defined by $I$ over $\overline{\mathbb{F}} . \mathbb{V}_{\overline{\mathbb{F}}}(I)$ is zero-dimensional if it has finitely many points. We say a point $\mathbf{a} \in \mathbb{V}_{\overline{\mathbb{F}}}(I)$ is $\mathbb{F}$-rational if $\mathbf{a} \in \mathbb{F}^{n}$. The proof of the next result follows from [Ier89] (see also [HW99]).

Lemma 2.2 ([Ier89]). There is a randomized algorithm that takes input $m,(n, d)$-polynomials $g_{1}, g_{2}, \ldots, g_{m}$ generating an ideal I of $\mathbb{F}[\mathbf{x}]$. If $\mathbb{V}_{\overline{\mathbb{F}}}(I)$ is zero-dimensional and all points in it are $\mathbb{F}$-rational then the algorithm computes all the points in $\mathbb{V}_{\overline{\mathbb{F}}}(I)$ with probability $1-\exp (-m n d \log q)$. The running time of the algorithm is $\left(m d^{n} \log q\right)^{O(1)} .24$

### 2.3 A few useful facts

We list down three claims (without proofs) that will be used in the later sections. A proof of the first can be given using interpolation. Proofs of the last two follow from applications of the Schwartz-Zippel lemma [Sch80,Zip79].

[^9]Claim 2.1. There is a deterministic algorithm that given blackbox access to a $(n, d)$-polynomial $f \in \mathbb{F}[\mathbf{x}]$, and a monomial $\mu$ of constant degree in $\mathbf{x}$, computes blackbox access to $\partial_{\mu} f$ in $(n d \log q)^{O(1)}$ time.

Claim 2.2. Let $f_{1}, f_{2}, \ldots, f_{m}$ be $\mathbb{F}$-linearly independent $(n, d)$-polynomials in $\mathbb{F}[\mathbf{x}]$. If $\mathbf{a}_{1}, \mathbf{a}_{2}, \ldots, \mathbf{a}_{m}$ are points in $\mathbb{F}^{n}$ chosen independently and uniformly at random, then the matrix $\left(f_{t}\left(\mathbf{a}_{s}\right)\right)_{s, t \in[m]}$ has rank $m$ over $\mathbb{F}$ with probability at least $1-\frac{d m}{q}$.

Claim 2.3. Let $X_{1} \cdot X_{2} \ldots X_{d}$ be a random ( $w, d, n$ )-matrix product over $\mathbb{F}$. If $n \geq w^{2}$ then $X_{1}, X_{2}, \ldots, X_{d}$ are full rank linear matrices and $\operatorname{det}\left(X_{1}\right), \operatorname{det}\left(X_{2}\right), \ldots, \operatorname{det}\left(X_{d}\right)$ are coprime irreducible polynomials with probability $1-(w d n)^{-\Omega(1)}$.

## 3 Average-case matrix factorization: Proof of Theorem 1

The algorithm in Theorem 1 is presented in Algorithm 1. To complete the analysis, given in Section 1.3.1, we need to argue the correctness of the key step of rearrangement of the matrices (Algorithm 3) by finding the last matrix (Algorithm 4). As the functioning of Algorithm 3 is already sketched out in Section 1.3.1, the reader may skip to Section 3.2. For completeness, we include an analysis of Algorithm 3 in the following subsection.

### 3.1 Rearranging the matrices

Recall, we have assumed $F$ is a $(w, d, n)$-matrix product $X_{1} \cdot X_{2} \ldots X_{d}$, where $X_{1}, X_{2}, \ldots, X_{d}$ are full rank linear matrices, and $\operatorname{det}\left(X_{1}\right), \operatorname{det}\left(X_{2}\right), \ldots, \operatorname{det}\left(X_{d}\right)$ are coprime irreducible polynomials. The inputs to Algorithm 3 are $d$ full rank linear matrices $Z_{1}, Z_{2}, \ldots, Z_{d}$ over $\mathbb{L}$ such that there are matrices $C_{i}, D_{i} \in G L(w, \mathbb{L})$ and a permutation $\sigma$ of $[d]$ satisfying $Z_{\sigma(i)}=C_{i} \cdot X_{i} \cdot D_{i}$ or $Z_{\sigma(i)}^{T}=C_{i} \cdot X_{i} \cdot D_{i}$ for every $i \in[d]$. Algorithm 3 iteratively determines $\sigma$ (implicitly) by repeatedly using Algorithm 4. The behavior of Algorithm 4 is summarized in the lemma below. For the lemma statement, assume $n \geq 2 w^{2}, Z$ is a full rank linear matrix over $\mathbb{L}$, and $F_{t}$ is a ( $w, t, n$ )matrix product $R_{1} \cdot R_{2} \ldots R_{t}$ over $\mathbb{L}$, where $t \leq d$. Also, $R_{1}, R_{2}, \ldots, R_{t}$ are full rank linear matrices, and $\operatorname{det}\left(R_{1}\right), \operatorname{det}\left(R_{2}\right), \ldots, \operatorname{det}\left(R_{t}\right)$ are coprime irreducible polynomials. Further, there are matrices $C, D \in \mathrm{GL}(w, \mathbb{L})$ and $i \in[t]$ such that $Z=C \cdot R_{i} \cdot D$ or $Z^{T}=C \cdot R_{i} \cdot D$.

Lemma 3.1. Algorithm 4 takes input $Z$ and blackbox access to the $w^{2}$ entries of $F_{t}$, and with probability $1-(w d n)^{-\Omega(1)}$ does this: If $Z=C \cdot R_{t} \cdot D$ then it outputs a $\tilde{D}=a D$ for an $a \in \mathbb{L}^{\times}$, and for all other cases - $Z=C \cdot R_{i} \cdot D$ or $Z^{T}=C \cdot R_{i} \cdot D$ for $i \in[t-1]$, or $Z^{T}=C \cdot R_{t} \cdot D$ - it outputs 'Failed'.

Algorithm 4 and the proof of Lemma 3.1 are presented in Section 3.2. We analyze Algorithm 3 below by tracing its steps:
 at the start of the loop the algorithm ensures $F_{t}$ is a $(w, t, n)$-matrix product $R_{1} \cdot R_{2} \ldots R_{t}{ }^{25}$ over $\mathbb{L}$, where $R_{1}, R_{2}, \ldots, R_{t}$ are full rank linear matrices and $\operatorname{det}\left(R_{1}\right), \operatorname{det}\left(R_{2}\right), \ldots, \operatorname{det}\left(R_{t}\right)$ are coprime irreducible polynomials. Further, there is a permutation $\sigma_{t}$ of $[t]$, and for every $i \in[t]$ there are matrices $C_{i}, D_{i} \in \mathrm{GL}(w, \mathbb{L})$ such that either $Z_{\sigma_{t}(i)}=C_{i} \cdot R_{i} \cdot D_{i}$ or $Z_{\sigma_{t}(i)}^{T}=C_{i} \cdot R_{i} \cdot D_{i}$. In the loop,

[^10]```
Algorithm 3 Rearrangement of the matrices
    INPUT: Blackbox access to \(F\), and \(w \times w\) full rank linear matrices \(Z_{1}, Z_{2}, \ldots, Z_{d}\) over \(\mathbb{L}\).
    OUTPUT: Linear matrices \(Y_{1}, Y_{2}, \ldots, Y_{d}\) over \(\mathbb{L}\) such that \(F=Y_{1} \cdot Y_{2} \cdots Y_{d}\).
    Set \(t=d, k=1\), and \(F_{d}=F\).
    while \(t>1\) do
        while \(k \leq t\) do
            Call Algorithm 4 on inputs \(F_{t}\) and \(Z_{k}\).
            if Algorithm 4 outputs \(\tilde{D}\) then
            Rename \(Z_{k}\) as \(Z_{t}\) and \(Z_{t}\) as \(Z_{k}\), and set \(\tilde{D}_{t}=\tilde{D} . \quad / * \sigma\) is determined implicitly. */
            Set \(M_{t}=Z_{t}\) and \(F_{t-1}=F_{t} \cdot \tilde{D}_{t} \cdot M_{t}^{-1}\).
            Set \(k=1\) and \(t=t-1\).
            Exit the inner loop.
            end if
            Call Algorithm 4 on inputs \(F_{t}\) and \(Z_{k}^{T}\).
            if Algorithm 4 outputs a \(\tilde{D}\) then
            Rename \(Z_{k}\) as \(Z_{t}\) and \(Z_{t}\) as \(Z_{k}\), and set \(\tilde{D}_{t}=\tilde{D} . \quad / * \sigma\) is determined implicitly. */
            Set \(M_{t}=Z_{t}^{T}\) and \(F_{t-1}=F_{t} \cdot \tilde{D}_{t} \cdot M_{t}^{-1}\).
            Set \(k=1\) and \(t=t-1\).
            Exit the inner loop.
            end if
            Set \(k=k+1\).
        end while
        if \(k=t+1\) then
            Exit the outer loop.
        end if
    end while
    if \(t \geq 2\) then
        Output 'Rearrangement not possible'.
    else
        Set \(Y_{1}=F_{1}\), and \(Y_{t}=M_{t} \cdot \tilde{D}_{t}^{-1}\) for all \(t \in[2, d]\). Output \(Y_{1}, \ldots, Y_{d}\).
    end if
```

the algorithm determines $\sigma_{t}(t)$ and whether $Z_{\sigma_{t}(t)}=C_{t} \cdot R_{t} \cdot D_{t}$ or $Z_{\sigma_{t}(t)}^{T}=C_{t} \cdot R_{t} \cdot D_{t}$.
Steps 4-21: Inside the inner loop, the algorithm calls Algorithm 4 on inputs $F_{t}, Z_{k}$ (step 5) and $\overline{F_{t}, Z_{k}^{T} \text { (step 13) for all } k \in[t] \text {. By Lemma 3.1, only when } k=\sigma_{t}(t) \text {, Algorithm } 4 \text { returns a } \tilde{D}=a_{t} D_{t}, ~ t h e r ~}$ for some $a_{t} \in \mathbb{L}^{\times}$. The renaming of $Z_{k}$ and $Z_{t}$ (in steps 7 and 15) ensures that we have a suitable permutation $\sigma_{t-1}$ of $[t-1]$ in the next iteration of the outer loop. The setting of $M_{t}$ (in steps 8 and 16) implies that $M_{t}=C_{t} \cdot R_{t} \cdot D_{t}$. Hence,

$$
F_{t-1}=F_{t} \cdot \tilde{D}_{t} \cdot M_{t}^{-1}=\left(R_{1} \cdot R_{2} \ldots R_{t-1}\right) \cdot\left(a_{t} C_{t}^{-1}\right)
$$

By reusing symbols and calling $R_{t-1} \cdot\left(a_{t} C_{t}^{-1}\right)$ as $R_{t-1}$, and $a_{t}^{-1} C_{t} \cdot D_{t-1}$ as $D_{t-1}$, we observe that the setup at step 2 is maintained in the next iteration of the outer loop.

Step 32: As $F_{t-1}=F_{t} \cdot \tilde{D}_{t} \cdot M_{t}^{-1}$ at every iteration of the outer loop, setting $Y_{t}=M_{t} \cdot \tilde{D}_{t}^{-1}$ implies $F_{t-1}=F_{t} \cdot Y_{t}^{-1}$ for every $t \in[d, 2]$. Therefore, $F=F_{d}=Y_{1} \cdots Y_{d}$.

### 3.2 Determining the last matrix: Proof of Lemma 3.1

We give an overview of the proof by first assuming that $Z$ is the 'last' matrix in the product $F_{t}$. The correctness of the idea is then made precise by tracing the steps of Algorithm 4.

Overview: Suppose $Z=C \cdot R_{t} \cdot D$, where $C, D \in G L(w, \mathbb{L})$. As $Z$ is a full rank linear matrix, we can assume the entries of $Z$ are distinct variables, by applying an invertible linear transformation. For any polynomial $h \in \mathbb{L}[\mathbf{x}], h \bmod \operatorname{det}(Z)$ can be identified with an element of $\mathbb{L}(\mathbf{x})^{26}$. Let $Z^{\prime}, F_{t}^{\prime} \in \mathbb{L}(\mathbf{x})^{w \times w}$ be obtained by reducing the entries of $Z$ and $F_{t}$, respectively, modulo $\operatorname{det}(Z)$. The coprimality of the determinants of $R_{1}, \ldots, R_{t}$ and their full rank nature imply,

$$
D \cdot \operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(Z^{\prime}\right)=\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(F_{t}^{\prime}\right),
$$

and these two kernels have dimensions one. A basis of $\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(Z^{\prime}\right)$ can be easily derived as $Z$ is known explicitly. However, we only have blackbox access to $F_{t}^{\prime}$. To leverage the above relation, we compute bases of $\operatorname{Kernel}_{\mathbb{L}}\left(F_{t}^{\prime}(\mathbf{a})\right)$ and $\operatorname{Kernel}_{\mathbb{L}}\left(Z^{\prime}(\mathbf{a})\right)$ for several random $\mathbf{a} \in_{r} \mathbb{F}^{n}$, and form two matrices $U, V \in \mathrm{GL}(w, \mathbb{L})$ from these bases so that $D$ equals $U \cdot V^{-1}$ (up to scaling by elements in $\mathbb{L}^{\times}$). Hereafter, Kernel $_{\mathbb{L}}$ will be denoted as Ker in the analysis of Algorithm 4.

Applying an invertible linear map (Step 2): The invertible linear transformation lets us assume that $\bar{Z}=\left(z_{l k}\right)_{l, k \in[w]}$, where $z_{l k}$ 's are distinct variables in $\mathbf{x}$.

Reducing $Z$ and $F_{t}$ modulo $\operatorname{det}(Z)$ (Step 5): The reduction of the entries of $Z$ and the blackbox entries of $F_{t} \operatorname{modulo} \operatorname{det}(Z)$ is achieved by the substitution,

$$
z_{11}=-\frac{\sum_{k=2}^{w} z_{1 k} \cdot N_{1 k}}{N_{11}}
$$

[^11]```
Algorithm 4 Determining the last matrix
    INPUT: Blackbox access to a \((w, t, n)\)-matrix product \(F_{t}\) and a full rank linear matrix \(Z\) over \(\mathbb{L}\).
    OUTPUT: A matrix \(\tilde{D} \in \mathrm{GL}(w, \mathbb{L})\), if \(Z\) is the 'last' matrix of the product \(F_{t}\).
    /* Applying an invertible linear map */
    . Let the first \(w^{2}\) variables in \(\mathbf{x}\) be \(\mathbf{z}=\left\{z_{l k}\right\}_{l, k \in[w]}\). Compute an invertible linear map \(A\) that
    maps the affine forms in \(Z\) to distinct \(\mathbf{z}\) variables, and apply \(A\) to the \(w^{2}\) blackbox entries of \(F_{t}\).
    Reusing symbols, \(Z=\left(z_{l k}\right)_{l, k \in[w]}\) and \(F_{t}\) is the matrix product after the transformation.
    3.
    4. \(/{ }^{*}\) Reducing \(Z\) and \(F_{t}\) modulo \(\operatorname{det}(Z) * /\)
    Let \(N_{l k}\) be the \((l, k)\)-th cofactor of \(Z\), for \(l, k \in[w]\). Substitute \(z_{11}=\frac{-\sum_{k=2}^{w} z_{1 k} N_{1 k}}{N_{11}}\) in \(Z\) and in the
    blackbox for \(F_{t}\). Call the matrices \(Z^{\prime}\) and \(F_{t}^{\prime}\) respectively after the substitution.
6.
7. /* Computing the kernels */
    for \(k=1\) to \(w+1\) do
    Choose \(\mathbf{a}_{k}, \mathbf{b}_{k} \in_{r} \mathbb{F}^{n}\). Compute bases of \(\operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{a}_{k}\right)\right), \operatorname{Ker}\left(Z^{\prime}\left(\mathbf{a}_{k}\right)\right), \operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{b}_{k}\right)\right), \operatorname{Ker}\left(Z^{\prime}\left(\mathbf{b}_{k}\right)\right)\).
        Pick non-zero \(\mathbf{u}_{k} \in \operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{a}_{k}\right)\right)\), \(\mathbf{v}_{k} \in \operatorname{Ker}\left(Z^{\prime}\left(\mathbf{a}_{k}\right)\right)\), \(\mathbf{w}_{k} \in \operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{b}_{k}\right)\right)\), \(\mathbf{s}_{k} \in \operatorname{Ker}\left(Z^{\prime}\left(\mathbf{b}_{k}\right)\right)\).
        If the computation fails (i.e., \(N_{11}\left(\mathbf{a}_{k}\right)=0\) or \(N_{11}\left(\mathbf{b}_{k}\right)=0\) ), or any of the kernels is empty,
        output 'Failed'.
    end for
11.
2. /* Extracting D from the kernels */
13. Compute \(\alpha_{k}, \beta_{k}, \gamma_{k}, \delta_{k} \in \mathbb{L}\) for \(k \in[w]\) such that \(\mathbf{u}_{w+1}=\sum_{k=1}^{w} \alpha_{k} \mathbf{u}_{k}, \mathbf{v}_{w+1}=\sum_{k=1}^{w} \beta_{k} \mathbf{v}_{k}, \mathbf{w}_{w+1}=\)
    \(\sum_{k=1}^{w} \gamma_{k} \mathbf{w}_{k}\) and \(\mathbf{s}_{w+1}=\sum_{k=1}^{w} \delta_{k} \mathbf{s}_{k}\). If the computation fails, or any of \(\alpha_{k}, \beta_{k}, \gamma_{k}, \delta_{k}\) is zero for
    some \(k \in[w]\), output 'Failed'.
14.
15. Set \(U, V, W, S \in \mathbb{L}^{w \times w}\) such that the \(k\)-th column of \(U, V, W, S\) are \(\frac{\alpha_{k} \cdot \mathbf{u}_{k}}{\beta_{k}}, \mathbf{v}_{k}, \frac{\gamma_{k} \cdot \mathbf{w}_{k}}{\delta_{k}}, \mathbf{s}_{k}\) respectively.
    If any of \(U, V, W, S \notin \mathrm{GL}(w, \mathbb{L})\), output 'Failed'.
16.
7. if \(U V^{-1} S W^{-1}\) is a scalar matrix then
        Set \(\tilde{D}=U \cdot V^{-1}\) and output \(\tilde{D}\).
    else
        Output 'Failed'. /* The check fails w.h.p if Z is not the 'last' matrix */
    end if
```

After the substitution, the matrices become $Z^{\prime}$ and $F_{t}^{\prime}=R_{1}^{\prime} \cdot R_{2}^{\prime} \ldots R_{t}^{\prime}$ respectively. As there is an $i \in[t]$ and $C, D \in G L(w, \mathbb{L})$ such that either $Z=C \cdot R_{i} \cdot D$ or $Z^{T}=C \cdot R_{i} \cdot D$, we have either $Z^{\prime}=C \cdot R_{i}^{\prime} \cdot D$ or $\left(Z^{\prime}\right)^{T}=C \cdot R_{i}^{\prime} \cdot D$ and hence $\operatorname{det}\left(Z^{\prime}\right)=\operatorname{det}\left(R_{i}^{\prime}\right)=\operatorname{det}\left(F_{t}^{\prime}\right)=0$.
Observation 3.1. 1. $\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(Z^{\prime}\right)=\operatorname{span}_{\mathbb{L}(\mathbf{x})}\left\{\left(N_{11} N_{12} \ldots N_{1 w}\right)^{T}\right\}$,
2. $\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(\left(Z^{\prime}\right)^{T}\right)=\operatorname{span}_{\mathbb{L}(\mathbf{x})}\left\{\left(N_{11} N_{21} \ldots N_{w 1}\right)^{T}\right\}$.

Hence, $\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(Z^{\prime}\right)$ has dimension one, and the observation below implies $\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(F_{t}^{\prime}\right)$ is also one dimensional. The proof follows from the coprimality of $\operatorname{det}\left(R_{1}\right), \operatorname{det}\left(R_{2}\right), \ldots, \operatorname{det}\left(R_{t}\right)$.

Observation 3.2. For all $j \in[t]$ and $j \neq i, \operatorname{det}\left(R_{j}^{\prime}\right) \neq 0$, and so the dimension of $\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(F_{t}^{\prime}\right)$ is one.

Computing the kernels (Steps 8-10): The following observation shows that the algorithm does not fail at step 9 with high probability. The proof is immediate from the above two observations and an application of the Schwartz-Zippel lemma.

Observation 3.3. Let $\mathbf{a}_{k}, \mathbf{b}_{k} \in_{r} \mathbb{F}^{n}$ for $k \in[w+1]$. Then, for every $k \in[w+1]$, and $\mathbf{a}=\mathbf{a}_{k}$ or $\mathbf{b}_{k}$,

1. $\operatorname{Ker}\left(Z^{\prime}(\mathbf{a})\right)=\operatorname{span}_{\mathbb{L}}\left\{\left(N_{11}(\mathbf{a}) N_{12}(\mathbf{a}) \ldots N_{1 w}(\mathbf{a})\right)^{T}\right\}$,
2. $\operatorname{Ker}\left(\left(Z^{\prime}(\mathbf{a})\right)^{T}\right)=\operatorname{span}_{\mathbb{L}}\left\{\left(N_{11}(\mathbf{a}) N_{21}(\mathbf{a}) \ldots N_{v 1}(\mathbf{a})\right)^{T}\right\}$,
and $\operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{a}_{k}\right)\right), \operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{b}_{k}\right)\right)$ are one dimensional subspaces of $\mathbb{L}^{w}$, with probability $1-(w d n)^{-\Omega(1)}$.

Extracting D from the kernels (Steps 13-21): We analyse these steps for three separate cases. The analysis shows that if $Z$ is the 'last' matrix then the algorithm succeeds with high probability, otherwise the test at step 17 fails with high probability.

Case a $\left[Z=C \cdot R_{t} \cdot D\right]$ : From Observation 3.2, $\operatorname{det}\left(R_{j}^{\prime}\left(\mathbf{a}_{k}\right)\right)$ and $\operatorname{det}\left(R_{j}^{\prime}\left(\mathbf{b}_{k}\right)\right)$ are nonzero with high probability, for all $j \in[t-1]$ and $k \in[w+1]$. Assuming this, the following holds for all $k \in[w+1]$ :

$$
\begin{align*}
D \cdot \operatorname{Ker}\left(Z^{\prime}\left(\mathbf{a}_{k}\right)\right) & =\operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{a}_{k}\right)\right), \\
D \cdot \operatorname{Ker}\left(Z^{\prime}\left(\mathbf{b}_{k}\right)\right) & =\operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{b}_{k}\right)\right) \tag{2}
\end{align*}
$$

Hence, at step 9, there are $\lambda_{k}, \rho_{k} \in \mathbb{L}^{\times}$such that

$$
D \cdot \mathbf{v}_{k}=\lambda_{k} \mathbf{u}_{k}, \quad D \cdot \mathbf{s}_{k}=\rho_{k} \mathbf{w}_{k} \quad \text { for } k \in[w+1] .
$$

Step 13 also succeeds with high probability due to the following claim (proof in Appendix A).
Claim 3.1. With probability $1-(w d n)^{-\Omega(1)}$, any subset of $w$ vectors in any of the sets $\left\{\mathbf{u}_{1}, \mathbf{u}_{2}, \ldots, \mathbf{u}_{w+1}\right\}$, $\left\{\mathbf{v}_{1}, \mathbf{v}_{2}, \ldots, \mathbf{v}_{w+1}\right\},\left\{\mathbf{w}_{1}, \mathbf{w}_{2}, \ldots, \mathbf{w}_{w+1}\right\}$, or $\left\{\mathbf{s}_{1}, \mathbf{s}_{2}, \ldots, \mathbf{s}_{w+1}\right\}$ are $\mathbb{L}$-linearly independent.

At this step, $\mathbf{v}_{w+1}=\sum_{k=1}^{w} \beta_{k} \mathbf{v}_{k}$ and $\mathbf{s}_{w+1}=\sum_{k=1}^{w} \delta_{k} \mathbf{s}_{k}$, and so by applying $D$ on both sides,

$$
\lambda_{w+1} \mathbf{u}_{w+1}=\sum_{k=1}^{w} \beta_{k} \lambda_{k} \mathbf{u}_{k}, \quad \rho_{w+1} \mathbf{w}_{w+1}=\sum_{k=1}^{w} \delta_{k} \rho_{k} \mathbf{w}_{k} .
$$

Also, $\mathbf{u}_{w+1}=\sum_{k=1}^{w} \alpha_{k} \mathbf{u}_{k}$ and $\mathbf{w}_{w+1}=\sum_{k=1}^{w} \gamma_{k} \mathbf{w}_{k}$. By Claim 3.1, none of the $\alpha_{k}, \beta_{k}, \gamma_{k}, \delta_{k}$ is zero and

$$
\frac{\lambda_{k}}{\lambda_{w+1}}=\frac{\alpha_{k}}{\beta_{k}}, \quad \frac{\rho_{k}}{\rho_{w+1}}=\frac{\gamma_{k}}{\delta_{k}}, \quad \text { for all } k \in[w] .
$$

From the construction of $U, V, W$ and $S$ at step 15,

$$
D \cdot V=\lambda_{w+1} U, \quad D \cdot S=\rho_{w+1} W
$$

and $U, V, W, S \in \mathrm{GL}(w, \mathbb{L})$ (by Claim 3.1). Therefore, $U V^{-1} S W^{-1}$ is a scalar matrix.
Case $\mathbf{b}\left[Z^{T}=C \cdot R_{t} \cdot D\right]$ : In this case, the check at step 17 fails with high probability. Suppose the algorithm passes steps 13 and 15 , and reaches step 17 . We show that $U V^{-1} S W^{-1}$ being a scalar matrix implies an event $\mathcal{E}$ that happens with a low probability. The event $\mathcal{E}$ can be derived as follows:

Let $M \stackrel{\text { def }}{=} U \cdot V^{-1}$, and $c \in \mathbb{L}^{\times}$such that $M=c W \cdot S^{-1}$. Assuming the invertibility of $R_{j}^{\prime}\left(\mathbf{a}_{k}\right)$ and $R_{j}^{\prime}\left(\mathbf{b}_{k}\right)$ for $j \in[t-1]$ (Observation 3.2), and as in Equation 2, the following holds for all $k \in[w+1]$.

$$
\begin{aligned}
D \cdot \operatorname{Ker}\left(\left(Z^{\prime}\left(\mathbf{a}_{k}\right)\right)^{T}\right) & =\operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{a}_{k}\right)\right), \\
D \cdot \operatorname{Ker}\left(\left(Z^{\prime}\left(\mathbf{b}_{k}\right)\right)^{T}\right) & =\operatorname{Ker}\left(F_{t}^{\prime}\left(\mathbf{b}_{k}\right)\right) .
\end{aligned}
$$

By Observation 3.3, we can assume the above four kernels are one-dimensional. Hence, at step 9 there are $\mathbf{p}_{k} \in \operatorname{Ker}\left(\left(Z^{\prime}\left(\mathbf{a}_{k}\right)\right)^{T}\right)$ and $\mathbf{q}_{k} \in \operatorname{Ker}\left(\left(Z^{\prime}\left(\mathbf{b}_{k}\right)\right)^{T}\right)$ satisfying $D \cdot \mathbf{p}_{k}=\mathbf{u}_{k}$ and $D \cdot \mathbf{q}_{k}=\mathbf{w}_{k}$, for every $k \in[w+1]$. Consider the $w \times w$ matrices $P$ and $Q$ such that the $k$-th column of these matrices are $\frac{\alpha_{k}}{\beta_{k}} \mathbf{p}_{k}$ and $\frac{\gamma_{k}}{\delta_{k}} \mathbf{q}_{k}$ respectively, where $\alpha_{k}, \beta_{k}, \gamma_{k}, \delta_{k}$ are the constants computed at step 13. Clearly, $D \cdot P=U$ and $D \cdot Q=W$, where $U, W$ are the matrices computed at step 15 .

As $M=c W \cdot S^{-1}$ (by assumption), we have $D^{-1} M S=c D^{-1} W=c Q$. Hence, for $k \in[w]$,

$$
D^{-1} M \cdot \mathbf{s}_{k}=\frac{c \gamma_{k}}{\delta_{k}} \mathbf{q}_{k}
$$

At step $13, \mathbf{w}_{w+1}=\sum_{k=1}^{w} \gamma_{k} \mathbf{w}_{k}$ and $\mathbf{s}_{w+1}=\sum_{k=1}^{w} \delta_{k} \mathbf{s}_{k}$. Multiplying $D^{-1}$ on both sides and $D^{-1} M$ on both sides of these two equations respectively,

$$
\begin{gather*}
\mathbf{q}_{w+1}=\sum_{k=1}^{w} \gamma_{k} \mathbf{q}_{k}, \quad \text { and } \quad D^{-1} M \cdot \mathbf{s}_{w+1}=\sum_{k=1}^{w} c \gamma_{k} \mathbf{q}_{k} . \\
\Rightarrow \quad D^{-1} M \cdot \mathbf{s}_{w+1}=c \mathbf{q}_{w+1} . \tag{3}
\end{gather*}
$$

From Observation 3.3, there are $\lambda_{1}, \lambda_{2} \in \mathbb{L}^{\times}$such that

$$
\left.\begin{array}{rl}
\mathbf{s}_{w+1} & =\lambda_{1} \cdot\left(N_{11}\left(\mathbf{b}_{w+1}\right)\right. \\
N_{12}\left(\mathbf{b}_{w+1}\right) \ldots & \ldots \\
\left.\mathbf{N}_{1 w}\left(\mathbf{b}_{w+1}\right)\right)^{T}, \\
\mathbf{q}_{w+1} & =\lambda_{2} \cdot\left(N_{11}\left(\mathbf{b}_{w+1}\right)\right.
\end{array} N_{21}\left(\mathbf{b}_{w+1}\right) \ldots N_{w 1}\left(\mathbf{b}_{w+1}\right)\right)^{T} .
$$

Let $D^{-1} M=\left(m_{l k}\right)_{l, k \in[w]}$. Using the above values of $\mathbf{s}_{w+1}$ and $\mathbf{q}_{w+1}$ in Equation 3 and restricting to the first two entries of the resulting column vectors, we have

$$
\lambda_{1}\left(\sum_{k=1}^{w} m_{1 k} N_{1 k}\left(\mathbf{b}_{w+1}\right)\right)=c \lambda_{2} N_{11}\left(\mathbf{b}_{w+1}\right), \quad \lambda_{1}\left(\sum_{k=1}^{w} m_{2 k} N_{1 k}\left(\mathbf{b}_{w+1}\right)\right)=c \lambda_{2} N_{21}\left(\mathbf{b}_{w+1}\right) .
$$

Thus we get the following relation,

$$
N_{21}\left(\mathbf{b}_{w+1}\right)\left(\sum_{k=1}^{w} m_{1 k} N_{1 k}\left(\mathbf{b}_{w+1}\right)\right)=N_{11}\left(\mathbf{b}_{w+1}\right)\left(\sum_{k=1}^{w} m_{2 k} N_{1 k}\left(\mathbf{b}_{w+1}\right)\right) .
$$

Event $\mathcal{E}$ is defined by the above equality, i.e. we say $\mathcal{E}$ has happened whenever the above equality holds. Now observe that $D^{-1} M$ is independent ${ }^{27}$ of the random bits used to choose $\mathbf{b}_{w+1}$. Hence, it is sufficient to show that the above equality happens with low probability over the randomness of $\mathbf{b}_{w+1}$, for any arbitrarily fixed $m_{11}, \ldots, m_{1 w}$ and $m_{21}, \ldots, m_{2 w}$ from $\mathbb{L}$. Moreover, as $D^{-1} M$ is invertible, we can assume - not all in $\left\{m_{11}, \ldots, m_{1 w}\right\}$ or $\left\{m_{21}, \ldots, m_{2 w}\right\}$ are zero. The following observation and Schwartz-Zippel lemma complete the proof in this case.

Observation 3.4. $N_{21}(\mathbf{z})\left(\sum_{k=1}^{w} m_{1 k} \cdot N_{1 k}(\mathbf{z})\right) \neq N_{11}(\mathbf{z})\left(\sum_{k=1}^{w} m_{2 k} \cdot N_{1 k}(\mathbf{z})\right)$ as polynomials in $\mathbb{F}[\mathbf{z}]$.
Proof. Suppose the two sides are equal. As $N_{21}(\mathbf{z})$ and $N_{11}(\mathbf{z})$ are irreducible and coprime polynomials, $N_{21}(\mathbf{z})$ must divide $\sum_{k=1}^{w} m_{2 k} \cdot N_{1 k}(\mathbf{z})$. But the two polynomials have the same degree and they are monomial disjoint, thereby giving us a contradiction.

Case c $\left[Z=C \cdot R_{i} \cdot D\right.$ or $Z^{T}=C \cdot R_{i} \cdot D$ for some $\left.i \in[t-1]\right]$ : Assume $Z=C \cdot R_{i} \cdot D$ for some $i \in[t-1]$. The case $Z^{T}=C \cdot R_{i} \cdot D$ can be argued similarly. Similar to Case b , we show that if the algorithm passes steps 13 and 15 , and reaches step 17 then $U V^{-1} S W^{-1}$ being a scalar matrix implies an event $\mathcal{E}$ that happens with very low probability. Hence, the check at step 17 fails with high probability. The event $\mathcal{E}$ can be derived as follows:

Let $M \stackrel{\text { def }}{=} U \cdot V^{-1}$, and $c \in \mathbb{L}^{\times}$be such that $M=c \cdot W S^{-1}$. From the construction of $W$ and $S$,

$$
\frac{c \gamma_{k}}{\delta_{k}} \mathbf{w}_{k}=M \cdot \mathbf{s}_{k}, \quad \text { for all } k \in[w],
$$

where $\gamma_{k}, \delta_{k}$ are as computed at step 13. Since $\mathbf{w}_{w+1}=\sum_{k=1}^{w} \gamma_{k} \mathbf{w}_{k}$ and $\mathbf{s}_{w+1}=\sum_{k=1}^{w} \delta_{k} \cdot \mathbf{s}_{k}$,

$$
c \cdot \mathbf{w}_{w+1}=M \cdot \mathbf{s}_{w+1} .
$$

Let $H \stackrel{\text { def }}{=} D^{-1} \cdot R_{i+1}^{\prime} \ldots R_{t}^{\prime}$. From Observation 3.2, the following holds,

$$
H^{-1} \cdot \operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(Z^{\prime}\right)=\operatorname{Kernel}_{\mathbb{L}(\mathbf{x})}\left(F_{t}^{\prime}\right)
$$

Let $\mathbf{n}=\left(N_{11}\left(\mathbf{b}_{w+1}\right) \quad N_{12}\left(\mathbf{b}_{w+1}\right) \ldots N_{1 w}\left(\mathbf{b}_{w+1}\right)\right)^{T}$. From Observation 3.3, and as $H\left(\mathbf{b}_{w+1}\right)$ is invertible with high probability over the random choice of $\mathbf{b}_{w+1}$, there are $\lambda_{1}, \lambda_{2} \in \mathbb{L}^{\times}$such that

$$
\begin{aligned}
\mathbf{w}_{w+1} & =\lambda_{1} H^{-1}\left(\mathbf{b}_{w+1}\right) \cdot \mathbf{n} \\
\mathbf{s}_{w+1} & =\lambda_{2} \mathbf{n} .
\end{aligned}
$$

Substituting the above values of $\mathbf{w}_{w+1}$ and $\mathbf{s}_{w+1}$ in $c \cdot \mathbf{w}_{w+1}=M \cdot \mathbf{s}_{w+1}$, we have

$$
c \lambda_{1} H^{-1}\left(\mathbf{b}_{w+1}\right) \cdot \mathbf{n}=\lambda_{2} M \cdot \mathbf{n}, \quad \Rightarrow \quad c \lambda_{1} \mathbf{n}=\lambda_{2} H\left(\mathbf{b}_{w+1}\right) \cdot M \cdot \mathbf{n} .
$$

[^12]Let $H \cdot M=\left(h_{l k}\right)_{l, k \in[w]}$. Restricting to the first two entries of the vectors in the above equality, we have

$$
\begin{aligned}
& c \lambda_{1} N_{11}\left(\mathbf{b}_{w+1}\right)=\lambda_{2}\left(\sum_{k=1}^{w} h_{1 k}\left(\mathbf{b}_{w+1}\right) \cdot N_{1 k}\left(\mathbf{b}_{w+1}\right)\right), \\
& c \lambda_{1} N_{12}\left(\mathbf{b}_{w+1}\right)=\lambda_{2}\left(\sum_{k=1}^{w} h_{2 k}\left(\mathbf{b}_{w+1}\right) \cdot N_{1 k}\left(\mathbf{b}_{w+1}\right)\right) .
\end{aligned}
$$

Hence, we get the following relation

$$
\begin{equation*}
N_{11}\left(\mathbf{b}_{w+1}\right) \cdot\left(\sum_{k=1}^{w} h_{2 k}\left(\mathbf{b}_{w+1}\right) \cdot N_{1 k}\left(\mathbf{b}_{w+1}\right)\right)=N_{12}\left(\mathbf{b}_{w+1}\right) \cdot\left(\sum_{k=1}^{w} h_{1 k}\left(\mathbf{b}_{w+1}\right) \cdot N_{1 k}\left(\mathbf{b}_{w+1}\right)\right) . \tag{4}
\end{equation*}
$$

Event $\mathcal{E}$ is defined by the above equality, that is $\mathcal{E}$ happens if the above equality is satisfied. Observe that the entries of the matrix product $H \cdot M=\left(h_{l k}\right)_{l, k \in[w]}$ are rational functions in $\mathbf{x}$ variables and are independent of the random bits used to choose $\mathbf{b}_{w+1}$. We show next the probability that the above equality holds is low over the randomness of $\mathbf{b}_{w+1}$.

The only implications of the average-case nature of $F_{t}$ that we have used in the proofs so far are: every $R_{i}$ is full rank and $\operatorname{det}\left(R_{1}\right), \ldots, \operatorname{det}\left(R_{t}\right)$ are mutually coprime with high probability. However, these two properties are not sufficient to ensure the uniqueness of the last matrix in the product (as mentioned in a remark after Theorem 1). In the following claim, we use one more effect of $F_{t}$ being a random matrix product which ensures the desired uniqueness of the last matrix.

Claim 3.2. If $E=Q_{1} \cdots Q_{\ell}$ is a random ( $w, \ell, m$ )-matrix product over $\mathbb{F}$, where $w^{2}+1 \leq m \leq n$ and $\ell \leq d$, then the entries of $E$ are $\mathbb{F}$-linearly independent with probability $1-(w d n)^{-\Omega(1)}$.

If the entries of $E$ are $\mathbb{F}$-linearly independent then they are also $\mathbb{L}$-linearly independent. We conclude the proof of Case c using the above claim (proof given in Appendix A).

Observation 3.5. Let $n \geq 2 w^{2}$. Then all the entries of $H \cdot M$ are nonzero polynomials after setting the variables in $\mathbf{z}_{1} \stackrel{\text { def }}{=}\left\{z_{11}, z_{21}, z_{31}, \ldots, z_{w 1}\right\}$ to zero, with probability $1-(w d n)^{-\Omega(1)}$.
Proof. $H \cdot M=D^{-1} \cdot R_{i+1}^{\prime} \ldots R_{t}^{\prime} \cdot M=\left(h_{l k}\right)_{l, k \in[w]}$. Recalling the substitution $z_{11}=\frac{-\sum_{k=1}^{w} z_{1 k} N_{1 k}}{N_{11}}$ at step 5, we observe that the rational function $h_{l k}$ becomes a polynomial under the setting $z_{11}=$ $z_{21}=\ldots=z_{w 1}=0^{28}$. Let $Q_{j}=\left(R_{j}\right)_{\mathbf{z}_{1}=0}$. By observing $\left(R_{j}\right)_{\mathbf{z}_{1}=0}=\left(R_{j}^{\prime}\right)_{\mathbf{z}_{1}=0}$, it follows that $(H \cdot M)_{\mathbf{z}_{1}=0}=D^{-1} \cdot Q_{i+1} \ldots Q_{t} \cdot M$. Moreover, $Q_{i+1} \cdot Q_{i+2} \ldots Q_{t}$ is a random $(w, t-i, n-w)-$ matrix product. By Claim 3.2, the entries of $Q_{i+1} \ldots Q_{t}$ are $\mathbb{L}$-linearly independent with high probability. Hence, none of the entries of $D^{-1} \cdot Q_{i+1} \ldots Q_{t} \cdot M$ is zero as $D, M \in \mathrm{GL}(\mathbb{L}, w)$.

Observation 3.6. $N_{11}(\mathbf{x}) \cdot\left(\sum_{k=1}^{w} h_{2 k}(\mathbf{x}) N_{1 k}(\mathbf{x})\right) \neq N_{12}(\mathbf{x}) \cdot\left(\sum_{k=1}^{w} h_{1 k}(\mathbf{x}) N_{1 k}(\mathbf{x})\right)$ as rational functions in $\mathbb{L}(\mathbf{x})$, with probability $1-(\text { wdn })^{-\Omega(1)}$.
Proof. Suppose $N_{11}(\mathbf{x}) \cdot\left(\sum_{k=1}^{w} h_{2 k}(\mathbf{x}) N_{1 k}(\mathbf{x})\right)=N_{12}(\mathbf{x}) \cdot\left(\sum_{k=1}^{w} h_{1 k}(\mathbf{x}) N_{1 k}(\mathbf{x})\right)$. By substituting $\mathbf{z}_{1}=$ 0 in the equation, the R.H.S becomes zero whereas the L.H.S reduces to $N_{11}^{2} \cdot\left(h_{21}\right)_{\mathbf{z}_{1}=0} \neq 0$ with high probability (from Observation 3.5).

[^13]Noting that the degrees of the numerator and the denominator of $h_{l k}$ are upper bounded by $w d$, we conclude that the equality in Equation 4 happens with a low probability over the randomness of $\mathbf{b}_{w+1}$.

## 4 Average-case ABP reconstruction: Proof of Theorem 2

The algorithm for average-case ABP reconstruction is presented in Algorithm 2, Section 1.3.2. The algorithm uses Algorithm 5 and Algorithm 6 during its execution - we present and analyze these two algorithms in the following subsections.

### 4.1 Computing the corner spaces

Let $f$ be the polynomial computed by a random $(w, d, n)$-ABP $X_{1} \cdot X_{2} \ldots X_{d}$ over $\mathbb{F}$, where $n \geq 4 w^{2}$.
Lemma 4.1. With probability $1-(w d n)^{-\Omega(1)}$ over the randomness of $f$, the following holds: Let $\mathbb{K} \supseteq \mathbb{F}$ be any field and $f=0 \bmod \left\langle l_{1}, \ldots, l_{k}\right\rangle$, where $l_{i}$ 's are linear forms in $\mathbb{K}[\mathbf{x}]$. Then $k \geq w$ and for $k=w$, the space $\operatorname{span}_{\mathbb{K}}\left\{l_{1}, \ldots, l_{w}\right\}$ equals the $\mathbb{K}$-span of either the linear forms in $X_{1}$ or the linear forms in $X_{d}$.

The above uniqueness of the corner spaces, $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ (defined in Section 1.3.2), helps compute them in Algorithm 5. The proof of the lemma is given at the end of this subsection.

Canonical bases of $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$ : For a set of variables $\mathbf{y} \subseteq \mathbf{x}$ and a linear form $g$ in $\mathbb{F}[\mathbf{x}]$, define $g(\mathbf{y}) \stackrel{\text { def }}{=} g_{\mathbf{x} \backslash \mathbf{y}=0}$. We say $g(\mathbf{y})$ is the linear form $g$ projected to the $\mathbf{y}$ variables. Let $x_{1}, \ldots, x_{w}$ and $v$ be a designated set of $w+1$ variables in $\mathbf{x}$, and $\mathbf{u}=\mathbf{x} \backslash\left\{x_{1}, \ldots, x_{w}, v\right\}$. With $n \geq 4 w^{2}$, a random $(w, d, n)$-ABP $X_{1} \cdot X_{2} \ldots X_{d}$ satisfies the following condition with probability $1-(w d n)^{-\Omega(1)}$ :
(*a) The linear forms in $X_{1}$ (similarly, $X_{d}$ ) projected to $x_{1}, \ldots, x_{w}$ are $\mathbb{F}$-linearly independent. If the above condition is satisfied then there is a $C \in G L(w, \mathbb{F})$ such that the linear forms in $X_{1} \cdot C$ are of the kind:

$$
\begin{equation*}
x_{i}-\alpha_{i} v-g_{i}(\mathbf{u}), \quad \text { for } i \in[w], \tag{5}
\end{equation*}
$$

where each $\alpha_{i} \in \mathbb{F}$ and $g_{i}$ is a linear form in $\mathbb{F}[\mathbf{u}]$. Thus, we can assume without loss of generality, the linear forms in $X_{1}$ are of the above kind. Similarly, the linear forms in $X_{d}$ are also of the kind:

$$
\begin{equation*}
x_{i}-\beta_{i} v-h_{i}(\mathbf{u}), \quad \text { for } i \in[w], \tag{6}
\end{equation*}
$$

where each $\beta_{i} \in \mathbb{F}$ and $h_{i}$ is a linear form in $\mathbb{F}[\mathbf{u}]$. Moreover, with probability $1-(w d n)^{-\Omega(1)}$ over the randomness of the ABP, the following condition is satisfied:
$\left.\mathbf{(}^{*} \mathbf{b}\right) \alpha_{1}, \ldots, \alpha_{w}$ and $\beta_{1}, \ldots, \beta_{w}$ are distinct elements in $\mathbb{F}$.
The task at hand for Algorithm 5 is to solve for $\alpha_{i}, g_{i}$ and $\beta_{j}, h_{j}$, for $i, j \in[w]$, assuming that conditions ( ${ }^{*} \mathrm{a}$ ) and ( ${ }^{*} \mathrm{~b}$ ) are satisfied. The bases defined by Equations 5 and 6 are canonical for $\mathcal{X}_{1}$ and $\mathcal{X}_{d}$.

We analyze the three main steps of Algorithm 5 next:

1. Partitioning the variables (Step 2): The only thing to note here is, if $n-(w+1)$ is not divisible by $4 w^{2}-(w+1)$ then we allow the last two sets $\mathbf{u}_{m-1}$ and $\mathbf{u}_{m}$ to overlap - the algorithm can be suitably adjusted in this case.
```
Algorithm 5 Computing the corner spaces
    INPUT: Blackbox access to a \(f\) computed by a random ( \(w, d, n\) )-ABP.
    OUTPUT: Bases of the two corner spaces \(\mathcal{X}_{1}\) and \(\mathcal{X}_{d}\) modulo which \(f\) is zero.
    . /* Partitioning the variables */
    Choose \(w+1\) designated variables \(x_{1}, x_{2}, \ldots, x_{w}, v\), and let \(\mathbf{u}=\mathbf{x} \backslash\left\{x_{1}, \ldots, x_{w}, v\right\}\). Partition \(\mathbf{u}\)
    into sets \(\mathbf{u}_{1}, \mathbf{u}_{2}, \ldots, \mathbf{u}_{m}\), each of size \(4 w^{2}-(w+1)\). .
3.
    /* Reduction to solving \(m\) systems of polynomial equations */
    for \(\ell=1\) to \(m\) do
        Set \(f_{\ell}=f_{\mathbf{u} \backslash \mathbf{u}_{\ell}=0}\).
    Solve for all possible \(\left(\alpha_{1}, \ldots, \alpha_{w}, g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, g_{w}\left(\mathbf{u}_{\ell}\right)\right)\), where each \(\alpha_{i} \in \mathbb{F}\) and \(g_{i}\left(\mathbf{u}_{\ell}\right)\) is a linear
        form in \(\mathbb{F}\left[\mathbf{u}_{\ell}\right]\) such that
            \(f_{\ell}=0 \bmod \left\langle x_{1}-\alpha_{1} v-g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, x_{w}-\alpha_{w} v-g_{w}\left(\mathbf{u}_{\ell}\right)\right\rangle\).
        if Step 7 does not return exactly two solutions for \(\left(\alpha_{1}, \ldots, \alpha_{w}, g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, g_{w}\left(\mathbf{u}_{\ell}\right)\right)\) then
            Output 'Failed'.
        else
            The solutions be \(\left(\alpha_{\ell 1}, \ldots, \alpha_{\ell w}, g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, g_{w}\left(\mathbf{u}_{\ell}\right)\right)\) and \(\left(\beta_{\ell 1}, \ldots, \beta_{\ell w}, h_{1}\left(\mathbf{u}_{\ell}\right), \ldots, h_{w}\left(\mathbf{u}_{\ell}\right)\right)\).
        end if
    end for
    /* Combining the solutions */
    if \(\left|\cup_{\ell \in[m]}\left\{\left(\alpha_{\ell 1}, \ldots, \alpha_{\ell w}\right),\left(\beta_{\ell 1}, \ldots, \beta_{\ell w}\right)\right\}\right| \neq 2\) then
        Output 'Failed'.
    else
        Without loss of generality, \(\left(\alpha_{\ell 1}, \ldots, \alpha_{\ell w}\right)=\left(\alpha_{1}, \ldots, \alpha_{w}\right)\) and \(\left(\beta_{\ell 1}, \ldots, \beta_{\ell w}\right)=\left(\beta_{1}, \ldots, \beta_{w}\right)\) for
        every \(\ell \in[m]\). Set \(g_{i}(\mathbf{u})=\sum_{\ell \in[w]} g_{i}\left(\mathbf{u}_{\ell}\right)\) and \(h_{i}(\mathbf{u})=\sum_{\ell \in[w]} h_{i}\left(\mathbf{u}_{\ell}\right)\) for every \(i \in[w]\).
        Return \(\left\{x_{i}-\alpha_{i} v-g_{i}(\mathbf{u})\right\}_{i \in[w]}\) and \(\left\{x_{i}-\beta_{i} v-h_{i}(\mathbf{u})\right\}_{i \in[w]}\) as the bases of \(\mathcal{X}_{1}\) and \(\mathcal{X}_{d}\).
    end if
```

2. Reduction to solving systems of polynomial equations (Steps 5-13): At step 7, the task of comput$\operatorname{ing}\left(\alpha_{1}, \ldots, \alpha_{w}, g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, g_{w}\left(\mathbf{u}_{\ell}\right)\right)$ such that

$$
f_{\ell}=0 \bmod \left\langle x_{1}-\alpha_{1} v-g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, x_{w}-\alpha_{w} v-g_{w}\left(\mathbf{u}_{\ell}\right)\right\rangle,
$$

can be reduced to solving for all $\mathbb{F}$-rational points of a system of polynomial equations over $\mathbb{F}$ as follows: Treat $\alpha_{1}, \ldots, \alpha_{w}$ and the $4 w^{3}-w(w+1)$ coefficients of $g_{1}\left(\mathbf{u}_{\ell}\right), \ldots, g_{w}\left(\mathbf{u}_{\ell}\right)$, say $\mathbf{w}$, as formal variables. Substitute $x_{i}=\alpha_{i} v+g_{i}\left(\mathbf{u}_{\ell}\right)$ for every $i \in[w]$ in the blackbox for $f_{\ell}$ and interpolate the resulting polynomial $p$ in the variables $\alpha_{1}, \ldots, \alpha_{w}, \mathbf{w}, v, \mathbf{u}_{\ell}$ with coefficients in $\mathbb{F}$. The interpolation, which can be done in $\left(d^{w^{3}} \log q\right)^{O(1)}$ time ${ }^{29}$, gives $p$ in dense representation (i.e. as a sum of monomials). Now by treating $p$ as a polynomial in the variables $v, \mathbf{u}_{\ell}$ with coefficients in $\mathbb{F}\left(\alpha_{1}, \ldots, \alpha_{w}, \mathbf{w}\right)$, and equating these coefficients to zero, we get a system of $d^{O\left(w^{2}\right)}$ polynomial equations in $O\left(w^{3}\right)$ variables with degree of each polynomial equation bounded by $d$. By Lemma 4.1 , such a system has exactly two solutions over $\overline{\mathbb{F}}$ and moreover, these two solution points are $\mathbb{F}$-rational. Hence, by applying Lemma 2.2, we can compute the two solutions for $\left(\alpha_{1}, \ldots, \alpha_{w}, \mathbf{w}\right)$ at step 7 , in $\left(d^{w^{3}} \log q\right)^{O(1)}$ time.
3. Combining the solutions (Steps 16-21): The correctness of the steps follows from condition (*b).

## Uniqueness of the corner spaces: Proof of Lemma 4.1

As $n \geq 4 w^{2}$, a random $(w, d, n)$-ABP $X_{1} \cdots X_{d}$ satisfies the following condition with probability $1-(w d n)^{-\Omega(1)}$ :
${ }^{(* *)}$ The linear forms in $X_{1}, X_{d}$ and any three or less of the other $X_{i}$ 's are $\mathbb{F}$-linearly independent.
So, it is sufficient to prove the following restatement of the lemma.
Lemma 4.1. Suppose $f$ is computed by a ( $w, d, n$ )-ABP $X_{1} \cdot X_{2} \cdots X_{d}$ satisfying the above condition ( ${ }^{* *)}$ ). If $f=0 \bmod \left\langle l_{1}, \ldots, l_{k}\right\rangle$, where $l_{i}$ 's are linear forms over $\mathbb{K} \supseteq \mathbb{F}$, then $k \geq w$ and for $k=w$, the space $\operatorname{span}_{\mathbb{K}}\left\{l_{1}, \ldots, l_{w}\right\}$ equals the $\mathbb{K}$-span of either the linear forms in $X_{1}$ or the linear forms in $X_{d}$.

We prove the lemma first for $d=3$, and then use this case to prove it for $d>3$.
Case $[d=3]$ : There is an $A \in \operatorname{GL}(n, \mathbb{F})$ such that $f(A \cdot \mathbf{x})$ is computed by $\left(y_{1} y_{2} \ldots y_{w}\right) \cdot\left(r_{i j}\right)_{i, j \in[w]}$. $\left(z_{1} z_{2} \ldots z_{w}\right)^{T}$, where $\mathbf{y}=\left\{y_{i}\right\}_{i \in[w]}, \mathbf{r}=\left\{r_{i j}\right\}_{i, j \in[w]}$ and $\mathbf{z}=\left\{z_{j}\right\}_{j \in[w]}$ are distinct variables in $\mathbf{x}$. If $f=0 \bmod \left\langle l_{1}, \ldots, l_{k}\right\rangle$, then $f(A \cdot \mathbf{x})=0 \bmod \left\langle l_{1}(A \cdot \mathbf{x}), \ldots, l_{k}(A \cdot \mathbf{x})\right\rangle$. Next, we show that if $f(A \cdot \mathbf{x})=0$ modulo $k^{\prime}$ linear forms $h_{1}, \ldots, h_{k^{\prime}} \in \mathbb{K}[\mathbf{y} \uplus \mathbf{z} \uplus \mathbf{r}]$ then $k^{\prime} \geq w$, and for $k^{\prime}=w$, the space $\operatorname{span}_{\mathbb{K}}\left\{h_{1}, \ldots, h_{w}\right\}$ equals either $\operatorname{span}_{\mathbb{K}}\left\{y_{1}, \ldots, y_{w}\right\}$ or $\operatorname{span}_{\mathbb{K}}\left\{z_{1}, \ldots, z_{w}\right\}$. It follows that $k \geq k^{\prime} \geq w$, and for $k=w$, the linear forms $l_{1}(A \cdot \mathbf{x}), \ldots, l_{w}(A \cdot \mathbf{x})$ must belong to $\mathbb{K}[\mathbf{y} \uplus \mathbf{z} \uplus \mathbf{r}]^{30}$ and hence $\operatorname{span}_{\mathbb{K}}\left\{l_{1}, \ldots, l_{w}\right\}$ equals the $\mathbb{K}$-span of either the linear forms in $X_{1}$ or the linear forms in $X_{d}$.

Reusing symbols, assume that $f$ is computed by $X_{1} \cdot X_{2} \cdot X_{3}$, where $X_{1}=\left(y_{1} y_{2} \ldots y_{w}\right), X_{2}=$ $\left(r_{i j}\right)_{i, j \in[w]}$ and $X_{3}=\left(z_{1} z_{2} \ldots z_{w}\right)^{T}$, and $f=0 \bmod \left\langle l_{1}, \ldots, l_{k}\right\rangle$, where $l_{i}$ 's are linear forms in $\mathbb{K}[\mathbf{y} \uplus \mathbf{z} \uplus \mathbf{r}]$. Suppose $k \leq w$; otherwise, we have nothing to prove. Consider the reduced Gröbner

[^14]basis ${ }^{31} G$ of the ideal $\left\langle l_{1}, \ldots, l_{k}\right\rangle$ with respect to the lexicographic monomial ordering defined by $\mathbf{y} \succ \mathbf{z} \succ \mathbf{r}$. There are sets $S_{\mathbf{y}}, S_{\mathbf{z}} \subseteq[w]$ and $S_{\mathbf{r}} \subseteq[w] \times[w]$, satisfying $\left|S_{\mathbf{y}}\right|+\left|S_{\mathbf{z}}\right|+\left|S_{\mathbf{r}}\right| \leq k$, such that $G$ consists of linear forms of the kind:
\[

$$
\begin{aligned}
y_{i}-g_{i}(\mathbf{y}, \mathbf{z}, \mathbf{r}) & \text { for } i \in S_{\mathbf{y}}, \\
z_{j}-h_{j}(\mathbf{z}, \mathbf{r}) & \text { for } j \in S_{\mathbf{z}} \\
r_{\ell e}-p_{\ell e}(\mathbf{r}) & \text { for }(\ell, e) \in S_{\mathbf{r}},
\end{aligned}
$$
\]

where $g_{i}, h_{j}$ and $p_{\ell e}$ are linear forms over $\mathbb{K}$ in their respective sets of variables. Let $X_{1}^{\prime}, X_{2}^{\prime}, X_{3}^{\prime}$ be the linear matrices obtained from $X_{1}, X_{2}, X_{3}$ respectively, by replacing $y_{i}$ by $g_{i}(\mathbf{y}, \mathbf{z}, \mathbf{r}), r_{\ell e}$ by $p_{\ell e}(\mathbf{r})$ and $z_{j}$ by $h_{j}(\mathbf{z}, \mathbf{r})$, for $i \in S_{\mathbf{y}},(\ell, e) \in S_{\mathbf{r}}$ and $j \in S_{\mathbf{z}}$. Then,

$$
\begin{equation*}
X_{1}^{\prime} \cdot X_{2}^{\prime} \cdot X_{3}^{\prime}=0 \tag{7}
\end{equation*}
$$

The dimension of the $\mathbb{K}$-span of the linear forms of $X_{1}^{\prime}$ is at least $\left(w-\left|S_{\mathbf{y}}\right|\right)$, that of $X_{2}^{\prime}$ is at least $\left(w^{2}-\left|S_{\mathbf{r}}\right|\right)$, and of $X_{3}^{\prime}$ is at least $\left(w-\left|S_{\mathbf{z}}\right|\right)$. Also, there are $C, D \in G L(w, \mathbb{K})$ such that $X_{1}^{\prime} \cdot C, D \cdot X_{3}^{\prime}$ are obtained ${ }^{32}$ from $X_{1}, X_{3}$ respectively, by replacing $y_{i}$ by $g_{i}(0, \mathbf{z}, \mathbf{r})$ and $z_{j}$ by $h_{j}(0, \mathbf{r})$, for $i \in S_{\mathbf{y}}$ and $j \in S_{z}$. Consider the following equation,

$$
\begin{equation*}
\left(X_{1}^{\prime} C\right) \cdot\left(C^{-1} X_{2}^{\prime} D^{-1}\right) \cdot\left(D X_{3}^{\prime}\right)=0 . \tag{8}
\end{equation*}
$$

By examining the L.H.S, we can conclude that for $s \in[w] \backslash S_{\mathbf{y}}$ and $t \in[w] \backslash S_{\mathbf{z}}$, the coefficient of the monomial $y_{s} z_{t}$ over $\mathbb{K}(\mathbf{r})$ is the ( $\left.s, t\right)$-th entry of $C^{-1} X_{2}^{\prime} D^{-1}$ which must be zero. Hence, the dimension of the $\mathbb{K}$-span of the linear forms in $C^{-1} X_{2}^{\prime} D^{-1}$ is at most $w^{2}-\left(w-\left|S_{\mathbf{y}}\right|\right)\left(w-\left|S_{\mathbf{z}}\right|\right)$. As the dimension of the $\mathbb{K}$-span of the linear forms in $X_{2}^{\prime}$ remains unaltered under left and right multiplications by elements in $\mathrm{GL}(w, \mathbb{K})$, we get the relation

$$
\begin{aligned}
w^{2}-\left|S_{\mathbf{r}}\right| & \leq w^{2}-\left(w-\left|S_{\mathbf{y}}\right|\right)\left(w-\left|S_{\mathbf{z}}\right|\right) \\
\Rightarrow\left(w-\left|S_{\mathbf{y}}\right|\right)\left(w-\left|S_{\mathbf{z}}\right|\right) & \leq\left|S_{\mathbf{r}}\right| \\
\Rightarrow w^{2}-\left(\left|S_{\mathbf{y}}\right|+\left|S_{\mathbf{z}}\right|\right) w+\left|S_{\mathbf{y}}\right| \cdot\left|S_{\mathbf{z}}\right| & \leq\left|S_{\mathbf{r}}\right| \\
\Rightarrow w^{2}-\left(w-\left|S_{\mathbf{r}}\right|\right) w+\left|S_{\mathbf{y}}\right| \cdot\left|S_{\mathbf{z}}\right| & \leq\left|S_{\mathbf{r}}\right|, \quad \text { as }\left|S_{\mathbf{y}}\right|+\left|S_{\mathbf{z}}\right|+\left|S_{\mathbf{r}}\right| \leq k \leq w \\
\Rightarrow\left|S_{\mathbf{r}}\right| w+\left|S_{\mathbf{y}}\right| \cdot\left|S_{\mathbf{z}}\right| & \leq\left|S_{\mathbf{r}}\right| .
\end{aligned}
$$

As $\left|S_{\mathbf{y}}\right|,\left|S_{\mathbf{z}}\right|,\left|S_{\mathbf{r}}\right| \geq 0$, we must have $\left|S_{\mathbf{r}}\right|=0$, and either $\left|S_{\mathbf{y}}\right|=0$ or $\left|S_{\mathbf{z}}\right|=0$.
Suppose $\left|S_{\mathbf{r}}\right|=\left|S_{\mathbf{z}}\right|=0$ (the case for $\left|S_{\mathbf{r}}\right|=\left|S_{\mathbf{y}}\right|=0$ is similar). Then, Equation 8 simplifies to

$$
\left(X_{1}^{\prime} C\right) \cdot\left(C^{-1} X_{2}\right) \cdot X_{3}=0
$$

If $k<w$ then there is a $y_{s}$ in $X_{1}$ that is not replaced while forming $X_{1}^{\prime} C$ from $X_{1}$. By examining the coefficient of $y_{s}$ over $\mathbb{K}(\mathbf{r}, \mathbf{z})$ in the L.H.S of the above equation, we arrive at a contradiction. Hence, $k=w$, in which case Equation 7 simplifies to

$$
X_{1}^{\prime} \cdot X_{2} \cdot X_{3}=0
$$

[^15]The entries of $X_{1}^{\prime}$ are linear forms in $\mathbf{z}$ and $\mathbf{r}$, and so $X_{1}^{\prime}=X_{1}^{\prime}(\mathbf{z})+X_{1}^{\prime}(\mathbf{r})$ where the entries of $X_{1}^{\prime}(\mathbf{z})$ (similarly, $X_{1}^{\prime}(\mathbf{r})$ ) are linear forms in $\mathbf{z}$ (respectively, $\mathbf{r}$ ). The above equation implies

$$
X_{1}^{\prime}(\mathbf{z}) \cdot X_{2} \cdot X_{3}=0 \quad \text { and } \quad X_{1}^{\prime}(\mathbf{r}) \cdot X_{2} \cdot X_{3}=0
$$

as the two L.H.S above are monomial disjoint. It is now easy to argue that $X_{1}^{\prime}(\mathbf{z})=X_{1}^{\prime}(\mathbf{r})=0$, implying $X_{1}^{\prime}=0$ and hence the reduced Gröbner basis $G$ is in fact $\left\{y_{1}, \ldots, y_{w}\right\}$.

Case [d>3]: As before, by applying an invertible transformation, we can assume that $X_{1}=$ $\left(y_{1} y_{2} \ldots y_{w}\right), X_{2}=\left(r_{i j}\right)_{i, j \in[w]}$ and $X_{d}=\left(z_{1} z_{2} \ldots z_{w}\right)^{T}$. Let $\mathbf{u}=\mathbf{x} \backslash(\mathbf{y} \uplus \mathbf{z} \uplus \mathbf{r})$ and $k \leq w$. Consider the reduced Gröbner basis $G$ of the ideal $\left\langle l_{1}, l_{2}, \ldots, l_{k}\right\rangle$ with respect to the lexicographic monomial ordering defined by $\mathbf{u} \succ \mathbf{y} \succ \mathbf{z} \succ \mathbf{r}$. There are sets $S_{\mathbf{u}} \subseteq\left[n-w^{2}-2 w\right], S_{\mathbf{y}}, S_{\mathbf{z}} \subseteq[w]$ and $S_{\mathbf{r}} \subseteq\left[w^{2}\right]$, satisfying $\left|S_{\mathbf{u}}\right|+\left|S_{\mathbf{y}}\right|+\left|S_{\mathbf{z}}\right|+\left|S_{\mathbf{r}}\right| \leq k$, such that $G$ consists of linear forms of the kind:

$$
\begin{aligned}
u_{m}-t_{m}(\mathbf{u}, \mathbf{y}, \mathbf{z}, \mathbf{r}) & \text { for } m \in S_{\mathbf{u}} \\
y_{i}-g_{i}(\mathbf{y}, \mathbf{z}, \mathbf{r}) & \text { for } i \in S_{\mathbf{y}} \\
z_{j}-h_{j}(\mathbf{z}, \mathbf{r}) & \text { for } j \in S_{\mathbf{z}} \\
r_{\ell e}-p_{\ell e}(\mathbf{r}) & \text { for }(\ell, e) \in S_{\mathbf{r}},
\end{aligned}
$$

where $t_{m}, g_{i}, h_{j}$ and $p_{\ell e}$ are linear forms over $\mathbb{K}$ in their respective sets of variables. Let $X^{\prime}$ be the matrix obtained from $X$ by replacing $u_{m}$ by $t_{m}(\mathbf{u}, \mathbf{y}, \mathbf{z}, \mathbf{r}), y_{i}$ by $g_{i}(\mathbf{y}, \mathbf{z}, \mathbf{r}), z_{j}$ by $h_{j}(\mathbf{z}, \mathbf{r})$, and $r_{\ell e}$ by $p_{\ell e}(\mathbf{r})$, for $m \in S_{\mathbf{u}}, i \in S_{\mathbf{y}}, j \in S_{\mathbf{z}}$, and $(\ell, e) \in S_{\mathbf{r}}$. Then,

$$
X_{1}^{\prime} \cdot X_{2}^{\prime} \cdot X_{3}^{\prime} \ldots X_{d}^{\prime}=0
$$

Let $X(\mathbf{u}) \stackrel{\text { def }}{=}(X)_{\mathbf{y}=\mathbf{z}=\mathbf{r}=0}$. By treating the L.H.S of the above equation as a polynomial in $\mathbf{u}$-variables with coefficients from $\mathbb{K}(\mathbf{y}, \mathbf{z}, \mathbf{r})$ and focusing on the degree- $(d-3)$ homogeneous component of this polynomial, we have

$$
\begin{equation*}
X_{1}^{\prime} \cdot X_{2}^{\prime} \cdot X_{3}^{\prime}(\mathbf{u}) \ldots X_{d-1}^{\prime}(\mathbf{u}) \cdot X_{d}^{\prime}=0 \tag{9}
\end{equation*}
$$

If $X_{3}^{\prime}(\mathbf{u}) \cdots X_{d-1}^{\prime}(\mathbf{u}) \in \operatorname{GL}(w, \mathbb{K}(\mathbf{u}))$ then there is a $\mathbf{c} \in \mathbb{F}^{|\mathbf{u}|}$ such that $C=X_{3}^{\prime}(\mathbf{c}) \cdots X_{d-1}^{\prime}(\mathbf{c}) \in$ $G L(w, \mathbb{K})$. Define

$$
f_{1}=X_{1} \cdot X_{2} \cdot C \cdot X_{d}
$$

and observe that Equation 9 implies $f_{1}$ is zero modulo the linear forms,

$$
\begin{aligned}
y_{i}-g_{i}(\mathbf{y}, \mathbf{z}, \mathbf{r}) & \text { for } i \in S_{\mathbf{y}} \\
z_{j}-h_{j}(\mathbf{z}, \mathbf{r}) & \text { for } j \in S_{\mathbf{z}} \\
r_{\ell e}-p_{\ell e}(\mathbf{r}) & \text { for }(\ell, e) \in S_{\mathbf{r}}
\end{aligned}
$$

By applying Case $\left[\mathrm{d}=3\right.$ ] on $f_{1}$, we get the desired conclusion, i.e. $k=w$ and the $\mathbb{K}$-span of the above linear forms (hence also that of $\left\{l_{1}, \ldots, l_{k}\right\}$ ) is either $\operatorname{span}_{\mathbb{K}}\left\{y_{1}, \ldots, y_{w}\right\}$ or $\operatorname{span}_{\mathbb{K}}\left\{z_{1}, \ldots, z_{w}\right\}$. So, suppose $X_{3}^{\prime}(\mathbf{u}) \cdots X_{d-1}^{\prime}(\mathbf{u}) \notin \mathrm{GL}(w, \mathbb{K}(\mathbf{u}))$ in Equation 9. Then, there is a $j \in[3, d-1]$ such that $\operatorname{det}\left(X_{j}^{\prime}(\mathbf{u})\right)=0$. Observe that $X_{i}^{\prime}(\mathbf{u})$ can be obtained from $X_{i}(\mathbf{u})$ by replacing $u_{m}$ by $t_{m}(\mathbf{u}, 0,0,0)$ for $m \in S_{\mathbf{u}}$. That is,

$$
X_{i}^{\prime}(\mathbf{u})=X_{i}(\mathbf{u}) \quad \bmod \left\langle\left\{u_{m}-t_{m}(\mathbf{u}, 0,0,0)\right\}_{m \in S_{\mathbf{u}}}\right\rangle, \quad \text { for every } i \in[3, d-1] .
$$

As $X_{j}(\mathbf{u})$ is full $\operatorname{rank}^{33}$ and $\operatorname{det}\left(X_{j}^{\prime}(\mathbf{u})\right)=0$, the fact below implies $\left|S_{\mathbf{u}}\right|=w,\left|S_{\mathbf{y}}\right|=\left|S_{\mathbf{z}}\right|=\left|S_{\mathbf{r}}\right|=0$.

[^16]Observation 4.1. If the symbolic determinant $\operatorname{Det}_{w}$ is zero modulo s linear forms then $s \geq w$.
Hence, Equation 9 simplifies to

$$
\begin{align*}
X_{1} \cdot X_{2} \cdot X_{3}^{\prime}(\mathbf{u}) \ldots X_{d-1}^{\prime}(\mathbf{u}) \cdot X_{d} & =0 \\
\Rightarrow X_{3}^{\prime}(\mathbf{u}) \cdots X_{d-1}^{\prime}(\mathbf{u}) & =0 \tag{10}
\end{align*}
$$

The above equality can not happen and this can be argued by applying induction on the number of matrices in the L.H.S of Equation 10:

Base case: $(d=4)$ The L.H.S of Equation 10 has one matrix $X_{3}^{\prime}(\mathbf{u})$. As $X_{3}(\mathbf{u})$ is full rank $^{33}$, it cannot vanish modulo $w$ linear forms.
Induction hypothesis: Equation 10 does not hold if the L.H.S has at most $d-4$ matrices. Inductive step: $(d>4)$ Suppose Equation 10 is true. As the $2 w^{2}$ linear forms in $X_{3}(\mathbf{u})$ and $X_{d-1}(\mathbf{u})$ are linearly independent ${ }^{33}$, by Observation 4.1, at least one of $X_{3}^{\prime}(\mathbf{u})$ and $X_{d-1}^{\prime}(\mathbf{u})$ is invertible. This gives a shorter product where we can apply the induction hypothesis to get a contradiction.

### 4.2 Finding the coefficients in the intermediate matrices

Following the notations in Section 1.3.2, $\mathbf{y}=\left\{y_{1}, \ldots, y_{w}\right\}$ and $\mathbf{z}=\left\{z_{1}, \ldots, z_{w}\right\}$ are subsets of $\mathbf{x}$, $\mathbf{r}=\mathbf{x} \backslash(\mathbf{y} \uplus \mathbf{z}), X_{1}^{\prime}=\left(y_{1} y_{2} \ldots y_{w}\right)$ and $X_{d}^{\prime}=\left(z_{1} z_{2} \ldots z_{w}\right)^{T}$. When Algorithm 2 reaches the third and final stage, it has blackbox access to a $f^{\prime} \in \mathbb{F}[\mathbf{x}]$ and linear matrices $S_{2}, \ldots, S_{d-1} \in \mathbb{L}[\mathbf{r}]{ }^{w \times w}$ returned by Algorithm 1, such that $S_{2} \cdot S_{3} \ldots S_{d-1}$ is the linear matrix factorization of a random $(w, d-2, n-2 w)$-matrix product $R_{2} \cdot R_{3} \ldots R_{d-1}$ over $\mathbb{F}$. Further, there exist linear matrices $T_{2}, \ldots, T_{d-1} \in \mathbb{L}[\mathbf{x}]^{w \times w}$ satisfying $\left(T_{k}\right)_{\mathbf{y}=0, \mathbf{z}=0}=S_{k}$ for every $k \in[2, d-1]$, such that $f^{\prime}$ is computed by the ABP $X_{1}^{\prime} \cdot T_{2} \ldots T_{d-1} \cdot X_{d-1}^{\prime}$. The task for Algorithm 6 is to efficiently compute the coefficients of the $\mathbf{y}$ and $\mathbf{z}$ variables in $T_{k}$. At a high level, this is made possible because of the uniqueness of such $T_{k}$ matrices: Indeed the analysis of Algorithm 6 shows that with high probability the coefficients of $\mathbf{y}$ and $\mathbf{z}$ in $T_{3}, \ldots, T_{d-2}$ are uniquely determined, and (if a certain canonical form is assumed then) the same is true for matrices $T_{2}$ and $T_{d-1}$.

Canonical form for $T_{2}$ and $T_{d-1}$ : Matrix $T_{2}$ is said to be in canonical form if for every $l \in[w]$ the coefficient of $y_{l}$ is zero in the linear form at the $(i, j)$-th entry of $T_{2}$, whenever $i>l$. Similarly, $T_{d-1}$ is in canonical form if for every $l \in[w]$ the coefficient of $z_{l}$ is zero in the linear form at the $(i, j)$-th entry of $T_{d-1}$ whenever $j>l$. It can be verified (see [KNST17]), if $f^{\prime}$ is computed by an ABP $X_{1}^{\prime} \cdot T_{2} \ldots T_{d-1} \cdot X_{d-1}^{\prime}$ then it is computed by another ABP where the corresponding $T_{2}$ and $T_{d-1}$ are in canonical form, and the other matrices remain unchanged.

Linear independence of minors of a random $A B P$ : The lemma given below is the reason Algorithm 6 is able to reduce the task of finding the coefficients of the $\mathbf{y}$ and $\mathbf{z}$ variables to solving linear equations. In the following discussion, the $i$-th row and $j$-th column of a matrix $M$ will be denoted by $M(i, *)$ and $M(*, j)$ respectively.

Let $R_{2} \cdot R_{3} \ldots R_{d-1}$ be a random ( $w, d-2, n-2 w$ )-matrix product in $\mathbf{r}$-variables over $\mathbb{F}$. For every $s, t \in[w], R_{2}(s, *) \cdot R_{3} \ldots R_{d-2} \cdot R_{d-1}(*, t)$ is a random $(w, d-2, n-2 w)$-ABP having a total of $w^{2}(d-4)+2 w$ linear forms in all the $R_{k}$ matrices. Let us index the linear forms ${ }^{34}$ by $\left[w^{2}(d-4)+\right.$

[^17]$2 w]$. We associate a polynomial $g_{e}^{(s, t)}$ with the $e$-th linear form, for every $e \in\left[w^{2}(d-4)+2 w\right]$, as follows: If the $e$-th linear form is the $(\ell, m)$-th entry of $R_{k}$ then
$$
g_{e}^{(s, t)}(\mathbf{r}) \stackrel{\text { def }}{=}\left[R_{2}(s, *) \cdot R_{3} \ldots R_{k-2} \cdot R_{k-1}(*, \ell)\right] \cdot\left[R_{k+1}(m, *) \cdot R_{k+2} \ldots R_{d-2} \cdot R_{d-1}(*, t)\right] .^{35}
$$

The polynomials $\left\{g_{e}^{(s, t)}: e \in\left[w^{2}(d-4)+2 w\right]\right\}$, will be called the minors of the $\operatorname{ABP} R_{2}(s, *)$. $R_{3} \ldots R_{d-2} \cdot R_{d-1}(*, t)$.

Lemma 4.2. With probability $1-(w d n)^{-\Omega(1)}$ over the randomness of $R_{2} \cdots R_{d-1}$ the following holds: For every $s, t \in[w]$, the minors $\left\{g_{e}^{(s, t)}: e \in\left[w^{2}(d-4)+2 w\right]\right\}$, are $\mathbb{F}$-linearly independent.

The proof of the lemma is given at the end of this section. Due to the uniqueness of factorization, the matrices $S_{2}, \ldots, S_{d-1}$ in Algorithm 2 are related to $R_{2}, \ldots, R_{d-1}$ as follows: There are $C_{i}, D_{i} \in$ $\mathrm{GL}(w, \mathbb{L})$ such that $S_{i}=C_{i} \cdot R_{i} \cdot D_{i}$, for every $i \in[2, d-1]$; moreover, there are $c_{2}, \ldots, c_{d-2} \in \mathbb{L}^{\times}$ satisfying $C_{2}=D_{d-1}=I_{w}, D_{i} \cdot C_{i+1}=c_{i} I_{w}$ for $i \in[2, d-2]$, and $\prod_{i=2}^{d-2} c_{i}=1$. Define minors of the ABP $S_{2}(s, *) \cdot S_{3} \ldots S_{d-2} \cdot S_{d-1}(*, t)$, for every $s, t \in[w]$, like above. The edges of the ABP are indexed by $\left[w^{2}(d-4)+2 w\right]$ and a polynomial $h_{e}^{(s, t)}$ is associated with the $e$-th linear form as follows: If the $e$-th linear form is the $(\ell, m)$-th entry of $S_{k}$ then

$$
\begin{equation*}
h_{e}^{(s, t)}(\mathbf{r}) \stackrel{\text { def }}{=}\left[S_{2}(s, *) \cdot S_{3} \ldots S_{k-2} \cdot S_{k-1}(*, \ell)\right] \cdot\left[S_{k+1}(m, *) \cdot S_{k+2} \ldots S_{d-2} \cdot S_{d-1}(*, t)\right] . \tag{11}
\end{equation*}
$$

It is a simple exercise to derive the following corollary from the lemma above.
Corollary 4.1. With probability $1-(w d n)^{-\Omega(1)}$ the following holds: For every $s, t \in[w]$, the minors $\left\{h_{e}^{(s, t)}: e \in\left[w^{2}(d-4)+2 w\right]\right\}$ are $\mathbb{L}$-linearly independent.

We are now ready to argue the correctness of Algorithm 6 by tracing its steps.

1. Computing the partial derivatives (Step 2): In this step, we compute all the third order partial derivatives of $f^{\prime}$ using Claim 2.1.
2. Computing almost all the coefficients of the $\mathbf{y}$ and $\mathbf{z}$ variables (Steps 6-13): Equations 12 and 13 are justified by treating $f^{\prime}$ as a polynomial in the $\mathbf{y}$ and $\mathbf{z}$ variables with coefficients from $\mathbb{L}(\mathbf{r})$, and examining the coefficients of $y_{s}^{2} z_{t}$ and $y_{s} z_{t}^{2}$ respectively. A linear system obtained at step 9 or step 11 has $w^{2}(d-4)+2 w$ variables and the same number of linear equations. Corollary 4.1, together with Claim 2.2, ensure that the square coefficient matrix of the linear system is invertible (with high probability), and hence the solution computed is unique. The uniqueness implies that the solutions obtained across multiple iterations of the loop do not conflict with each other ${ }^{36}$. This also shows that the matrices $T_{3}, \ldots, T_{4}$ are unique. By the end of this stage, the coefficients of $\mathbf{y}$ and $\mathbf{z}$ variables are computed for all the linear forms, except for the coefficients of $y_{l}$ in $T_{2}(s, *)$ for $l>s$, and the coefficients of $z_{l}$ in $T_{d-1}(*, t)$ for $l>t$. These coefficients are retrieved in the next stage.
[^18]```
Algorithm 6 Computing the coefficients of \(\mathbf{y}\) and \(\mathbf{z}\) variables in \(T_{k}\)
    INPUT: Blackbox access to \(f^{\prime}\) and linear matrices \(S_{2}, \ldots, S_{d-1} \in \mathbb{L}[\mathbf{r}]^{w \times w}\).
```

    OUTPUT: Linear matrices \(T_{2}, T_{3}, \ldots, T_{d-1} \in \mathbb{L}[\mathbf{x}]^{w \times w}\) such that \(f^{\prime}\) is computed by \(\mathbf{y} \cdot T_{2}\).
    \(T_{3} \ldots T_{d-1} \cdot \mathbf{z}^{T}\), satisfying \(\left(T_{k}\right)_{\mathbf{y}=0, \mathbf{z}=0}=S_{k}\) for every \(k \in[2, d-1]\).
    . \(/ *\) Computing the partial derivatives */
    Compute blackbox access to \(\left(\frac{\partial f^{\prime}}{\partial y_{s} y_{l} z_{t}}\right)_{\mathbf{y}=0, \mathbf{z}=0}\) and \(\left(\frac{\partial f^{\prime}}{\partial y_{s} z_{l} z_{t}}\right)_{\mathbf{y}=0, \mathbf{z}=0}\) for all \(s, l, t \in[w]\).
    3. For every \(s, t \in[w]\), let \(\left\{h_{e}^{(s, t)}: e \in\left[w^{2}(d-4)+2 w\right]\right\}\) be the minors of the ABP \(S_{2}(s, *)\).
    \(S_{3} \ldots S_{d-2} \cdot S_{d-1}(*, t)\), as defined in Equation 11.
    4.
    /* Computing almost all the coefficients of the \(\mathbf{y}\) and z variables in \(T_{k}{ }^{*} /\)
    Set \(E=w^{2}(d-4)+2 w\).
    for every \(s, t \in[w]\) do
        Pick \(\mathbf{a}_{1}, \ldots, \mathbf{a}_{E} \in \in_{r} \mathbb{F}^{|\mathbf{r}|}\) independently.
    9. Solve the linear system over \(\mathbb{L}\) defined by
    $$
\begin{equation*}
\sum_{e \in[E]} c_{e} \cdot h_{e}^{(s, t)}\left(\mathbf{a}_{i}\right)=\left(\frac{\partial f^{\prime}}{\partial y_{s}^{2} z_{t}}\right)_{\mathbf{y}=0, \mathbf{z}=0}\left(\mathbf{a}_{i}\right), \quad \text { for } i \in[E] \tag{12}
\end{equation*}
$$

for a unique solution of $\left\{c_{e}\right\}_{e \in[E]}$. If the coefficient matrix is not invertible, output 'Failed'.
10. For every $e \in[E]$, set the solution value of $c_{e}$ as the coefficient of $y_{s}$ in the $e$-th linear form of the ABP $T_{2}(s, *) \cdot T_{3} \ldots T_{d-2} \cdot T_{d-1}(*, t)$.
11. Solve the linear system over $\mathbb{L}$ defined by

$$
\begin{equation*}
\sum_{e \in[E]} d_{e} \cdot h_{e}^{(s, t)}\left(\mathbf{a}_{i}\right)=\left(\frac{\partial f^{\prime}}{\partial y_{s} z_{t}^{2}}\right)_{\mathbf{y}=0, \mathbf{z}=0}\left(\mathbf{a}_{i}\right), \quad \text { for } i \in[E], \tag{13}
\end{equation*}
$$

for a unique solution of $\left\{d_{e}\right\}_{e \in[E]}$.
12. For every $e \in[E]$, set the solution value of $d_{e}$ as the coefficient of $z_{t}$ in the $e$-th linear form of the $\operatorname{ABP} T_{2}(s, *) \cdot T_{3} \ldots T_{d-2} \cdot T_{d-1}(*, t)$.
end for
14.
15. /* Computing the remaining y and z coefficients in $T_{2}$ and $T_{d-1} * /$
6. for every $s, t \in[w]$ do
17. For every $l>s$, compute the coefficients of $y_{l}$ in the linear forms in $T_{2}(s, *)$ by setting up a linear system similar to Equation 12, but with the R.H.S replaced by $\frac{\partial f^{\prime}}{\partial y_{s} y_{l} z_{1}}$.
18. For every $l>t$, compute the coefficients of $z_{l}$ in the linear forms in $T_{d-1}(*, t)$ by setting up a linear system similar to Equation 13, but with the R.H.S replaced by $\frac{\partial f^{\prime}}{\partial y_{1} z_{z_{t}}}$.
end for
20.
21. The coefficients of the $\mathbf{r}$ variables in the linear forms in $T_{k}$ remain the same as that in $S_{k}$, for all $k \in[2, d-1]$. Output $T_{2}, T_{3}, \ldots T_{d-1}$.
3. Computing the remaining $\mathbf{y}$ and $\mathbf{z}$ coefficients in $T_{2}$ and $T_{d-1}$ (Steps 16-19): For an $s \in[w]$, consider the following minors of $S_{2}(s, *) \cdot S_{3} \ldots S_{d-2} \cdot S_{d-1}(*, 1)$ :

$$
S_{3}(m, *) \cdot S_{4} \ldots S_{d-2} \cdot S_{d-1}(*, 1) \quad \text { for all } m \in[w] .
$$

Without loss of generality, let these minors be $h_{1}^{(s, 1)}, \ldots, h_{w}^{(s, 1)}$. Let $l>s$. By treating $f^{\prime}$ as a polynomial in the $\mathbf{y}, \mathbf{z}$ variables, with coefficients from $\mathbb{L}(\mathbf{r})$, and examining the coefficient of $y_{s} y_{l} z_{1}$ in $f^{\prime}$, we arrive at the equation,

$$
\sum_{e=1}^{w} c_{e} \cdot h_{e}^{(s, 1)}+K(\mathbf{r})=\left(\frac{\partial f^{\prime}}{\partial y_{s} y_{l} z_{1}}\right)_{\mathbf{y}=0, \mathbf{z}=0},
$$

where $c_{1}, \ldots, c_{w}$ are the unknown coefficients of $y_{l}$ in the linear forms of $T_{2}(s, *)$, and $K(\mathbf{r})$ is a known linear combination of some other minors. The fact that $K(\mathbf{r})$ is known at step 17 follows from this observation - while forming a monomial $y_{s} y_{l} z_{1}$, we either choose $y_{s}$ from $X_{1}^{\prime}$ and $y_{l}$ from $T_{2}(s, *)$ or $T_{3}, \ldots, T_{d-1}(*, 1)$, or $y_{l}$ from $X_{1}^{\prime}$ and $y_{s}$ from $T_{3}, \ldots, T_{d-1}(*, 1)$. In the latter case, we are using the fact that $T_{2}$ is in canonical form, and so $y_{s}$ does not appear in $T_{2}(l, *)$. As the coefficients of $y_{s}, y_{l}$ in $T_{3}, \ldots, T_{d-1}(*, 1)$ are known from the computation in steps $6-13$, we conclude that $K(\mathbf{r})$ in known. Thus, we can solve for $c_{1}, \ldots, c_{w}$ by plugging in $w$ random points in place of the $\mathbf{r}$ variables and setting up a linear system in $w$ variables. Corollary 4.1 and Claim 2.2 imply the $w \times w$ coefficient matrix of the system is invertible, and hence the solution for $c_{1}, \ldots, c_{w}$ is unique. The correctness of step 18 can be argued similarly, and this finally implies that $T_{2}$ and $T_{d-1}$ (in canonical form) are unique.

## Linear independence of minors: Proof of Lemma 4.2

We have to show that the minors of $R_{2}(s, *) \cdot R_{3} \ldots R_{d-2} \cdot R_{d-1}(*, t)$ are $\mathbb{F}$-linearly independent with high probability, for every $s, t \in[w]$, where $R_{2} \cdot R_{3} \ldots R_{d-1}$ is a random ( $w, d-2, n-2 w$ )matrix product. We will prove it for a fixed $s, t \in[w]$, and then by union bound the result will follow for every $s, t \in[w]$. As $n \geq 4 w^{2}$, we have $n-2 w \geq 3 w^{2}$. So, it is sufficient to show the linear independence of the minors of a random $(w, d, n)$-ABP $X_{1} \cdot X_{2} \ldots X_{d}$ in x-variables, for $n \geq 3 w^{2}$.

Treat the coefficients of the linear forms in $X_{1}, \ldots, X_{d}$ as formal variables. In particular,

$$
\begin{equation*}
X_{1}=\sum_{i=1}^{n} U_{i}^{(1)} x_{i}, \quad X_{k}=\sum_{i=1}^{n} U_{i}^{(k)} x_{i} \text { for } k \in[2, d-1], \quad X_{d}=\sum_{i=1}^{n} U_{i}^{(d)} x_{i}, \tag{14}
\end{equation*}
$$

where $U_{i}^{(1)}$ and $U_{i}^{(d)}$ are row and column vectors of length $w$ respectively, $U_{i}^{(k)}$ is a $w \times w$ matrix, and the entries of these matrices are distinct $\mathbf{u}$-variables. We will denote the $(\ell, m)$-th entry of $U_{i}^{(k)}$ by $U_{i}^{(k)}(\ell, m)$, and the $m$-th entry of $U_{i}^{(d)}$ by $U_{i}^{(d)}(m)$. From the above equations, $X_{1} \cdot X_{2} \ldots X_{d}$ is a $(w, d, n)$-ABP over $\mathbb{F}(\mathbf{u})$. We will show in the following claim that the minors of this ABP are $\mathbb{F}(\mathbf{u})$ linearly independent. As the coefficients of the $\mathbf{x}$-monomials of these minors are polynomials (in fact, multilinear polynomials) of degree $d-1$ in the $\mathbf{u}$-variables, an application of the SchwartzZippel lemma implies $\mathbb{F}$-linear independence of the minors (with high probability) when the $\mathbf{u}$ variables are set randomly to elements in $\mathbb{F}$ (as is done in a random ABP over $\mathbb{F}$ ).

Claim 4.1. The minors of $X_{1} \cdot X_{2} \ldots X_{d}$ are $\mathbb{F}(\mathbf{u})$-linearly independent.

Proof. We will prove by induction on $d$.
Base case ( $d=3$ ): Clearly, if the minors are $\mathbb{F}$-linearly independent after setting the $\mathbf{u}$-variables to some $\mathbb{F}$-elements then the minors are also $\mathbb{F}(\mathbf{u})$-linearly independent before the setting. As $n \geq w^{2}+2 w$, it is possible to set the $\mathbf{u}$-variables in $X_{1}, X_{2}, X_{3}$ such that the entries of these matrices (after the setting) become distinct $\mathbf{x}$-variables. The minors of this $\mathbf{u}$-evaluated ABP $X_{1} \cdot X_{2} \cdot X_{3}$ are monomial disjoint and so $\mathbb{F}$-linearly independent.

Inductive step: Split the $w^{2}(d-2)+2 w$ minors of $X_{1} \cdot X_{2} \ldots X_{d}$ into two sets: The first set $G_{1}$ consists of minors $g_{e}$, for $e \in\left[w^{2}(d-3)+2 w\right]$, such that the $e$-th linear form is the $(\ell, m)$-th entry of some matrix $X_{k}$ satisfying $k \neq d$ and if $k=d-1$ then $m=w$. The second set $G_{2}$ consists of minors $g_{e}$, for $e \in\left[w^{2}(d-3)+2 w+1, w^{2}(d-2)+2 w\right]$, such that the $e$-th linear form is either the $(\ell, m)$-th entry of $X_{d-1}$ for $m \neq w$, or the $\ell$-th entry of $X_{d}$. Set $G_{1}$ has $p=w^{2}(d-3)+2 w$ minors and $G_{2}$ has $w^{2}$ minors.

Suppose $\mu_{1}, \ldots, \mu_{p}$ are monomials in $\mathbf{x}$-variables of degree $d-2$. Imagine a $\left(w^{2}(d-2)+2 w\right) \times$ $\left(w^{2}(d-2)+2 w\right)$ matrix $M$ whose rows are indexed by the minors in $G_{1}$ and $G_{2}$, and columns by monomials $\mu_{1} x_{1}, \mu_{2} x_{1}, \ldots, \mu_{p} x_{1}$ and $x_{2}^{d-1}, x_{3}^{d-1}, \ldots, x_{w^{2}+1}^{d-1}$, The $(g, \sigma)$-th entry of $M$ contains the coefficient of the monomial $\sigma$ in $g$, this coefficient is a multilinear polynomial in the $\mathbf{u}$-variables. In a sequence of observations, we show that there exist $\mu_{1}, \ldots, \mu_{p}$ such that $\operatorname{det}(M) \neq 0$.

Consider the variable $u \stackrel{\text { def }}{=} U_{1}^{(d)}(w)$. The following observations are easy to verify.
Observation 4.2. 1. Variable $u$ does not appear in any of the monomials of the $(g, \sigma)$-th entry of $M$ if $g \in G_{2}$ or $\sigma \in\left\{x_{2}^{d-1}, \ldots, x_{w w^{2}+1}^{d-1}\right\}$.
2. Variable $u$ appears in some monomials of the $(g, \sigma)$-th entry of $M$ if $g \in G_{1}$ and $\sigma \in\left\{\mu_{1} x_{1}, \ldots, \mu_{p} x_{1}\right\}$, irrespective of $\mu_{1}, \ldots, \mu_{p}$.
Observation 4.3. Let $g \in G_{1}$ and $\sigma \in\left\{\mu_{1} x_{1}, \ldots, \mu_{p} x_{1}\right\}$. If we treat the $(g, \sigma)$-th entry of $M$ as a polynomial in $u$ with coefficients from $\mathbb{F}[\mathbf{u} \backslash u]$ then the coefficient of $u$ does not depend on the variables:
(a) $U_{i}^{(d)}(j)$ for $j \neq w$ and $i \in[n]$,
(b) $U_{i}^{(d)}(w)$ for $i \in[2, n]$,
(c) $U_{i}^{(d-1)}(\ell, m)$ for $\ell, m \in[w]$ with $m \neq w$, and $i \in[n]$.

Denote the union of the $\mathbf{u}$-variables specified in (a), (b) and (c) of the above observation by $\mathbf{v}$.
Observation 4.4. The set $\left\{g_{v=0}: g \in G_{1}\right\}$ equals the set $\left\{h \cdot u x_{1}: h\right.$ is a minor of $\left.X_{1} \cdot X_{2} \ldots X_{d-1}(*, w)\right\}$.
By the induction hypothesis, the minors of $X_{1} \cdot X_{2} \ldots X_{d-1}(*, w)$, say $h_{1}, \ldots, h_{p}$, are $\mathbb{F}(\mathbf{u})$-linearly independent. Hence there are $p$ monomials in $\mathbf{x}$-variables of degree $d-2$ such that $h_{1}, \ldots, h_{p}$, when restricted to these monomials, are $\mathbb{F}(\mathbf{u})$-linearly independent. These $p$ monomials are our choices for $\mu_{1}, \ldots, \mu_{p}$. Let $N$ be the $p \times p$ matrix with rows indexed by $h_{1}, \ldots, h_{p}$ and columns by $\mu_{1}, \ldots, \mu_{p}$, and $N(h, \mu)$ contains the coefficient of the monomial $\mu$ in $h$. Then, $\operatorname{det}(N) \neq 0$. Under these settings, we have the following observation (which can be derived easily from the above).

Observation 4.5. The coefficient of $u^{p}$ in $\operatorname{det}(M)$, when treated as a polynomial in $u$ with coefficients from $\mathbb{F}[\mathbf{u} \backslash u]$, is $\operatorname{det}(N) \cdot \operatorname{det}\left(M_{0}\right)$, where $M_{0}$ is the submatrix of $M$ defined by rows indexed by $\left\{g: g \in G_{2}\right\}$ and columns by $x_{2}^{d-1}, \ldots, x_{w^{2}+1}^{d-1}$.
The next observation completes the proof of the claim by showing $\operatorname{det}(M) \neq 0$.
Observation 4.6. $\operatorname{det}\left(M_{0}\right) \neq 0$.
The proof of the above follows by noticing that $M_{0}$ looks like $\left(f_{i}\left(\mathbf{u}_{j}\right)\right)_{i, j \in\left[w^{2}\right]}$, where $\mathbf{u}_{1}, \ldots, \mathbf{u}_{w w^{2}}$ are some disjoint subsets of the $\mathbf{u}$-variables and $f_{1}, \ldots, f_{w^{2}}$ are $\mathbb{F}$-linearly independent polynomials. The observation then follows from Claim 2.2.

## 5 Equivalence test for determinant over finite fields

We prove Theorem 3 in this section. It is known that the affine equivalence test can be reduced to equivalence test [Kay12], as briefly explained below.

Reduction to equivalence test: Suppose $f$ is a $(n, w)$-polynomial that is affine equivalent to $\operatorname{Det}_{w}$, where $n \geq w^{2}$. The following claim reduces the number of variables from $n$ to $w^{2}$. A proof can be found in [Kay12] (see also Algorithm 8 and Claim 2.3 in [KNST17]).

Claim 5.1. There is a randomized algorithm that takes input blackbox access to $f(\mathbf{x})$ and with probability $1-\frac{n^{\circ}(1)}{q}$ outputs a matrix $C \in G L(n, \mathbb{F})$ such that $f(C \cdot \mathbf{x})$ is a $\left(w^{2}, w\right)$-polynomial. The algorithm runs in $(n \log q)^{O(1)}$ time.
Suppose $\mathbf{y} \subseteq \mathbf{x}$ is the set of $w^{2}$ variables appearing in $f(C \cdot \mathbf{x})$, and let $g(\mathbf{y})$ be the degree- $w$ homogeneous component of $f(C \cdot \mathbf{x})$ which must be equivalent to $\operatorname{Det}_{w}$. By using an equivalence test for $\operatorname{Det}_{w}$, we can compute a $Q \in G L\left(w^{2}, \mathbb{L}\right)$ such that $g(\mathbf{y})=\operatorname{Det}_{w}(Q \cdot \mathbf{y})$, implying $g(\mathbf{x})=\operatorname{Det}_{w}\left(Q^{\prime} \cdot \mathbf{x}\right)$ where $Q^{\prime} \in \mathbb{L}^{w^{2} \times n}$ is obtained by padding $Q$ with $\left(n-w^{2}\right)$ all-zero columns. Now observe that there is an $\mathbf{a} \in \mathbb{F}^{n}$ such that $f(C \cdot \mathbf{x})=g(\mathbf{x}+\mathbf{a})$; the translation equivalence test in the claim below returns a $\mathbf{c} \in \mathbb{F}^{n}$ such that $f(C \cdot \mathbf{x})=g(\mathbf{x}+\mathbf{c})$. Hence, $f(C \cdot \mathbf{x})=\operatorname{Det}_{w}\left(Q^{\prime} \mathbf{x}+Q^{\prime} \cdot \mathbf{c}\right)$ implying $f(\mathbf{x})=\operatorname{Det}_{w}\left(Q^{\prime} C^{-1} \mathbf{x}+Q^{\prime} \cdot \mathbf{c}\right)$. The algorithm in Theorem 3 returns $B=Q^{\prime} C^{-1}$ and $\mathbf{b}=Q^{\prime} \cdot \mathbf{c}$.

Claim 5.2. Let $f(\mathbf{x})=g(\mathbf{x}+\mathbf{a})$, where $f, g$ are $(n, d)$-polynomials and $\mathbf{a} \in \mathbb{F}^{n}$. There is randomized algorithm that takes blackbox access to $f$ and $g$ and with probability $1-\frac{(n d)^{\circ(1)}}{q}$ computes $a \mathbf{c} \in \mathbb{F}^{n}$ such that $f(\mathbf{x})=g(\mathbf{x}+\mathbf{c})$.
See [Kay12, DdOS14] (also Algorithm 9 and Lemma 2.1 in [KNST17]) for proofs of the claim.
For the rest of this section, set $n=w^{2}$. The equivalence test for $\operatorname{Det}_{w}$ is done in two steps: In the first step, the problem is reduced to the simpler problem of PS-equivalence testing. The second step then solves the PS-equivalence test. A $\left(w^{2}, w\right)$-polynomial $f \in \mathbb{L}[\mathbf{x}]$ is PS-equivalent to Det ${ }_{w}$ if there is a permutation matrix $P$ and a diagonal matrix $S \in \mathrm{GL}\left(w^{2}, \mathbb{L}\right)$ such that $f=\operatorname{Det}_{w}(P S \cdot \mathbf{x})$.

Lemma 5.1 ( [Kay12]). There is a randomized algorithm that takes input blackbox access to $f$, which is PSequivalent to $\operatorname{Det}_{w}$, and with probability $1-\frac{w^{0(1)}}{q}$ outputs a permutation matrix $P$ and a diagonal matrix $S \in \operatorname{GL}\left(w^{2}, \mathbb{L}\right)$ such that $f=\operatorname{Det}_{w}(P S \cdot \mathbf{x})$. The algorithm runs in $(w \log q)^{O(1)}$ time.

It is in the first step where our algorithm differs from (and slightly simplifies) [Kay12]. This reduction to PS-equivalence testing is given in Section 5.2. As in [Kay12], the algorithm uses the structure of the group of symmetries and the Lie algebra of $\mathrm{Det}_{w}$. An estimate of the probability that a random element of the Lie algebra of $g_{\text {Det }_{w v}}$ has all its eigenvalues in $\mathbb{L}$ (Lemma 5.4) is key to the simplification in the first step.

### 5.1 Group of symmetries and Lie algebra of determinant

We state a few well known facts and claims about the Lie algebra and the group of symmetries of Det $_{w}$. Proofs of these can be found in [Kay12,KNST17] and the references therein.

Definition 5.1. The group of symmetries of an $n$-variate polynomial $f$, denoted as $\mathscr{G}_{f}$, consists of matrices $A \in \mathrm{GL}(n, \mathbb{F})$ such that $f(\mathbf{x})=f(A \cdot \mathbf{x})$.
$\operatorname{Det}_{w}(\mathbf{x})$ is the determinant of the symbolic matrix $X=\left(x_{i j}\right)_{i, j \in[w]}$, where $\mathbf{x}=\left\{x_{i j}\right\}_{i, j \in[w]}$. Let $A(X)$ denote the $w \times w$ linear matrix obtained by applying a transformation $A \in \mathbb{F}^{w^{2} \times w^{2}}$ on $\mathbf{x}$.
Fact 1. An $A \in \mathrm{GL}\left(w^{2}, \mathbb{F}\right)$ is in $\mathscr{G}_{\operatorname{Det}_{w}}$ if and only if there are two matrices $S, T \in \mathrm{SL}(w, \mathbb{F})$ such that either $A(X)=S \cdot X \cdot T$ or $A(X)=S \cdot X^{T} \cdot T$.

Definition 5.2. The Lie algebra of a polynomial $f \in \mathbb{F}\left[x_{1}, x_{2}, \ldots, x_{n}\right]$, denoted as $\mathfrak{g}_{f}$, is the set of all $n \times n$ matrices $E=\left(e_{i j}\right)_{i, j \in[n]}$ in $\mathbb{F}^{n \times n}$ satisfying

$$
\sum_{i, j \in[n]} e_{i j} x_{j} \cdot \frac{\partial f}{\partial x_{i}}=0 .
$$

To express the Lie algebra of $\operatorname{Det}_{w}$, order the variables of $\mathbf{x}$ in row major fashion and call them $x_{1}, \ldots, x_{n}$. Let $\mathcal{Z}_{w}$ be the $\mathbb{F}$-linear space of all $w \times w$ traceless matrices over $\mathbb{F}, \mathcal{L}_{\text {row }}$ be the space $\mathcal{Z}_{w} \otimes I_{w}=\left\{Z \otimes I_{w}: Z \in \mathcal{Z}_{w}\right\}$, and $\mathcal{L}_{\text {col }}$ the space $I_{w} \otimes \mathcal{Z}_{w}=\left\{I_{w} \otimes Z: Z \in \mathcal{Z}_{w}\right\}$.
Fact 2. $\mathfrak{g}_{\operatorname{Det}_{w}}=\mathcal{L}_{\text {row }} \oplus \mathcal{L}_{\text {col }}$.
It follows that the dimension of $\mathfrak{g}_{\operatorname{Det}_{w}}$ over $\mathbb{F}$ is $2 w^{2}-2$.
Fact 3. Let $f, g$ be $n$-variate polynomials such that there is an $A \in G L(n, \mathbb{F})$ satisfying $f=g(A \cdot \mathbf{x})$. Then $\mathfrak{g}_{f}=A^{-1} \cdot \mathfrak{g}_{g} \cdot A=\left\{A^{-1} \cdot L \cdot A \mid L \in \mathfrak{g}_{g}\right\}$.
Claim 5.3. There is a randomized algorithm that given blackbox access to a $(n, d)$-polynomial $f$ over $\mathbb{F}$, computes an $\mathbb{F}$-basis of $\mathfrak{g}_{f}$ with probability $1-\frac{(n d)^{\circ(1)}}{q}$. The algorithm runs in $(n d \log q)^{O(1)}$ time.
From Fact 2, it is easy to observe that $\mathfrak{g}_{\text {Det }_{w}}$ contains a diagonal matrix with distinct elements on the diagonal. The next claim can be proved using this observation.
Claim 5.4. Let $L_{1}, \ldots, L_{2 w^{2}-2}$ be an $\mathbb{F}$-basis of $\mathfrak{g}_{\text {Det }_{w^{\prime}}}$ and $L=\sum_{i=1}^{2 w^{2}-2} \alpha_{i} \cdot L_{i}$, where $\alpha_{1}, \ldots, \alpha_{2 w^{2}-2} \in_{r} \mathbb{F}$ are picked independently. Then, the characteristic polynomial of $L$ is square-free with probability $1-\frac{w^{0(1)}}{q}$. The following lemma is the main technical contribution of this section.

Lemma 5.2. Let $L_{1}, \ldots, L_{2 w^{2}-2}$ be an $\mathbb{F}$-basis of $\mathfrak{g}_{\operatorname{Det}_{w}}$, and $L=\sum_{i=1}^{2 w w^{2}-2} \alpha_{i} \cdot L_{i}$, where $\alpha_{1}, \ldots, \alpha_{2 w^{2}-2} \in_{r} \mathbb{F}$ are picked independently. Then, the characteristic polynomial of $L$ is square-free and splits completely over $\mathbb{L}$ with probability at least $\frac{1}{2 w^{2}}$.

Proof. Let $h(y)$ be the characteristic polynomial of $L$. From Claim 5.4, $h$ is square-free with probability $1-\frac{w^{0(1)}}{q}$. From Fact $2, L=L_{1}+L_{2}$ where $L_{1} \in \mathcal{L}_{\text {row }}$ and $L_{2} \in \mathcal{L}_{\text {col }}$. As $L$ is uniformly distributed over $\mathfrak{g}_{\text {Det }}$, so is $L_{1}$ over $\mathcal{L}_{\text {row }}$ and $L_{2}$ over $\mathcal{L}_{\text {col }}$. In other words, if $L_{1}=Z_{1} \otimes I_{w}$ and $L_{2}=I_{w} \otimes Z_{2}$ then $Z_{1}, Z_{2}$ are both uniformly (and independently) distributed over $\mathcal{Z}_{w}$. If the characteristic polynomial of $Z_{1}$ (similarly $Z_{2}$ ) is irreducible over $\mathbb{F}$ then the eigenvalues of $Z_{1}$ (respectively, $Z_{2}$ ) lie in $\mathbb{L}$ and are distinct. If this happens for both $Z_{1}$ and $Z_{2}$ then there are $D_{1}, D_{2} \in G L(w, \mathbb{L})$ such that $D_{1}^{-1} Z_{1} D_{1}$ and $D_{2}^{-1} Z_{2} D_{2}$ are diagonal matrices. This further implies,

$$
\left(D_{1}^{-1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}^{-1}\right) \cdot L \cdot\left(I_{w} \otimes D_{2}\right) \cdot\left(D_{1} \otimes I_{w}\right)
$$

is a diagonal matrix, due to the observation below.
Observation 5.1. For any $M, N \in \overline{\mathbb{F}}^{w \times w},\left(M \otimes I_{w}\right)$ and $\left(I_{w} \otimes N\right)$ commutes. Also, if $M, N \in \operatorname{GL}(w, \overline{\mathbb{F}})$ then $\left(M \otimes I_{w}\right)^{-1}=\left(M^{-1} \otimes I_{w}\right)$ and $\left(I_{w} \otimes N\right)^{-1}=\left(I_{w} \otimes N^{-1}\right)$.

Thus, if we show that the characteristic polynomial of $Z \in_{r} \mathcal{Z}_{w}$ is irreducible with probability $\delta$ then with probability at least $\delta^{2}$ the characteristic polynomial of $L$ splits completely over $\mathbb{L}$. Much like the proof of Claim 5.4, it can be shown that the characteristic polynomial of $Z \in_{r} \mathcal{Z}_{w}$ is square-free with probability $1-\frac{w^{0(1)}}{q}$. Hence, if the characteristic polynomial of $Z \in_{r} \mathcal{Z}_{w}^{\prime}$, where $\mathcal{Z}_{w}^{\prime} \subset \mathcal{Z}_{w}$ consists of matrices with distinct eigenvalues in $\overline{\mathbb{F}}$, is irreducible with probability $\rho$ then $\delta \geq \rho \cdot\left(1-\frac{w^{\circ}(1)}{q}\right)$. Next, we lower bound $\rho$.

Let $\mathcal{P}$ be the set of monic, degree- $w$, square-free polynomials in $\mathbb{F}[y]$ with the coefficient of $y^{w-1}$ equal to zero. Define a map $\phi$ from $\mathcal{Z}_{w}^{\prime}$ to $\mathcal{P}$,

$$
\phi: \quad Z \mapsto \text { characteristic polynomial of } Z .
$$

The map $\phi$ is onto as the companion matrix of $p(y) \in \mathcal{P}$ belongs to its pre-image under $\phi$. Let $\phi^{-1}(p(y))$ be the set of matrices in $\mathcal{Z}_{w}^{\prime}$ that map to $p$.
Claim 5.5. Let $p(y) \in \mathcal{P}$. Then

$$
\frac{\left(q^{w}-1\right) \cdot\left(q^{w}-q\right) \ldots\left(q^{w}-q^{w-1}\right)}{q^{w}} \leq\left|\phi^{-1}(p(y))\right| \leq \frac{\left(q^{w}-1\right) \cdot\left(q^{w}-q\right) \ldots\left(q^{w}-q^{w-1}\right)}{q^{w}\left(1-\frac{w}{q}\right)} .
$$

Proof. Let $C_{p}$ be the companion matrix of $p(y)$. If the characteristic polynomial of a $Z \in \mathcal{Z}_{w}^{\prime}$ equals $p(y)$ then there is an $E \in G L(w, \mathbb{F})$ such that $Z=E \cdot C_{p} \cdot E^{-1}$, as the eigenvalues of $C_{p}$ are distinct in $\overline{\mathbb{F}}$. Moreover, for any $E \in G L(w, \mathbb{F}), E \cdot C_{p} \cdot E^{-1} \in \mathcal{Z}_{w}^{\prime}$ has characteristic polynomial $p(y)$. Hence, $\phi^{-1}(p(y))=\left\{E \cdot C_{p} \cdot E^{-1} \mid E \in G L(w, \mathbb{F})\right\}$. Suppose $E, F \in \mathrm{GL}(w, \mathbb{F})$ such that $F \cdot C_{p} \cdot F^{-1}=E \cdot C_{p} \cdot E^{-1}$. Then $E^{-1} F$ commutes with $C_{p}$. Since $C_{p}$ has distinct eigenvalues in $\overline{\mathbb{F}}$, $E^{-1} F$ can be expressed as a polynomial in $C_{p}$, say $h\left(C_{p}\right)$, of degree at most $(w-1)$ with coefficients from $\mathbb{F}$. Conversely, if $h \in \mathbb{F}[y] \leq(w-1) 37$ and $h\left(C_{p}\right)$ is invertible then $F=E \cdot h\left(C_{p}\right)$ is such that $F \cdot C_{p} \cdot F^{-1}=E \cdot C_{p} \cdot E^{-1}$. As $h_{1}\left(C_{p}\right) \neq h_{2}\left(C_{p}\right)$ for distinct $h_{1}, h_{2} \in \mathbb{F}[y]^{\leq(w-1)}$, we have

$$
\left|\phi^{-1}(p(y))\right|=\frac{|\mathrm{GL}(w, \mathbb{F})|}{\mid\left\{h \in \mathbb{F}[y]: \operatorname{deg}(h) \leq(w-1) \text { and } h\left(C_{p}\right) \in \mathrm{GL}(w, \mathbb{F})\right\} \mid}
$$

[^19]The numerator is exactly $\left(q^{w}-1\right) \cdot\left(q^{w}-q\right) \ldots\left(q^{w}-q^{w-1}\right)$, and the denominator is trivially upper bounded by $q^{w}$. A lower bound on the denominator can be worked out as follows: Let $\lambda_{1}, \ldots, \lambda_{w} \in \overline{\mathbb{F}}$ be the distinct eigenvalues of $C_{p}$. If $h(y)=a_{w-1} y^{w-1}+a_{w-2} y^{w-2}+\ldots+a_{0} \in \mathbb{F}[y]$, then $h\left(\lambda_{1}\right), \ldots, h\left(\lambda_{w}\right)$ are the eigenvalues of $h\left(C_{p}\right)$. Observe that

$$
\begin{aligned}
& \operatorname{Pr}_{h \in_{r} \mathbb{F}[y]^{\leq(w-1)}} \quad\left\{h\left(\lambda_{i}\right)=0, \text { for some fixed } i \in[w]\right\} \leq \frac{1}{q^{\prime}} \\
\Rightarrow & \operatorname{Pr}_{h \in_{r} \mathbb{F}[y]^{\leq(w-1)}} \quad\left\{h\left(\lambda_{i}\right)=0, \text { for any } i \in[w]\right\} \leq \frac{w}{q}, \\
\Rightarrow & \operatorname{Pr}_{h \in \in_{r} \mathbb{F}[y]^{\leq(w-1)}} \quad\left\{h\left(C_{p}\right) \in \mathrm{GL}(w, \mathbb{F})\right\} \geq 1-\frac{w}{q} .
\end{aligned}
$$

Hence, the denominator is lower bounded by $q^{w}\left(1-\frac{w}{q}\right)$.
Let $\rho_{p}=\frac{\left|\phi^{-1}(p(y))\right|}{\left|\mathcal{Z}_{w}^{\prime}\right|}$, the probability that $p(y)$ is the characteristic polynomial of $Z \in_{r} \mathcal{Z}_{w}^{\prime}$. From Claim 5.5, it follows that

$$
\left|\mathcal{Z}_{w}^{\prime}\right| \leq \frac{\left(q^{w}-1\right) \cdot\left(q^{w}-q\right) \ldots\left(q^{w}-q^{w-1}\right)}{q^{w}\left(1-\frac{w}{q}\right)} \cdot|\mathcal{P}| \Rightarrow 1-\frac{w}{q} \leq \rho_{p} \cdot|\mathcal{P}|
$$

We show in the next claim that a $p \in_{r} \mathcal{P}$ is irreducible over $\mathbb{F}$ with probability at least $\frac{1}{\bar{w}}\left(1-\frac{2}{q^{w / 2}}\right)$, implying the characteristic polynomial of $Z \in_{r} \mathcal{Z}_{w}^{\prime}$ is irreducible over $\mathbb{F}$ with probability $\rho \geq$ $\frac{1}{w}\left(1-\frac{2}{q^{w / 2}}\right)\left(1-\frac{w}{q}\right)$. Therefore, the probability that the characteristic polynomial of $Z \in_{r} \mathcal{Z}_{w}$ is irreducible over $\mathbb{F}$ is $\delta \geq \frac{1}{w}\left(1-\frac{2}{q^{w / 2}}\right)\left(1-\frac{w}{q}\right)\left(1-\frac{w^{O(1)}}{q}\right)$. As $q \geq w^{7}$, the probability that the characteristic polynomial of $L \in_{r} \mathfrak{g}_{\text {Det }_{w}}$ splits completely over $\mathbb{L}$ is at least $\delta^{2} \geq \frac{1}{2 w^{2}}$.
Claim 5.6. A polynomial $p \in_{r} \mathcal{P}$ is irreducible over $\mathbb{F}$ with probability at least $\frac{1}{w}\left(1-\frac{2}{q^{w / 2}}\right)$.
Proof. Let $\mathcal{F}$ be the set of monic, degree- $w$, square-free polynomials in $\mathbb{F}[y]$. The difference between $\mathcal{F}$ and $\mathcal{P}$ is that a polynomial in $\mathcal{P}$ additionally has coefficient of $y^{w-1}$ equal to zero. We argue in the next paragraph that the fraction of $\mathbb{F}$-irreducible polynomials in $\mathcal{F}$ and in $\mathcal{P}$ are the same. As irreducible polynomials are square-free, the number of irreducible polynomials in $\mathcal{F}$ is at least $\frac{q^{w}-2 q^{w / 2}}{w}$ [vzGG03]. Hence, the fraction of irreducible polynomials in $\mathcal{F}$ is at least $\frac{1}{w}\left(1-\frac{2}{q^{w / 2}}\right)$.

Define a map $\Psi$ from $\mathcal{F}$ to $\mathcal{P}$ as follows: For a $u(y)=y^{w}+a_{w-1} y^{w-1}+\ldots+a_{0} \in \mathcal{F}$, define $\Psi(u)=u\left(y-\frac{a_{w-1}}{w}\right)$. Observe that the coefficient of $y^{w-1}$ in $\Psi(u)$ is zero. It is also an easy exercise to show that $\Psi\left(u_{1}\right)=\Psi\left(u_{2}\right)$ if and only if there exists an $a \in \mathbb{F}$ such that $u_{1}(y)=u_{2}(y+a)$. As $u(y)$ is irreducible over $\mathbb{F}$ if and only if $u(y+a)$ is irreducible over $\mathbb{F}$, for $a \in \mathbb{F}$, the fraction of $\mathbb{F}$-irreducible polynomials in $\mathcal{F}$ is the same as that in $\mathcal{P}$.

This completes the proof of Lemma 5.2.

### 5.2 Reduction to PS-equivalence testing

Algorithm 7 gives a reduction to PS-equivalence testing for $\operatorname{Det}_{w}$. Suppose the input to the algorithm is a blackbox access to $f=\operatorname{Det}_{w}(A \cdot \mathbf{x})$, where $A \in \mathrm{GL}\left(w^{2}, \mathbb{F}\right)$. We argue the correctness of the algorithm by tracing its steps:

```
Algorithm 7 Reduction to PS-equivalence
    INPUT: Blackbox access to a \(\left(w^{2}, w\right)\)-polynomial \(f \in \mathbb{F}[\mathbf{x}]\) that is equivalent to \(\operatorname{Det}_{w}\) over \(\mathbb{F}\).
    OUTPUT: A \(D \in G L\left(w^{2}, \mathbb{L}\right)\) such that \(f(D \cdot \mathbf{x})\) is PS-equivalent to \(\operatorname{Det}_{w}\) over \(\mathbb{L}\).
    Compute an \(\mathbb{F}\)-basis of \(\mathfrak{g}_{f}\). Let \(\left\{F_{1}, F_{2}, \ldots F_{2 w^{2}-2}\right\}\) be the basis. Set \(j=1\).
    for \(j=1\) to \(w^{3} \log q\) do
        Pick \(\alpha_{1}, \ldots, \alpha_{2 w^{2}-2} \in r \mathbb{F}\) independently. Set \(F=\sum_{i \in\left[2 w^{2}-2\right]} \alpha_{i} \cdot F_{i}\).
        Compute the characteristic polynomial \(h\) of \(F\). Factorize \(h\) into irreducible factors over \(\mathbb{L}\).
        if \(h\) is square-free and splits completely over \(\mathbb{L}\) then
            Use the roots of \(h\) to compute a \(D \in \mathrm{GL}\left(w^{2}, \mathbb{L}\right)\) such that \(D^{-1} \cdot F \cdot D\) is diagonal.
            Exit loop.
        else
            Set \(j=j+1\).
        end if
    end for
    if No \(D\) found at step 7 in the loop then
        Output 'Failed'.
    else
        Output D.
    end if
```

Step 1: An $\mathbb{F}$-basis of $\mathfrak{g}_{f}$ can be computed efficiently using Claim 5.3.
Step 3-12: At step 4 an element $F$ of $\mathfrak{g}_{f}$ is chosen uniformly at random. By Fact $3, F=A^{-1} \cdot L \cdot A$, where $L$ is a random element of $\mathfrak{g}_{\text {Det }_{w}}$. Lemma 5.2 implies, in every iteration of the loop, $h$ (at step 5 ) is square-free and splits completely over $\mathbb{L}$ with probability at least $\frac{1}{2 w^{2}}$. Since the loop has $w^{3} \log q$ iterations, the algorithm finds an $h$ that is square-free and splits completely over $\mathbb{L}$, with probability at least $1-\frac{1}{q}$. Assume that the algorithm succeeds in finding such an $h$, and suppose $\lambda_{1}, \ldots, \lambda_{w^{2}} \in \mathbb{L}$ are the distinct roots of $h$. The algorithm finds a $D$ in step 7 by picking a random solution of the linear system obtained from the relation $F \cdot D=D \cdot \operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{w w^{2}}\right)$ treating the entries of $D$ as formal variables. We argue next that $f(D \cdot \mathbf{x})$ is PS-equivalent to $\operatorname{Det}_{w}$ over $\mathbb{L}$.

By Fact $2, L=L_{1}+L_{2}$ where $L_{1} \in \mathcal{L}_{\text {row }}$ and $L_{2} \in \mathcal{L}_{\text {col }}$. In other words, there are $Z_{1}, Z_{2} \in \mathcal{Z}_{w}$ such that $L_{1}=Z_{1} \otimes I_{w}$ and $L_{2}=I_{w} \otimes Z_{2}$. It is easy to verify, if $L$ has distinct eigenvalues then so do $Z_{1}$ and $Z_{2}$. Hence, there are $D_{1}, D_{2} \in G L(w, \overline{\mathbb{F}})$ such that $D_{1} Z_{1} D_{1}^{-1}$ and $D_{2} Z_{2} D_{2}^{-1}$ are both diagonal, implying

$$
M \stackrel{\text { def }}{=}\left(D_{1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}\right) \cdot L \cdot\left(D_{1}^{-1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}^{-1}\right)
$$

is diagonal (by Observation 5.1) with distinct diagonal entries. Also,

$$
\begin{aligned}
D^{-1} \cdot F \cdot D & =(A D)^{-1} \cdot L \cdot(A D) \\
& =\left(\left(D_{1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}\right) \cdot A D\right)^{-1} \cdot M \cdot\left(\left(D_{1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}\right) \cdot A D\right)
\end{aligned}
$$

As both $D^{-1} \cdot F \cdot D$ and $M$ are diagonal matrices with distinct diagonal entries, it must be that

$$
\left(D_{1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}\right) \cdot A D=P \cdot S,
$$

where $P$ is a permutation matrix and $S \in \mathrm{GL}\left(w^{2}, \overline{\mathbb{F}}\right)$ is a diagonal matrix. Now observe that $\operatorname{Det}_{w}\left(\left(D_{1} \otimes I_{w}\right) \cdot \mathbf{x}\right)=\beta \cdot \operatorname{Det}_{w}(\mathbf{x})$ and $\operatorname{Det}_{w}\left(\left(I_{w} \otimes D_{2}\right) \cdot \mathbf{x}\right)=\gamma \cdot \operatorname{Det}_{w}(\mathbf{x})$, for $\beta, \gamma \in \overline{\mathbb{F}} \backslash\{0\}$. Hence,

$$
\begin{aligned}
\operatorname{Det}_{w}(P \cdot S \cdot \mathbf{x}) & =\operatorname{Det}_{w}\left(\left(D_{1} \otimes I_{w}\right) \cdot\left(I_{w} \otimes D_{2}\right) \cdot A D \cdot \mathbf{x}\right) \\
& =\beta \gamma \cdot \operatorname{Det}_{w}(A D \cdot \mathbf{x}) \\
& =\beta \gamma \cdot f(D \cdot \mathbf{x}) \\
\Rightarrow f(D \cdot \mathbf{x}) & =\operatorname{Det}_{w}\left(P \cdot S^{\prime} \cdot \mathbf{x}\right)
\end{aligned}
$$

where $S^{\prime} \in \mathrm{GL}\left(w^{2}, \overline{\mathbb{F}}\right)$ is also diagonal. Therefore, $f(D \cdot \mathbf{x})$ is $P S$-equivalent to $\operatorname{Det}_{w}$ over $\overline{\mathbb{F}}$. As $f(D \cdot \mathbf{x}) \in \mathbb{L}[\mathbf{x}]$, it is a simple exercise to show that $f(D \cdot \mathbf{x})$ must be $P S$-equivalent to $\operatorname{Det}_{w}$ over $\mathbb{L}$.

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## A Proof of two claims in Section 3

Claim 3.1 (restated): With probability $1-(w d n)^{-\Omega(1)}$, any subset of $w$ vectors in any of the sets $\left\{\mathbf{u}_{1}, \mathbf{u}_{2}\right.$, $\left.\ldots, \mathbf{u}_{w+1}\right\},\left\{\mathbf{v}_{1}, \mathbf{v}_{2}, \ldots, \mathbf{v}_{w+1}\right\},\left\{\mathbf{w}_{1}, \mathbf{w}_{2}, \ldots, \mathbf{w}_{w+1}\right\}$, or $\left\{\mathbf{s}_{1}, \mathbf{s}_{2}, \ldots, \mathbf{s}_{w+1}\right\}$ are L-linearly independent.

Proof. From Observation 3.3, for the sets $\left\{\mathbf{v}_{1}, \mathbf{v}_{2}, \ldots, \mathbf{v}_{w+1}\right\}$ and $\left\{\mathbf{s}_{1}, \mathbf{s}_{2}, \ldots, \mathbf{s}_{w+1}\right\}$ it is sufficient to show that any $w$ columns of the $w \times(w+1)$ matrices $\left(N_{1 i}\left(\mathbf{a}_{j}\right)\right)_{i \in[w], j \in[w+1]}$ and $\left(N_{1 i}\left(\mathbf{b}_{j}\right)\right)_{i \in[w], j \in[w+1]}$ are $\mathbb{L}$-linearly independent with high probability. As the cofactors $N_{11}, \ldots, N_{1 w}$ are $\mathbb{L}$-linearly independent, the above follows from Claim 2.2. For the sets $\left\{\mathbf{u}_{1}, \mathbf{u}_{2}, \ldots, \mathbf{u}_{w+1}\right\}$ and $\left\{\mathbf{w}_{1}, \mathbf{w}_{2}, \ldots, \mathbf{w}_{w+1}\right\}$, it follows from Equation 2 that there are $\lambda_{k}, \rho_{k} \in \mathbb{L}^{\times}$such that $D \cdot \mathbf{v}_{k}=\lambda_{k} \mathbf{u}_{k}$ and $D \cdot \mathbf{s}_{k}=\rho_{k} \mathbf{w}_{k}$ for all $k \in[w+1]$. Since $D$ is invertible, the claim follows for these two sets as well.
Claim 3.2 (restated): If $E=Q_{1} \cdots Q_{\ell}$ is a random $(w, \ell, m)$-matrix product over $\mathbb{F}$, where $w^{2}+1 \leq$ $m \leq n$ and $\ell \leq d$, then the entries of $E$ are $\mathbb{F}$-linearly independent with probability $1-(w d n)^{-\Omega(1)}$.

Proof. Treat the coefficients of the linear forms in $Q_{1}, Q_{2}, \ldots, Q_{\ell}$ as distinct formal variables. In particular

$$
Q_{k}=\sum_{i=1}^{m} U_{i}^{(k)} x_{i} \text { for } k \in[\ell],
$$

where the $U_{i}^{(k)}$ 's are $w \times w$ matrices and the entries of these matrices are distinct $\mathbf{u}$-variables. The entries of the matrix product $E$ are polynomials in the $\mathbf{x}$-variables over $\mathbb{F}(\mathbf{u})$. If we show the $w^{2}$ entries of $E$ are $\mathbb{F}(\mathbf{u})$-linearly independent then an application of Schwartz-Zippel lemma implies the statement of the claim. On the other hand, to show that the entries of $E$ are $\mathbb{F}(\mathbf{u})$-linearly independent, it is sufficient to show that the entries are $\mathbb{F}$-linearly independent under a setting of the $\mathbf{u}$-variables to $\mathbb{F}$ elements. Consider such a setting: For every $k \in[\ell] \backslash\{1\}$, let $U_{w^{2}+1}^{(k)}=I_{w}$ and $U_{i}^{(k)}=0$ for all $i \in[m] \backslash\left\{w^{2}+1\right\}$. Let $U_{i}^{(1)}=0$ for all $i \geq w^{2}+1$ and set $U_{1}^{(1)}, \ldots, U_{w^{2}}^{(1)}$ in a way so that the linear forms in $\sum_{i=1}^{w^{2}} U_{i}^{(1)} x_{i}$ are $\mathbb{F}$-linearly independent. It is straightforward to check that the entries of $E$ under this setting are $\mathbb{F}$-linearly independent.


[^0]:    ${ }^{1}$ It is worth noting that an average-case reconstruction algorithm for ABPs does not necessarily subsume a result on average-case reconstruction of formulas as the distributions of the inputs may be incomparable.
    ${ }^{2}$ For circuit classes whose known lower bound proofs do not fit in the natural proof framework, the situation is less clear. Examples of such classes are ACC ${ }^{0}$ [Wil14] and monotone circuits [Raz85]. A hardness result for polynomial-time learning of monotone circuits is known assuming the existence of one-way functions [DLM ${ }^{+}$08].

[^1]:    ${ }^{3}$ Unlike PAC learning, in the algebraic setting we need to reconstruct a circuit that computes the input polynomial exactly instead of approximately (as two polynomial functions differ at too many points). If we insist on exact learning in the Boolean setting (which is closely related to the compression problem) then the best known output circuit size for $\mathrm{AC}^{0}$ and $\mathrm{AC}^{0}[p]$ functions is exponential in the number of variables [CKK ${ }^{+} 15$, Sri15, CIKK16]. On the other hand, reconstruction algorithms have the power of making membership queries.
    ${ }^{4}$ Even in the Boolean setting, similar average-case relaxations of learning problems have been studied, particularly for DNFs [LSW06,JLSW08].

[^2]:    ${ }^{5}$ A more general way to define an ABP (in Definition 1.1) is to consider matrices of varying dimensions, i.e. the $i$-th matrix has dimension $w_{i} \times w_{i+1}$, and $w_{1}=w_{d+1}=1$. In this case, size of the ABP is the quantity $\sum_{i=1}^{d+1} w_{i}$. Equivalently, an ABP can be defined as a layered directed acyclic graph, in which case size is the number of nodes in the graph.
    ${ }^{6} \mathrm{An} \mathrm{ABP} X_{1} \cdot X_{2} \ldots X_{d}$ is homogeneous, if every entry in every partial product $X_{1} \cdot X_{2} \ldots X_{i}$ is a homogeneous polynomial.
    ${ }^{7} \mathrm{~A}$ trivial brute-force algorithm to reconstruct a $(w, d, n)$ - ABP over $\mathbb{F}_{q}$ takes time $q^{\Theta\left(w^{2} d n\right)}$. By 'nontrivial' reconstruction, we mean an algorithm that takes time exponentially better than the trivial complexity. Note that we can interpolate a polynomial computed by a $(w, d, n)$ - ABP in $\left(d^{n} \log q\right)^{O(1)}$ time, but knowing the coefficients of the polynomial does not give us any immediate information about the ( $w, d, n$ )-ABP that computes it. Hence, if we want a $(w, d, n)$-ABP representation for the input polynomial then even a $\left(d^{n} \log q\right)^{O(1)}$ time reconstruction algorithm is nontrivial as $d^{n} \ll q^{\Theta\left(w^{2} d n\right)}$. The complexity of our algorithm is $\left(d^{w^{3}} n \log q\right)^{O(1)}$ which is exponentially better than the trivial complexity $q^{\Theta\left(w^{2} d n\right)}$ for $w=O(n)$.

[^3]:    ${ }^{8}$ The probability is taken over the input distribution and the random bits used by the algorithm, if it is randomized.
    ${ }^{9} \mathbb{L}$ can constructed from a basis of $\mathbb{F}_{q}$ using a randomized algorithm running in $(w \log q)^{O(1)}$ time [vzGG03].
    ${ }^{10}$ We thank Rohit Gurjar for showing us a similar example.

[^4]:    ${ }^{11}$ The algorithm in [KNST17] works over both $\mathbb{Q}$ and $\mathbb{F}_{q}$, whereas ours is over $\mathbb{F}_{q}$.
    ${ }^{12}$ [KNST17] has running time polynomial in all the relevant parameters, and it also works if the width $w$ is varying along the ABP .

[^5]:    ${ }^{13} \operatorname{det}\left(X_{i}\right)$ is the determinant of the $w \times w$ matrix $X_{i}$.

[^6]:    ${ }^{14}$ Given in Section 5.
    ${ }^{15}$ i.e., if $C_{i} \cdot X_{i} \cdot D_{i}=C_{i}^{\prime} \cdot X_{i} \cdot D_{i}^{\prime}$, where $X_{i}$ is a full rank matrix, then $C_{i}^{\prime}=\alpha C_{i}$ and $D_{i}^{\prime}=\alpha^{-1} D_{i}$ for some $\alpha \in \mathbb{L}^{\times}$
    ${ }^{16}$ which also gives the uniqueness of factorization mentioned in the remark after Theorem 1

[^7]:    ${ }^{17}$ For a field $\mathbb{K} \supseteq \mathbb{F}$, we say $f$ is zero modulo a $\mathbb{K}$-linear space $\mathcal{X}=\operatorname{span}_{\mathbb{K}}\left\{l_{1}, \ldots, l_{w}\right\}$, where $l_{i}$ 's are linear forms in $\mathbb{K}[\mathbf{x}]$, if $f$ is in the ideal of $\mathbb{K}[\mathbf{x}]$ generated by $\left\{l_{1}, \ldots, l_{w}\right\}$. This is also denoted by $f=0 \bmod \left\langle l_{1}, \ldots, l_{w}\right\rangle$.
    ${ }^{18}$ the algebraic closure of $\mathbb{F}$
    ${ }^{19}$ The matrix $X_{i}^{\prime}$ with the $\mathbf{y}$ and $\mathbf{z}$ variables in its linear forms substituted to zero.
    ${ }^{20}$ by identifying the $1 \times 1$ matrix of the R.H.S with the entry of the matrix

[^8]:    ${ }^{21}$ over the randomness of the input $f$
    ${ }^{22}$ This lemma makes our reduction to PS-equivalence simpler than [Kay12], enabling the equivalence test to work over finite fields.
    ${ }^{23}$ In [Kay12], a basis of the centralizer of $F$ in $\mathfrak{g}_{g}$ is computed first and then a $D \in \mathrm{GL}\left(w^{2}, \mathbb{C}\right)$ is obtained that simulaneously diagonalizes this basis.

[^9]:    ${ }^{24} \mathrm{~A}$ similar result, but for homogeneous $g_{1}, \ldots, g_{m}$, follows from [Laz01].

[^10]:    ${ }^{25}$ For $t=d, R_{i}=X_{i}$ for all $i \in[d]$.

[^11]:    ${ }^{26} \operatorname{det}(Z)$ being multilinear, there is an injective ring homomorphism from $\mathbb{L}[\mathbf{x}] /(\operatorname{det}(Z))$ to $\mathbb{L}(\mathbf{x})$ via a simple substitution map taking a variable to a rational function.

[^12]:    ${ }^{27}$ One way of seeing this is that $D^{-1} M$ is already fixed before $\mathbf{b}_{w+1}$ is chosen.

[^13]:    ${ }^{28} z_{11}$ does not even appear in $h_{l k}$.

[^14]:    ${ }^{29}$ As the individual degrees of the variables in $p$ are bounded by $d$, we only need $|\mathbb{F}|>d$ to carry out this interpolation.
    ${ }^{30}$ Otherwise, we will have $f(A \cdot \mathbf{x})=0$ modulo less than $w$ linear forms in $\mathbb{K}[\mathbf{y} \uplus \mathbf{z} \uplus \mathbf{r}]$.

[^15]:    ${ }^{31}$ See [CLO07]. Equivalently, think of the set of linear forms obtained from a reduced row echelon form of the coefficient matrix of $l_{1}, \ldots, l_{k}$.
    ${ }^{32}$ via row and column operations on $X_{1}^{\prime}$ and $X_{3}^{\prime}$, respectively

[^16]:    ${ }^{33}$ which follows from condition (**)

[^17]:    ${ }^{34}$ by picking an arbitrarily fixed ordering among the linear forms

[^18]:    ${ }^{35}$ by identifying the $1 \times 1$ matrix of the R.H.S with the entry of the matrix
    ${ }^{36}$ For instance, the coefficients of $y_{s}$ in the linear forms in $T_{2}(s, *), T_{3}, \ldots, T_{d-2}$ get computed repeatedly at step 9 for every value of $t \in[w]$ - uniqueness ensures that we always get the same values for these coefficients.

[^19]:    ${ }^{37}$ the set of polynomials in $\mathbb{F}[y]$ of degree at most $w-1$.

