# An Efficient Sparse Regularity Concept 

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#### Abstract

Let $\mathbf{A}$ be a $0 / 1$ matrix of size $m \times n$, and let $p$ be the density of $\mathbf{A}$ (i.e., the number of ones divided by $m \cdot n$ ). We show that A can be approximated in the cut norm within $\varepsilon \cdot m n p$ by a sum of cut matrices (of rank 1 ), where the number of summands is independent of the size $m \cdot n$ of $\mathbf{A}$, provided that $\mathbf{A}$ satisfies a certain boundedness condition. This decomposition can be computed in polynomial time. This result extends the work of Frieze and Kannan [16] to sparse matrices. As an application, we obtain efficient $1-\varepsilon$ approximation algorithms for "bounded" instances of MAX CSP problems. Key words: approximation algorithms, regularity lemma, matrix decomposition, cut norm, random discrete structures.


## 1 Introduction and Results

For many fundamental optimization problems there are NP-hardness of approximation results known, showing that not only is it NP-hard to compute the optimum exactly, but even to approximate the optimum within a factor bounded away from 1. For instance, in the MAX $k$-SAT problem it is NP-hard to achieve an approximation ratio better than $1-2^{-k}$ [20]. Furthermore, it is NP-hard to approximate MAX CUT within better than $16 / 17 \approx 0.94118$ [20, 25] (which can be tightened to $\approx 0.87856$ under a stronger hypothesis [21]).

Frieze and Kannan [16] showed that the situation is much better for dense problem instances. For example, if $G=(V, E)$ is a graph on $n$ vertices of density $p=2 n^{-2}|E|$, then its MAX CUT can be approximated within a factor of $1-\varepsilon$ in time $\operatorname{poly}\left(\exp \left((\varepsilon p)^{-2}\right) \cdot n\right)$. Hence, if $p>\delta$ for some fixed number $\delta>0$, then this algorithm has a polynomial running time. Similarly, if $F$ is a $k$-SAT formula with at least $\delta 2^{k}\binom{n}{k}$ clauses (i.e., at least a constant fraction of all possible clauses is present), then the maximum number of simultaneously satisfiable clauses can be approximated within $1-\varepsilon$ in polynomial time for any fixed $\varepsilon>0$.

The key ingredient in [16] is an algorithm for approximating a dense matrix $\mathbf{A}$ by a sum of a bounded number of "cut matrices". Applied to the adjacency matrix of a graph, this yields the aforementioned algorithm for MAX CUT.

[^0]Moreover, an extension of this matrix algorithm to $k$-dimensional tensors yields the approximation algorithms for
 $m \times n$, and let $0 \leq p \leq 1$ be the density of $\mathbf{A}$, i.e., the number of ones in $\mathbf{A}$ divided by $m \cdot n$. A cut matrix is a matrix $\mathbf{D}$ such that there are sets $S \subset[m], T \subset[n]$ and a number $d$ such that the entry $\mathbf{D}_{i j}$ is equal to $d$ if $(i, j) \in S \times T$ and 0 otherwise. We denote such a matrix by $\mathbf{D}=\operatorname{CUT}(d, S, T)$, and observe that cut matrices have rank one. The cut norm of a $m \times n$ matrix $\mathbf{M}=\left(\mathbf{M}_{i j}\right)_{i \in[m], j \in[n]}$ is

$$
\|\mathbf{M}\|_{\square}=\max _{S \subset[m], T \subset[n]}|\mathbf{M}(S, T)|, \quad \text { where } \mathbf{M}(S, T)=\sum_{(s, t) \in S \times T} \mathbf{M}_{s t} .
$$

Frieze and Kannan proved that for any $\mathbf{A}$ and any $\varepsilon>0$ there exist cut matrices $\mathbf{D}_{1}, \ldots, \mathbf{D}_{s}$ such that

$$
\left\|\mathbf{A}-\left(\mathbf{D}_{1}+\cdots+\mathbf{D}_{s}\right)\right\|_{\square}<\varepsilon \cdot m n
$$

where $s \leq c \varepsilon^{-2}$ for a constant $c>0$. Indeed, such a decomposition can be computed in time $\varepsilon^{-2} \cdot \operatorname{poly}(m n)$ (or even in "constant" expected time $O\left(\varepsilon^{-2} \cdot \operatorname{polylog}(1 / \varepsilon)\right)$ by sampling). Hence, if $p \geq \delta$ for some fixed $\delta>0$, i.e., if $\mathbf{A}$ is a dense matrix, then setting $\varepsilon^{\prime}=\varepsilon p$ we can use this algorithm to find a decomposition of $\mathbf{A}$ within $\varepsilon\|A\|_{\square}=\varepsilon \cdot m n p$ efficiently by a sum of at most $c \varepsilon^{\prime-2}=c(\varepsilon p)^{-2} \leq c(\varepsilon \delta)^{-2}$ cut matrices. The crucial point here is that the number of cut matrices is bounded independently of the size $m \cdot n$ of $\mathbf{A}$.

The goal of the present paper is to extend this result to sparse matrices, where the density $p$ of $\mathbf{A}$ is no longer bounded below by a fixed number. Thus, in asymptotic terms, we are interested in $p=o(1)$ as $m, n \rightarrow \infty$. Clearly, in this case the bound $c(\varepsilon p)^{-2}$ on the number of cut matrices in the decomosition guaranteed by [16] is no longer "constant", but grows with the size $m \cdot n$ of $\mathbf{A}$. Of course, we cannot expect to obtain the same results in the sparse as in the dense case for arbitrary sparse matrices; for in the light of the aforementioned hardness results this would imply $\mathrm{P}=\mathrm{NP}$. Hence, our main result is that even in the sparse case a $0 / 1$ matrix $\mathbf{A}$ (or, more generally, a $k$-dimensional tensor) can be approximated in the cut norm by a sum of cut matrices with a number of summands independent of $m, n$, and $p$, provided that $\mathbf{A}$ satisfies a certain boundedness condition. This condition basically requires that $\mathbf{A}$ does not feature relatively large, extraordinarily dense spots. In addition, we shall use these decomposition results to obtain ( $1-\varepsilon$ )-approximation algorithms for instances of MAX CSP problems that have a suitable boundedness property. As we shall see, in a sense these results mediate between the "average" and the worst case analysis of algorithms.

### 1.0.1 Outline.

In this section we state our results and discuss related work. Section 2 contains a few preliminaries, and in Section 3 we present the algorithms and their analyses for decomposing matrices and graphs. Further, in Section 4 we deal with $k$-dimensional tensors. Then, in Section 5 we apply the tensor algorithm to approximate MAX CSP problems. Finally, Section 6 contains a few examples, which link our results to the "average case" analysis of algorithms.

### 1.1 Approximating $0 / 1$ matrices

Let $\mathbf{A}$ be a $0 / 1$ matrix of size $m \times n$ and density $p$. Given $C, \gamma>0$, we say that $\mathbf{A}$ is $(C, \gamma)$-bounded if for any two sets $S \subset[m]$ and $T \subset[n]$ of sizes $|S| \geq \gamma m,|T| \geq \gamma n$ we have

$$
\begin{equation*}
\mathbf{A}(S, T)=\sum_{(s, t) \in S \times T} \mathbf{A}_{s t} \leq C \cdot|S| \cdot|T| \cdot p \tag{1}
\end{equation*}
$$

In words, for any two sufficiently large sets $S, T$ the number $\mathbf{A}(S, T)$ of ones in the square $S \times T$ must not exceed the number $|S| \cdot|T| \cdot p$ that we would expect if $S, T$ were random sets by more than a factor of $C$.

Theorem 1 There is an algorithm ApxMatrix, absolute constants $\zeta, \zeta^{\prime}>0$, and a polynomial $\Pi$ such that the following holds. Suppose that $0<\varepsilon<\frac{1}{2}, C>1$. Let

$$
\begin{equation*}
\kappa=\frac{\zeta C^{2}}{\varepsilon^{2}} \quad \text { and } \quad \gamma=\gamma(\epsilon, C)=\frac{\zeta^{\prime} \epsilon}{2^{10 \kappa} C} \tag{2}
\end{equation*}
$$

If $\mathbf{A}$ is a $(C, \gamma)$-bounded 0/1 matrix, then in time $\kappa \cdot \Pi(m \cdot n)$, ApxMatrix $(\mathbf{A}, C, \varepsilon)$ outputs cut matrices $\mathbf{D}_{1}, \ldots, \mathbf{D}_{s}$ such that $s \leq \kappa$ and $\left\|\mathbf{A}-\left(\mathbf{D}_{1}+\cdots+\mathbf{D}_{s}\right)\right\|_{\square} \leq \epsilon\|\mathbf{A}\|_{\square}$.

We emphasize that the upper bound $\kappa$ on the number of cut matrices depends only on $C$ and $\varepsilon$, but not on the size of A or the density $p$. Moreover, while for the sake of similicity we assume that ApxMatrix is given the boundedness parameter $C$ as an input, this can easily be avoided by performing a binary search (details omitted).

Given the $0 / 1$ matrix $\mathbf{A}$ and partitions $\mathcal{S}$ of $[m]$ and $\mathcal{T}$ of $[n]$, we define a matrix $\mathbf{A}_{\mathcal{S} \times \mathcal{T}}$ as follows. If $s \in S \in \mathcal{S}$ and $t \in T \in \mathcal{T}$, then the corresponding entry $\left(\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right)_{s, t}$ equals $|S|^{-1}|T|^{-1} \mathbf{A}(S, T)$. Hence, on each square $S \times T$ the matrix $\mathbf{A}_{\mathcal{S} \times \mathcal{T}}$ is constant, and the value it takes is just the average of $\mathbf{A}$ over that square.

Corollary 1 There is an algorithm PartMatrix and a polynomial $\Pi$ that satisfy the following. Suppose that $\varepsilon, C>$ 0 , let $\kappa, \gamma$ be as in (2), and assume that $\mathbf{A}$ is a $(C, \gamma)$-bounded $0 / 1$ matrix of size $m \times n$. Then in time $2^{\kappa} \cdot \Pi(m \cdot n)$ PartMatrix $(\mathbf{A}, C, \varepsilon)$ computes partitions $\mathcal{S}$ of $[m]$ and $\mathcal{T}$ of $[n]$ such that $\left\|\mathbf{A}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \leq 2 \varepsilon\|\mathbf{A}\|_{\square}$. The number of classes in each partition $\mathcal{S}, \mathcal{T}$ is at most $2^{\kappa}$.

### 1.2 Weak regular partitions of graphs

Let $G=(V, E)$ be a graph on $n$ vertices, and let $0 \leq p \leq 1$ be such that $|E|=n^{2} p / 2$; we refer to $p$ as the density of $G$. Moreover, we assume that $V=[n]$. In addition, let $\mathbf{A}=\mathbf{A}(G)$ be the adjacency matrix of $G$. Then we say that $G$ is $(C, \gamma)$-bounded if $\mathbf{A}$ has this property. Thus, if $G$ is $(C, \gamma)$-bounded, then for any two sets $S, T \subset V$ of size at least $\gamma n$ we have $e_{G}(S, T) \leq C \gamma|S||T| p$, where $e_{G}(S, T)$ is the number of $S$ - $T$-edges in $G$.

We call a partition $\mathcal{V}$ of $V$ a weak $\varepsilon$-regular partition of $G$ if $\left\|\mathbf{A}-\mathbf{A}_{\mathcal{V} \times \mathcal{V}}\right\|_{\square} \leq \varepsilon\|\mathbf{A}\|_{\square}=2 \varepsilon|E|$. Hence, if, for instance, $S, T \subset V$ are disjoint sets of vertices, then the number $\mathbf{A}(S, T)$ of $S$ - $T$-edges is within $2 \varepsilon|E|$ of $\mathbf{A}_{\mathcal{V} \times \mathcal{V}}(S, T)$. As we shall see below, this definition is related to the notion of regular partitions introduced by Szemerédi.

Corollary 2 There is an algorithm WeakPartition and a polynomial $\Pi$ that satisfy the following. Suppose that $C>1,0<\varepsilon<\frac{1}{2}$, and let $\kappa, \gamma>0$ be as in (2), and let $G=(V, E)$ be a $(C, \gamma)$-bounded graph on $n$ vertices. Then WeakPartition $(G, C, \varepsilon)$ computes a weak $4 \varepsilon$-regular partition of $G$ in time $2^{2 \kappa} \cdot \Pi(n)$. This partition has at most $2^{2 \kappa}$ classes.

### 1.3 Approximating $k$-dimensional $0 / 1$ tensors

A $k$-dimensional tensor is a map $\mathbf{M}: R_{1} \times R_{2} \times \cdots \times R_{k} \rightarrow \mathbf{R}$, where $R_{1}, \ldots, R_{k}$ are finite index sets. Moreover, extending the matrix case in the obvious way to $k$ dimensions, we say that a tensor $\mathbf{C}: R_{1} \times R_{2} \times \cdots \times R_{k} \rightarrow \mathbf{R}$ is a cut tensor if there exist sets $S_{i} \subseteq R_{i}$ for $i=1,2, \ldots, k$ and a real number $d$ such that

$$
\mathbf{C}\left(i_{1}, i_{2}, \ldots, i_{k}\right)= \begin{cases}d & \text { if }\left(i_{1}, i_{2}, \ldots, i_{k}\right) \in S_{1} \times S_{2} \times \cdots \times S_{k} \\ 0 & \text { otherwise }\end{cases}
$$

In this case we write $\mathbf{C}=\operatorname{CUT}\left(d, S_{1}, \ldots, S_{k}\right)$. Further, we define the cut norm of a tensor as

$$
\|\mathbf{M}\|_{\square}=\max _{S_{i} \subseteq R_{i}}\left|\mathbf{M}\left(S_{1}, S_{2}, \ldots, S_{k}\right)\right|, \quad \text { where } \quad \mathbf{M}\left(S_{1}, \ldots, S_{k}\right)=\sum_{\left(s_{1}, \ldots, s_{k}\right) \in S_{1} \times \cdots \times S_{k}} \mathbf{M}\left(s_{1}, \ldots, s_{k}\right) .
$$

Let $\mathbf{A}: R_{1} \times R_{2} \times \cdots \times R_{k} \rightarrow\{0,1\}$ be a $0 / 1$ tensor. Set $k_{1}=\lfloor k / 2\rfloor$. Then letting $\mathcal{R}=R_{1} \times R_{2} \times \cdots \times R_{k_{1}}$ and $\mathcal{C}=R_{k_{1}+1} \times R_{k_{1}+2} \times \cdots \times R_{k}$, we define a (2-dimensional) matrix $\mathbf{B}=\mathbf{B}(\mathbf{A}): \mathcal{R} \times \mathcal{C} \rightarrow\{0,1\}$ by

$$
\begin{equation*}
\mathbf{B}\left(\left(i_{1}, i_{2}, \ldots, i_{k_{1}}\right),\left(i_{k_{1}+1}, i_{k_{1}+2}, \ldots, i_{k}\right)\right)=\mathbf{A}\left(i_{1}, i_{2}, \ldots, i_{k}\right) \tag{3}
\end{equation*}
$$

We say that $\mathbf{A}$ is $(C, \gamma)$-bounded if $\mathbf{B}(\mathbf{A})$ has this property.

Theorem 2 There are an algorithm ApxTensor, a polynomial $\Pi$ and a constant $\Gamma>0$ such that the following is true. Suppose that $C>1$ and $0<\varepsilon<\frac{1}{2}$. Let

$$
\gamma=\exp \left(-\Gamma(C / \varepsilon)^{2}\right)
$$

If $\mathbf{A}: R_{1} \times R_{2} \times \cdots \times R_{k} \rightarrow\{0,1\}$ is a $(C, \gamma)$-bounded $0 / 1$ tensor, then ApxTensor $(\mathbf{A}, C, \varepsilon)$ outputs cut tensors

$$
\mathbf{D}_{i}=C U T\left(d_{i}, S_{i}^{1}, \ldots, S_{i}^{k}\right) \quad\left(S_{i}^{1} \subset R_{1}, \ldots, S_{i}^{k} \subset R_{k}\right)
$$

for $i=1, \ldots, s$ with $s \leq(\Gamma C / \epsilon)^{2(k-1)}$ such that

$$
\left\|\mathbf{A}-\left(\mathbf{D}_{1}+\cdots+\mathbf{D}_{s}\right)\right\|_{\square} \leq \epsilon\|\mathbf{A}\|_{\square}
$$

Moreover, $\sum_{i=1}^{s} d_{i}^{2} \leq(C p)^{2} \Gamma^{2 k}$. The running time is $\left(2^{(C / \varepsilon)^{2}}+(C / \varepsilon)^{3 k}\right) \cdot \Pi\left(\left|R_{1} \times \cdots \times R_{k}\right|\right)$.
If $\mathcal{R}_{1}, \ldots, \mathcal{R}_{k}$ are partitions of $R_{1}, \ldots, R_{k}$, then we define a tensor $\mathbf{A}_{\mathcal{R}_{1} \times \cdots \times \mathcal{R}_{k}}: R_{1} \times \cdots \times R_{k} \rightarrow[0,1]$ as follows: if $t_{i} \in \rho_{i} \in \mathcal{R}_{i}$ for $i=1, \ldots, k$, then we set

$$
\mathbf{A}_{\mathcal{R}_{1} \times \cdots \times \mathcal{R}_{k}}\left(t_{1}, \ldots, t_{k}\right)=\frac{\mathbf{A}\left(\rho_{1}, \ldots, \rho_{k}\right)}{\prod_{i=1}^{k}\left|\rho_{i}\right|}=\frac{\sum_{\left(v_{1}, \ldots, v_{k}\right) \in \rho_{1} \times \cdots \times \rho_{k}} \mathbf{A}\left(v_{1}, \ldots, v_{k}\right)}{\prod_{i=1}^{k}\left|\rho_{i}\right|}
$$

In words, on every rectangle $\rho_{1} \times \cdots \times \rho_{k}$ made up of partition classes $\rho_{i} \in \mathcal{R}_{i}$ the entry of $\mathbf{A}_{\mathcal{R}_{1} \times \cdots \times \mathcal{R}_{k}}$ is the average of A over that rectangle.

Corollary 3 There are an algorithm Part Tensor, a polynomial $\Pi$ and a constant $\Gamma>0$ such that the following is true. Suppose that $C>0$ and $0<\varepsilon<\frac{1}{2}$. Let $\gamma=\exp \left(-\Gamma(C / \varepsilon)^{2}\right)$. If $\mathbf{A}: R_{1} \times R_{2} \times \cdots \times R_{k} \rightarrow\{0,1\}$ is a $(C, \gamma)$-bounded $0 / 1$ tensor, then Part $\operatorname{Tensor}(\mathbf{A}, C, \varepsilon)$ computes partitions $\mathcal{R}_{1}, \ldots, \mathcal{R}_{k}$ of $R_{1}, \ldots, R_{k}$ such that

$$
\left\|\mathbf{A}-\mathbf{A}_{\mathcal{R}_{1} \times \cdots \times \mathcal{R}_{k}}\right\|_{\square}<\varepsilon\|\mathbf{A}\|_{\square} .
$$

Each of the partitions $\mathcal{R}_{i}$ consists of at most $\exp \left((\Gamma C / \varepsilon)^{2(k-1)}\right)$ classes. The running time is bounded by

$$
\left[\exp \left((\Gamma C / \varepsilon)^{2(k-1)}\right)+(C / \varepsilon)^{3 k}\right] \Pi\left(n^{k}\right)
$$

### 1.4 More General Coefficient Values

It is not absolutely necessary to assume that our matrices and tensors are $0 / 1$ valued. For convenience, we will only describe the case where $\mathbf{A}$ is an $m \times n$ matrix (i.e., 2 -dimensional). It will be apparant how to generalise the results to higher dimensions.

We can in fact assume that our coefficient entries are in $\{0,1,2, \ldots, d\}$ for some positive integer $d=O(1)$. We write $\mathbf{A}=\mathbf{A}^{(1)}+\cdots+\mathbf{A}^{(d)}$ where $\mathbf{A}^{(r)}(i, j)=\overrightarrow{1}_{\mathbf{A}_{i, j} \geq r / d}$ for $r=1, \ldots, d$. Let $p$ be the density of $\mathbf{A}^{(1)}$ i.e. the proportion of non-zero values. We assume that $\mathbf{A}^{(1)}$ is $(C, \gamma)$-bounded, with $\gamma=\gamma(\epsilon / d, C)$ as in Theorem 1. Then $\mathbf{A}^{(2)}, \ldots, \mathbf{A}^{(d)}$ have the following property: if $S \subset[m]$ and $T \subset[n]$ are sets such that $|S| \geq \gamma m$ and $|T| \geq \gamma n$, then $A^{(j)}(S, T) \leq C|S \times T| p(j=2, \ldots, d)$. Therefore, we can apply the algorithm ApxMatrix from Theorem 1 to each of these matrices and replace each $\mathbf{A}^{(i)}$ by a sum of cut-matrices that is within cut-norm $\epsilon m n p / d$ of its cut-norm. The sum of these cut-matrix approximations is then within $\epsilon m n p$ cut-norm of $\mathbf{A}$.

### 1.5 An approximation algorithm for bounded MAX CSPs

Let $V=\left\{x_{1}, \ldots, x_{n}\right\}$ be a set of $n$ Boolean variables. A (binary) $k$-constraint over $V$ is a map $\phi:\{0,1\}^{V_{\phi}} \rightarrow\{0,1\}$ that is not identically zero, where $V_{\phi} \subset V$ is a set of size $k$. For an assignment $\sigma \in\{0,1\}^{V}$ we let $\phi(\sigma)=$ $\phi(\sigma(x))_{x \in V_{\phi}}$. Further, a $k$-CSP instance over $V$ is a set $\mathcal{F}$ of $k$-constraints over $V$, and we define

$$
\operatorname{OPT}(\mathcal{F})=\max _{\sigma \in\{0,1\}^{V}} \sum_{\phi \in \mathcal{F}} \phi(\sigma)
$$

We let $\Psi=\Psi_{k}$ be the set of all $2^{2^{k}}-1$ non-zero maps $\{0,1\}^{k} \rightarrow\{0,1\}$. Let $\psi \in \Psi$ and let $\phi:\{0,1\}^{V_{\phi}} \rightarrow\{0,1\}$ be a $k$-constraint, where $V_{\phi}=\left\{x_{i_{1}}, \ldots, x_{i_{k}}\right\}$ with $1 \leq i_{1}<\cdots<i_{k} \leq n$. Then we say that $\phi$ is of type $\psi$ if for any $\sigma: V_{\phi} \rightarrow\{0,1\}$ we have $\psi\left(\sigma\left(x_{i_{1}}\right), \ldots, \sigma\left(x_{i_{k}}\right)\right)=\phi(\sigma)$. With this notion we can represent a $k$-CSP instance $\mathcal{F}$ by a family $\left(\mathbf{A}_{\mathcal{F}}^{\psi}\right)_{\psi \in \Psi}$ of $2^{2^{k}}-1 k$-tensors as follows. We let

$$
\mathbf{A}_{\mathcal{F}}^{\psi}\left(i_{1}, \ldots, i_{k}\right)= \begin{cases}1 & \text { if there is a } \phi \in \mathcal{F} \text { of type } \psi \text { with } V_{\phi}=\left\{x_{i_{1}}, \ldots, x_{i_{k}}\right\} \\ 0 & \text { otherwise }\end{cases}
$$

Further, we say that $\mathcal{F}$ is $(C, \gamma)$-bounded if all tensors $\mathbf{A}_{\mathcal{F}}^{\psi}$ are $(C, \gamma)$-bounded $(\psi \in \Psi)$.
Theorem 3 There are an algorithm ApxCSP, a constant $\Gamma>0$, and a polynomial $\Pi$ such that for any $k, C>1$, $0<\varepsilon<\frac{1}{2}$ there is a number $n_{0}=n_{0}(C, \varepsilon, k)$ such that the following is true. Let

$$
\gamma=\exp \left(-\Gamma 2^{-2^{k}-2 k-2}(C / \varepsilon)^{2}\right)
$$

If $\mathcal{F}$ is $(C, \gamma)$-bounded $k$-CSP instance over $V=\left\{x_{1}, \ldots, x_{n}\right\}$ for some $n \geq n_{0}$, then $\operatorname{ApxCSP}(\mathcal{F}, C, \varepsilon)$ outputs an assignment $\sigma: V \rightarrow\{0,1\}$ such that

$$
\sum_{\phi \in \mathcal{F}} \phi(\sigma) \geq(1-\varepsilon) O P T(\mathcal{F})
$$

The running time is at most $\Pi\left[\exp \left(k 2^{k} 2^{2^{k}}(C / \varepsilon)^{2 k} \ln (C / \varepsilon)\right) n^{k}\right]$.

### 1.6 Related work

### 1.6.1 Approximating dense matrices and tensors.

As mentioned earlier, Frieze and Kannan [16] dealt with dense matrices and tensors. More precisely, they showed that for a tensor $\mathbf{A}: R_{1} \times \cdots \times R_{k} \rightarrow[0,1]$ and an $\varepsilon>0$ one can compute cut tensors $\mathbf{D}_{1}, \ldots, \mathbf{D}_{s}$ such that $\left\|\mathbf{A}-\sum_{i=1}^{s} \mathbf{D}_{i}\right\|_{\square}<\varepsilon\left|R_{1} \times \cdots \times R_{k}\right|$ in time $O\left(\varepsilon^{2(1-k)} \operatorname{poly} \log (1 / \varepsilon)\right)$ with $s \leq O(\varepsilon)^{2(1-k)}$ as $\varepsilon \rightarrow 0$. Let us point out two things.

1. The running time of their algorithm depends only on $\varepsilon$, and not on the size of $\mathbf{A}$. This is achieved by randomization. Basically the algorithm just works with a bounded (by a function of $\varepsilon$ only) size sample of the input data, and produces an implicit representation of the desired decomposition. (Further results of this type can be found in Arora, Karger and Karpinski [5], Fernandez de la Vega [12], Goldwasser, Goldreich and Ron [18], Alon, de la Vega, Kannan and Karpinski [2] and de la Vega, Kannan, Karpinski and Vempala [13]).) Of course, if $\mathbf{A}: R_{1} \times \cdots \times R_{k} \rightarrow\{0,1\}$ is a sparse $0 / 1$ tensor with density $p=\|\mathbf{A}\|_{\square} /\left|R_{1} \times \cdots \times R_{k}\right|=o(1)$ as the problem size $N=\left|R_{1} \times \cdots \times R_{k}\right| \rightarrow \infty$, then this sampling approach cannot yield an approximation within $\varepsilon N p$. For any constant sized sample of $\mathbf{A}$ is likely to be just identically 0 . Therefore, in the present work we do not aim for a sublinear running time.
2. The error term $\varepsilon\left|R_{1} \times \cdots \times R_{k}\right|$ does not account for the density of $\mathbf{A}$. For example, suppose that $\mathbf{A}$ is the adjacency matrix of a graph $G=(V, E)$ on $n$ vertices with density $p=2 n^{-2}|E|$. Then the algorithm from [16] can be used to compute a cut norm approximation of $\mathbf{A}$ to within $\varepsilon n^{2}$ for any $\varepsilon>0$. Hence, we can use this approximation to solve graph partitioning problems such as MAX CUT within an additive error of $\varepsilon n^{2}$ (edges). This is why this approach is limited to dense problem instances: if the total number of edges is of lower order than $n^{2}$, then an approximation within an additive $\varepsilon n^{2}$ for a fixed $\varepsilon>0$ is of little value. For similar reasons the techniques of [16] only apply to dense problem instances of $k$-ary MAX CSP problems, i.e., instances with at least $\Omega\left(n^{k}\right)$ constraints, where $n$ is the number of variables.

In spite of these differences, some of the algorithms that we present are very similar to those from [16]. Thus, our main contribution is to analyze these algorithms on sparse matrices/graphs/tensors. For instance, the matrix approximation algorithm for Theorem 1 is almost identical to the procedure described in [16, Section 4.1]. The only difference
is that [16] employs as a subroutine a combinatorial procedure for approximating the cut norm of a given $m \times n$ matrix within an additive error of $\varepsilon m n$, whereas here we need to approximate the cut norm within a constant multiplicative factor. To this end, we rely on an algorithm of Alon and Naor [4] (which is based on semidefinite programming). Nonetheless, as we shall see in Section 3 new ideas are necessary to analyze, e.g., the number of cut matrices that are necessary to approximate the input matrix $\mathbf{A}$ within the desired $\varepsilon\|A\|_{\square}$ in the cut norm (rather than within $\varepsilon m n$ ).

### 1.6.2 Szemerédi's regularity lemma.

Theorem 2 and the concept of weak regular partitions is related to Szemerédi's well-known regularity lemma [24]. While [24] only deals with "dense" graphs, Kohayakawa [22] and Rödl [23] independently extended the regularity lemma to the sparse case; for a comprehensive survey on the subject see Gerke and Steger [17]. They showed that for any $\varepsilon>0$ and any $C>0$ there is a number $\gamma$ such any $(C, \gamma)$-bounded graph has a regular partition $\left(V_{1}, \ldots, V_{s}\right)$ in the following sense.

- We have $\left|V_{i}-n / s\right| \leq 1$ for all $i$.
- All but $\varepsilon s^{2}$ pairs $\left(V_{i}, V_{j}\right)$ satisfy the following. For any two sets $S \subset V_{i}, T \subset V_{j}$ of size $|S| \geq \varepsilon\left|V_{i}\right|,|T| \geq \varepsilon\left|V_{j}\right|$ we have

$$
\begin{equation*}
\left|e_{G}(S, T)-\frac{|S \times T|}{\left|V_{i} \times V_{j}\right|} \cdot e_{G}\left(V_{i}, V_{j}\right)\right| \leq \varepsilon e_{G}\left(V_{i}, V_{j}\right) \tag{4}
\end{equation*}
$$

The number $s$ of classes is bounded by a function $\mathcal{T}(C / \varepsilon)$, i.e., it is independent of $n$. This is the key fact that makes Szemerédi's lemma so useful in extremal combinatorics. However, from an algorithmic perspective the bound $\mathcal{T}(C / \varepsilon)$ is somewhat disappointing, because it is a tower function of height $(C / \varepsilon)^{3}$ :

$$
\left.\begin{array}{llll} 
& & & \\
& 2 & &
\end{array}\right\}(C / \varepsilon)^{3}
$$

In fact, there is an infinite family of graphs for which the number of classes in the smallest $\varepsilon$-regular Szemerédi partition is a tower of height $C / \varepsilon$ [19]. Moreover, the number $\gamma$ required in the boundedness condition is as tiny as $\mathcal{T}\left((C / \varepsilon)^{3}\right)^{-1}$.

While [22, 23, 24] focus on proving that a regular partition exists, [1,3] deal with algorithmic versions of the regularity lemma. In the dense case (i.e., $|E|=\Omega\left(n^{2}\right)$ ) there is a purely combinatorial algorithm [3] with running time $\mathcal{T}\left(\varepsilon^{-3}\right) \cdot \operatorname{poly}(n)$. Moreover, an algorithm for the sparse case was presented in [1]; the running time is $\mathcal{T}\left((C / \varepsilon)^{-3}\right) \cdot \operatorname{poly}(n)$ for $(C, \gamma)$-bounded graphs, and the algorithm is based on the semidefinite programming algorithm for approximating the cut norm from [4]. For instance, this yields an algorithm for approximating the MAX CUT on $(C, \gamma)$-bounded graphs within $1-\varepsilon$ in time $\mathcal{T}\left((C / \varepsilon)^{3}\right) \cdot \operatorname{poly}(n)$.

Corollary 2 relates to [1] as follows. While the "strong" regularity condition (4) takes into account the "microscopic" edge distribution within (almost) each pair ( $V_{i}, V_{j}$ ), the "weak" regularity concept from Corollary 2 just provides a "macroscopic" approximation w.r.t. the cut norm. This approximation is sufficiently strong for algorithmic applications such as MAX CUT (but it would not suffice for applications in extremal combinatorics that rely on the "counting lemma"). In effect, the algorithm is more efficient. Indeed, instead of scaling as a tower function $\mathcal{T}\left((C / \varepsilon)^{5}\right)$, the running time of the algorithm WeakPartition from Corollary 2 is bounded by $\exp \left(O(C / \varepsilon)^{2}\right)$ in terms of $C, \varepsilon$. Although this may still seem impractical, this is just a worst-case upper bound, and it is quite conceivable that it is practically much easier to find a good approximation in the cut norm than a good regular partition. Besides, as Theorem 1 shows, one can approximate a $(C, \gamma)$-bounded adjacency matrix by a sum of $O(C / \varepsilon)^{2}$ cut matrices (if the actual partition of the vertex set is not needed), thus avoiding the exponential dependence on $C / \varepsilon$. Similarly, the parameter $\gamma$ required in the boundedness condition is just $\gamma=\exp \left(-O(C / \varepsilon)^{2}\right)$, rather than $\gamma=1 / \mathcal{T}\left((C / \varepsilon)^{3}\right)$ as in [1]. Consequently, Corollary 2 applies to a larger class of graphs.

A further novel aspect here is that we extend our results to $k$-dimensional tensors (or $k$-uniform hypergraphs). This point is not addressed in [1].

## 2 Preliminaries

An important ingredient to the algorithm ApxMatrix for Theorem 1 is the the following algorithmic version of Grothendieck's inequality from Alon and Naor [4].

Theorem 4 There is a polynomial time algorithm and a number $\alpha_{0}>0$ that has the following property. Given a $m \times n$ matrix $\mathbf{M}$, the algorithm outputs sets $S \subset[m]$ and $T \subset[n]$ such that $|\mathbf{M}(S, T)| \geq \alpha_{0}\|\mathbf{M}\|_{\square}$.

Alon and Naor present a randomized algorithm with $\alpha_{0}>0.56$, and a deterministic one with $\alpha_{0} \geq 0.03$.
The algorithm ApxTensor for Theorem 2 employs an algorithm FKTensor from [16] as a subroutine.
Theorem 5 There are a polynomial $\Pi_{F K}$, an algorithm $F K T e n s o r$ and a number $\Gamma_{F K}>0$ such that the following is true. Suppose that $\mathbf{M}: R_{1} \times \cdots \times R_{k} \rightarrow[0,1]$ is a tensor and let $0<\delta<1$. Then FKTensor $(\mathbf{M}, \delta)$ outputs cut tensors $\mathbf{D}_{1}, \ldots, \mathbf{D}_{s}$ such that $\left\|\mathbf{M}-\mathbf{D}_{1}-\cdots-\mathbf{D}_{s}\right\|_{\square} \leq \delta \prod_{i=1}^{k}\left|R_{i}\right|$ and $s \leq\left(\Gamma_{F K} / \delta\right)^{2(k-1)}$. Moreover, $\sum_{i=1}^{s}\left\|\mathbf{D}_{i}\right\|_{\infty}^{2} \leq \Gamma_{F K}^{k}$, and the running time is at most $\delta^{-3 k} \Pi_{F K}\left(\prod_{i=1}^{k}\left|R_{i}\right|\right)$.

Actually Frieze and Kannan have a slightly stronger statement [16, Section 6] (better running time), but the above is sufficient for our purposes and easier to state.

If $\mathbf{M}$ is a real $m \times n$ matrix, then we let

$$
\|\mathbf{M}\|_{F}=\sqrt{\sum_{i=1}^{m} \sum_{j=1}^{n} \mathbf{M}_{i j}^{2}}
$$

signify the Frobenius norm of M. Moreover, if $G$ is a graph, then we denote the vertex set of $G$ by $V(G)$ and the edge set by $E(G)$. For sets $S, T \subset V(G)$ we let $e_{G}(S, T)$ signify the number of $S$ - $T$-edges of $G$, and $e_{G}(S)$ signifies the number of edges spanned by $S$.

Suppose that $X$ is a set and that $\mathcal{P}_{1}, \mathcal{P}_{2}$ are partitions of $X$. We say that $\mathcal{P}_{1}$ is coarser than $\mathcal{P}_{2}$ if each class of $\mathcal{P}_{2}$ is contained in a class of $\mathcal{P}_{1}$. If $\mathcal{S}$ is an arbitrary set of subsets of $X$, then there is a unique partition $\mathcal{P}$ of $X$ such that

1. each set in $X$ is a union of classes of $\mathcal{P}$,
2. $\mathcal{P}$ is coarser than any other partition that satisfies 1 .

This partition $\mathcal{P}$ has at most $2^{|\mathcal{S}|}$ classes.

## 3 Approximating and partitioning $0 / 1$ matrices and graphs

### 3.1 Proof of Theorem 1

Let $C>1$ and $0<\varepsilon<\frac{1}{2}$. Moreover, let $\alpha_{0}$ be the constant from Theorem 4 and set

$$
\kappa=\frac{513 C^{2}}{\varepsilon^{2} \alpha_{0}}, \quad \gamma=\frac{\varepsilon \alpha_{0}}{2^{10 \kappa} C}, \quad \gamma^{\prime}=2^{\kappa} \gamma
$$

Throughout this section we assume that $\mathbf{A}$ is a $(C, \gamma)$ bounded $0 / 1$ matrix of size $m \times n$.

```
Algorithm 6 ApxMatrix \((\mathbf{A}, C, \varepsilon)\)
Input: A \(0 / 1\) matrix \(\mathbf{A}\) of size \(m \times n\), numbers \(C, \varepsilon>0\).
Output: A sequence of cut matrices.
```

1. $\quad$ Set $\mathbf{A}_{0}=\mathbf{A}$.
2. For $j=0,1,2, \ldots, \kappa$ do
3. $\quad$ Compute sets $S_{j+1}, T_{j+1}$ of sizes $\left|S_{j+1}\right| \geq m / 2,\left|T_{j+1}\right| \geq n / 2$ s.t. $\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right| \geq \alpha_{0}\left\|\mathbf{A}_{j}\right\|_{\square} / 4$.
4. If $\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right|<\alpha_{0} \varepsilon m n p / 4$ and $j \geq 1$, then output the cut matrices $\mathbf{D}_{1}, \ldots, \mathbf{D}_{j}$ and halt. Else,
5. Compute

$$
d_{j+1}=\frac{\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)}{\left|S_{j+1}\right|\left|T_{j+1}\right|},
$$

set $\mathbf{D}_{j+1}=\operatorname{CUT}\left(d_{j+1}, S_{j+1}, T_{j+1}\right)$, and let $\mathbf{A}_{j+1}=\mathbf{A}_{j}-\mathbf{D}_{j+1}$.
6. Output "failure".

In order to approximate the given $0 / 1$ matrix $\mathbf{A}$ by a sum $\mathbf{D}_{1}+\cdots+\mathbf{D}_{j}$ of cut matrices, ApxMatrix proceeds as follows. After $j$ iterations, $\mathbf{A}_{j}=\mathbf{A}-\sum_{i=1}^{j} \mathbf{D}_{j}$ is the "error term" that results from approximating $\mathbf{A}$ by $\sum_{i=1}^{j} \mathbf{D}_{j}$. Thus, the goal is to eventually achieve an error term $\mathbf{A}_{j}$ that has a small cut norm. Therefore, Step 3 computes sets $S_{j+1}, T_{j+1}$ of rows and columns such that $\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right|$ is a good approximation of the cut norm of $\mathbf{A}_{j}$. If the term $\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right|$ (and hence the cut norm of $\left.\mathbf{A}_{j}\right)$ is small, then Step 4 terminates and outputs the cut matrices $\mathbf{D}_{1}, \ldots, \mathbf{D}_{j}$. Otherwise, $S_{j+1}, T_{j+1}$ witness a set of rows/columns on which $\sum_{i=1}^{j} \mathbf{D}_{j}$ does not provide a good enough approximation. Therefore, Step 5 adds a further "patch" $\mathbf{D}_{j+1}$, which is a cut matrix whose value on $S_{j+1} \times T_{j+1}$ is just the average $d_{j+1}$ of $\mathbf{A}_{j}$ over that square (note that $d_{j+1}$ may be negative). This ensures that $\mathbf{A}_{j+1}\left(S_{j+1}, T_{j+1}\right)=0$, and thus takes care of the discrepancy witnessed by $S_{j+1}, T_{j+1}$.

If the algorithm outputs cut matrices $\mathbf{D}_{1}, \ldots, \mathbf{D}_{j}$, then clearly

$$
\left\|\mathbf{A}-\left(\mathbf{D}_{1}+\cdots+\mathbf{D}_{j}\right)\right\|_{\square}=\left\|\mathbf{A}_{j}\right\|_{\square} \leq \varepsilon m n p=\varepsilon\|\mathbf{A}\|_{\square},
$$

because of the halting condition in Step 4. Hence, in order to establish Theorem 1, we need to prove that
(a) Step 3 of ApxMatrix can be implemented by a polynomial time algorithm,
(b) the halting condition in Step 4 is satisfied for some $1 \leq j \leq \kappa$.

Proposition 1 In Step 3 the sets $S_{j+1}, T_{j+1}$ can be computed in time poly $(m n)$.
Proof To obtain $S_{j+1}, T_{j+1}$, we use the polynomial time algorithm from Theorem 4, which yields sets $S_{j+1}^{\prime}, T_{j+1}^{\prime}$ such that $\left|\mathbf{A}_{j}\left(S_{j+1}^{\prime}, T_{j+1}\right)^{\prime}\right| \geq \alpha_{0}\left\|\mathbf{A}_{j}\right\|_{\square}$. If $\left|S_{j+1}^{\prime}\right| \geq n / 2$ then we take $S_{j+1}=S_{j+1}^{\prime}$. If $\left|S_{j+1}^{\prime}\right|<n / 2$ then since

$$
\mathbf{A}\left(R, T_{j+1}^{\prime}\right)=\mathbf{A}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)+\mathbf{A}\left(R \backslash S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)
$$

we get $\max \left\{\left|\mathbf{A}\left(R, T_{j+1}^{\prime}\right)\right|, \mathbf{A}\left(R \backslash S_{j+1}^{\prime}, T_{j+1}^{\prime}\right) \mid\right\} \geq \alpha_{0}\left\|\mathbf{A}_{0}\right\|_{\square} / 2$. We can therefore take either $R$ or $R \backslash S_{j+1}^{\prime}$ as our set $S_{j+1}$ and note it is at least $n / 2$ in size. We perform the same operation to get $T_{j+1}$, losing (at most) another factor 2 in the approximation.

With respect to (b), we will study the Frobenius norm of $\mathbf{A}_{j}$. Namely, it is not difficult to show that $\left\|\mathbf{A}_{j}\right\|_{F}^{2} \leq$ $\|\mathbf{A}\|_{F}^{2}\left(1-j \cdot \alpha_{0}^{2} \varepsilon^{2} p / 4\right)$. Since trivially $\left\|\mathbf{A}_{j}\right\|_{F} \geq 0$, this implies that the total number of iterations is at most $4 /\left(\alpha_{0}^{2} \varepsilon^{2} p\right)$. Hence, if $p$ is bounded from below by a constant, then this argument shows that the total number of iterations is bounded by a number that does not depend on $n, m$. In fact, this is the basic argument used to establish the matrix decomposition theorem in [16, Section 4.1].

But in the present work we do not assume that $p$ remains bounded away from 0 by a number independent of $n, m$. In effect, the aforementioned argument does not apply. As it turns out, the problem is that the above argument just uses the trivial lower bound $\left\|\mathbf{A}_{j}\right\|_{F}^{2} \geq 0$. By contrast, the basic idea here is to use the boundedness condition to establish $\left\|\mathbf{A}_{F}\right\|^{2}\left(1-C^{2} p\right)$ as a lower bound. Indeed, if we could show that $\left\|\mathbf{A}_{j}\right\|_{F}^{2} \geq\left\|\mathbf{A}_{F}\right\|^{2}\left(1-C^{2} p\right)$ for all $j$, then the bound $\left\|\mathbf{A}_{j}\right\|_{F}^{2} \leq\|\mathbf{A}\|_{F}^{2}\left(1-j \cdot \alpha_{0}^{2} \varepsilon^{2} p / 4\right)$ would imply that the number of iterations is at most $4 C^{2} /\left(\alpha_{0}^{2} \varepsilon^{2}\right)$, and thus independent of $m, n, p$.

However, we can't quite use the boundedness condition to prove that $\left\|\mathbf{A}_{j}\right\|_{F}^{2} \geq\left\|\mathbf{A}_{F}\right\|^{2}\left(1-C^{2} p\right)$. The reason is that the boundedness condition only applies to "sufficiently large" sets, i.e., sets of size at least $\gamma n$. Therefore, to show that ApxMatrix stops after at most $\kappa$ iterations, we will consider slightly different sequences of matrices $\mathbf{D}_{j}^{\prime}, \mathbf{A}_{j}^{\prime}$, to
which the boundedness condition applies. The matrices $\mathbf{D}_{j}^{\prime}, \mathbf{A}_{j}^{\prime}$ will be "close" to $\mathbf{D}_{j}, \mathbf{A}_{j}$ in cut norm, and to bound the number of iterations we are going to investigate the Frobenius norm of $\mathbf{A}_{j}^{\prime}$.

We construct the matrices $\mathbf{D}_{j}^{\prime}, \mathbf{A}_{j}^{\prime}$ as follows. Let us assume (for contradiction) that ApxMatrix outputs "failure", i.e., the number of iterations performed by Steps $2-5$ is $\kappa$. Then during these $\kappa$ iterations the algorithm constructed sets $S_{1}, \ldots, S_{\kappa}$ of rows and $T_{1}, \ldots, T_{\kappa}$ of columns. Let $\mathcal{S}$ be the coarsest partition of the set $[\mathrm{m}]$ of row indices such that each $S_{i}$ is a union of classes of $\mathcal{S}$ (thus, $\mathcal{S}$ consists of the classes of the Venn diagram of the sets $S_{1}, \ldots, S_{\kappa}$ ). Similarly, let $\mathcal{T}$ be the coarsest partition of the columns set $[n]$ such that every $T_{i}$ is a union of classes of $\mathcal{T}$. Clearly, both $\mathcal{S}$ and $\mathcal{T}$ have at most $2^{\kappa}$ classes. The reason why the boundedness condition does not imply directly that $\left\|\mathbf{A}_{j}\right\|_{F}^{2} \geq\left\|\mathbf{A}_{F}\right\|^{2}\left(1-C^{2} p\right)$ is that some classes of $\mathcal{S}$ and $\mathcal{T}$ may have size less than $\gamma m$ or $\gamma n$. Therefore, we let

$$
R_{0}=\bigcup_{S \in \mathcal{S}:|S|<\gamma m} S, \quad C_{0}=\bigcup_{T \in \mathcal{T}:|T|<\gamma n} T
$$

comprise the "small" classes of the partitions $\mathcal{S}, \mathcal{T}$. Setting $\gamma^{\prime}=2^{\kappa} \gamma$, we have

$$
\begin{equation*}
\left|R_{0}\right| \leq \gamma^{\prime} m, \quad\left|C_{0}\right| \leq \gamma^{\prime} n \tag{5}
\end{equation*}
$$

Further, let $\mathbf{A}_{0}^{\prime}=\mathbf{A}^{\prime}$ be the matrix obtained from $\mathbf{A}$ by replacing all rows in $R_{0}$ and all columns in $C_{0}$ by 0 . In addition, define inductively sets $S_{j}^{\prime}=S_{j} \backslash R_{0}$ and $T_{j}^{\prime}=T_{j} \backslash T_{0}$ and

$$
d_{j+1}^{\prime}=\frac{\mathbf{A}_{j}^{\prime}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)}{\left|S_{j+1}^{\prime}\right|\left|T_{j+1}^{\prime}\right|}, \quad \mathbf{D}_{j+1}^{\prime}=\operatorname{CUT}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}, d_{j+1}^{\prime}\right), \quad \mathbf{A}_{j+1}^{\prime}=\mathbf{A}_{j}^{\prime}-\mathbf{D}_{j+1}^{\prime}
$$

Let $\mathcal{S}^{\prime}$ be the coarsest partition of $[m] \backslash R_{0}$ such that each $S_{j}^{\prime}$ is a union of classes of $\mathcal{S}^{\prime}$, and define a partition $\mathcal{T}^{\prime}$ of $[n] \backslash C_{0}$ analogously w.r.t. the sets $T_{j}^{\prime}$. Then the construction of the sets $S_{j}^{\prime}, T_{j}^{\prime}$ readily implies:

Fact 7 All classes of $\mathcal{S}^{\prime}$ (resp. $\mathcal{T}^{\prime}$ ) have size at least $\gamma m$ (resp. $\gamma n$ ).
Consequently, we can use the boundedness condition to infer the following.
Lemma 1 For all $1 \leq j \leq \kappa$ we have $\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2} \geq\left\|\mathbf{A}^{\prime}\right\|_{F}^{2}\left(1-2 C^{2} p\right)$.
Proof Let $\mathbf{M}=\sum_{i=1}^{j} \mathbf{D}_{j}^{\prime}$. Then for any two sets $S \in \mathcal{S}^{\prime}, T \in \mathcal{T}^{\prime}$ the matrix $\mathbf{M}$ is constant on the square $S \times T$, because every $\mathbf{D}_{j}^{\prime}$ is a cut matrix on the square $S_{j}^{\prime} \times T_{j}^{\prime}$, and $S_{j}^{\prime}, T_{j}^{\prime}$ are unions of classes of $\mathcal{S}^{\prime}, \mathcal{T}^{\prime}$. Thus, letting $m_{S \times T}$ signify the value that $\mathbf{M}$ takes on $S \times T$, we obtain

$$
\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2}=\left\|\mathbf{A}^{\prime}-\mathbf{M}\right\|_{F}^{2}=\sum_{S \in \mathcal{S}^{\prime}, T \in \mathcal{T}^{\prime}} \sum_{(v, w) \in S \times T}\left(\mathbf{A}^{\prime}(v, w)-m_{S \times T}\right)^{2} .
$$

For any $S \in \mathcal{S}^{\prime}, T \in \mathcal{T}^{\prime}$ the $\operatorname{sum} \sum_{(v, w) \in S \times T}\left(\mathbf{A}^{\prime}(v, w)-m_{S \times T}\right)^{2}$ is minimized iff

$$
m_{S \times T}=m_{S \times T}^{*}=\mathbf{A}^{\prime}(S, T) /(|S| \cdot|T|)
$$

Therefore,

$$
\begin{aligned}
\sum_{(v, w) \in S \times T}\left(\mathbf{A}^{\prime}(v, w)-m_{S \times T}\right)^{2} & \geq \sum_{(v, w) \in S \times T}\left(\mathbf{A}^{\prime}(v, w)-m_{S \times T}^{*}\right)^{2} \\
& =\sum_{(v, w) \in S \times T} \mathbf{A}^{\prime}(v, w)^{2}-2 \mathbf{A}^{\prime}(S, T) m_{S \times T}^{*}+m_{S \times T}^{* 2}|S| \cdot|T| \\
& =\sum_{(v, w) \in S \times T} \mathbf{A}^{\prime}(v, w)^{2}-m_{S \times T}^{* 2}|S| \cdot|T| .
\end{aligned}
$$

Since $\mathbf{A}^{\prime}$ is $(C, \gamma)$ bounded and because $|S| \geq \gamma m,|T| \geq \gamma n$ by Lemma 7, we get $m_{S \times T}^{*} \leq C p$. Hence,

$$
\begin{equation*}
\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2} \geq\left\|\mathbf{A}^{\prime}\right\|_{F}^{2}-(C p)^{2} m n . \tag{6}
\end{equation*}
$$

Finally, using the fact that $\mathbf{A}$ and $\mathbf{A}^{\prime}$ are 0,1 matrices, we have

$$
\begin{aligned}
\|\mathbf{A}\|_{F}^{2}-\left\|\mathbf{A}^{\prime}\right\|_{F}^{2} & =\sum_{i=1}^{m} \sum_{j=1}^{n} \mathbf{A}_{i j}^{2}-\mathbf{A}_{i j}^{\prime 2}=\mathbf{A}\left(R_{0},[n]\right)+\mathbf{A}\left([m], C_{0}\right)-\mathbf{A}\left(R_{0}, C_{0}\right) \\
& \leq \mathbf{A}\left(R_{0},[n]\right)+\mathbf{A}\left([m], C_{0}\right) \stackrel{(5)}{\leq} 2 C \gamma^{\prime} m n p<m n p / 2
\end{aligned}
$$

whence $\left\|\mathbf{A}^{\prime}\right\|_{F}^{2} \geq\|\mathbf{A}\|_{F}^{2} / 2=m n p / 2$. Thus, the assertion follows from (6).
To show that our assumption that ApxMatrix performs at least $\kappa$ iterations yields a contradiction, we shall derive the following upper bound on $\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2}$.

Lemma 2 For all $1 \leq j \leq \kappa$ we have $\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2} \leq\left\|\mathbf{A}^{\prime}\right\|_{F}^{2}\left(1-j \cdot \alpha_{0}^{2} \varepsilon^{2} p / 256\right)$.
Combining Lemmas 1 and 2 and setting $j=\kappa$, we conclude $2 C^{2} \geq \kappa \cdot \alpha_{0}^{2} \varepsilon^{2} / 256$, which contradicts our choice of $\kappa$ (cf. (2)). This completes the proof of Theorem 1.

In the rest of this section we prove Lemma 2. The following lemma shows that $\mathbf{A}^{\prime}$ is close to $\mathbf{A}$ in cut norm. Recall that $\gamma^{\prime}=2^{\kappa} \gamma$.

Lemma 3 We have $\left\|\mathbf{A}-\mathbf{A}^{\prime}\right\|_{\square} \leq 2 C \gamma^{\prime} m n p$.
Proof For any sets $S \subset[m], T \subset[n]$ we have

$$
\left|\mathbf{A}(S, T)-\mathbf{A}^{\prime}(S, T)\right| \leq \mathbf{A}\left(R_{0},[n]\right)+\mathbf{A}\left([m], C_{0}\right)
$$

If $\left|R_{0}\right| \geq \gamma m$, then we let $R_{0}^{\prime}=R_{0}$; otherwise, let $R_{0}^{\prime} \supset R_{0}$ be any superset of size $\gamma m$. Then the fact that $\mathbf{A}$ is $(C, \gamma)$-bounded implies that for any $T \subseteq[n],|T| \geq \gamma n$,

$$
\begin{equation*}
\mathbf{A}\left(R_{0}, T\right) \leq \mathbf{A}\left(R_{0}^{\prime}, T\right) \leq C\left|R_{0}^{\prime}\right| n p \leq C \gamma^{\prime} m n p, \quad \text { from (5) } \tag{7}
\end{equation*}
$$

Applying the same argument to $\mathbf{A}\left([m], C_{0}\right)$ gives the result.
In addition, we also need to show that the matrices $\mathbf{D}_{j}, \mathbf{D}_{j}^{\prime}$ are close in cut norm. To this end, we derive a bound on the coefficients $d_{j}$ from Step 5 of ApxMatrix.

Lemma $4\left|d_{j}\right| \leq 2^{j} C p$ for all $1 \leq j \leq \kappa$.
Proof The proof is by induction on $j$. For $j=1$ we have

$$
d_{1}=\frac{\mathbf{A}_{0}\left(S_{1}, T_{1}\right)}{\left|S_{1}\right|\left|T_{1}\right|} \leq C p
$$

because $\mathbf{A}_{0}\left(S_{1}, T_{1}\right) \leq C p\left|S_{1}\right|\left|T_{1}\right|$ by the boundedness condition. Furthermore, assuming that $\left|d_{i}\right| \leq 2^{i} C p$ for all $i \leq j$, we can bound $d_{j+1}$ as follows.

$$
\begin{aligned}
\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right| & =\left|\mathbf{A}_{0}\left(S_{j+1}, T_{j+1}\right)+\sum_{i=1}^{j} d_{i}\right| S_{j+1} \cap S_{i}| | T_{j+1} \cap T_{i}| | \\
& \leq\left|S _ { j + 1 } \left\|T _ { j + 1 } \left|\cdot\left[C p+\sum_{i=1}^{j}\left|d_{i}\right|\right] \leq\left|S_{j+1}\left\|T_{j+1}\left|C p \sum_{i=0}^{j} 2^{i} \leq 2^{j+1}\right| S_{j+1}\right\| T_{j+1}\right| C p\right.\right.\right.
\end{aligned}
$$

Thus, $\left|d_{j+1}\right|=\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right| /\left(\left|S_{j+1}\right|\left|T_{j+1}\right|\right) \leq 2^{j+1} C p$.
Lemma 4 implies the following bound on the cut norm of $\mathbf{D}_{j}-\mathbf{D}_{j}^{\prime}$.
Corollary 4 For all $1 \leq j \leq \kappa$ we have $\left\|\mathbf{D}_{j}-\mathbf{D}_{j}^{\prime}\right\|_{\square} \leq 2^{8 j} C \gamma^{\prime} m n p$.

Proof We proceed by induction. The definitions of $d_{j}$ and $d_{j}^{\prime}$ imply that

$$
\begin{align*}
&\left|d_{j}^{\prime}-d_{j}\right|=\left|\frac{\mathbf{A}_{j-1}^{\prime}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)}{\left|S_{j}^{\prime}\right|\left|T_{j}^{\prime}\right|}-\frac{\mathbf{A}_{j-1}\left(S_{j}, T_{j}\right)}{\left|S_{j}\right|\left|T_{j}\right|}\right|=\frac{\left|\left|S_{j}\right|\right| T_{j}\left|\mathbf{A}_{j-1}^{\prime}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)-\left|S_{j}^{\prime}\right|\right| T_{j}^{\prime}\left|\mathbf{A}_{j-1}\left(S_{j}, T_{j}\right)\right|}{\left|S_{j}\right|\left|S_{j}^{\prime}\right|\left|T_{j}\right|\left|T_{j}^{\prime}\right|} \\
& \leq \frac{\left|\mathbf{A}_{j-1}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)-\mathbf{A}_{j-1}\left(S_{j}, T_{j}\right)\right|}{\left|S_{j}^{\prime}\right|\left|T_{j}^{\prime}\right|}+\frac{\left|\left(\left|S_{j}\right|\left|T_{j}\right|-\left|S_{j}^{\prime}\right|\left|T_{j}^{\prime}\right|\right) \mathbf{A}_{j-1}\left(S_{j}, T_{j}\right)\right|}{\left|S_{j}\right|\left|S_{j}^{\prime}\right|\left|T_{j}\right|\left|T_{j}^{\prime}\right|} \\
& \quad+\frac{\left|\mathbf{A}_{j-1}^{\prime}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)-\mathbf{A}_{j-1}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)\right|}{\left|S_{j}^{\prime}\right|\left|T_{j}^{\prime}\right|} \\
& \leq 16(m n)^{-1}\left(\left|\mathbf{A}_{j-1}\left(R_{0}, T_{j}\right)\right|+\left|\mathbf{A}_{j-1}\left(S_{j}, C_{0}\right)\right|\right) \\
& \quad+64(m n)^{-2}\left|\mathbf{A}_{j-1}\left(S_{j}, T_{j}\right)\right|\left(\left|R_{0}\right|\left|C_{0}\right|+\left|R_{0}\right|\left|T_{j}\right|+\left|S_{j}\right|\left|C_{0}\right|\right) \\
&+16(m n)^{-1}\left|\mathbf{A}_{j-1}^{\prime}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)-\mathbf{A}_{j-1}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)\right| \tag{8}
\end{align*}
$$

because $\left|S_{j}\right| \geq m / 2,\left|T_{j}\right| \geq n / 2$ by construction and $\left|S_{j}^{\prime}\right| \geq\left|S_{j}\right|-\left|R_{0}\right| \geq m / 4,\left|T_{j}\right| \geq\left|T_{j}\right|-\left|C_{0}\right| \geq n / 4$ by (5).
Observe that for $j=1$ the third term in (8) equals 0 . Furthermore, the boundedness condition and (5) imply that $\mathbf{A}\left(R_{0},[n]\right)+\mathbf{A}\left([n], C_{0}\right) \leq 2 C \gamma^{\prime} m n p$, (see (7)) and so we conclude that $\left|d_{1}^{\prime}-d_{1}\right| \leq C \gamma^{\prime} p(32+192)=224 C \gamma^{\prime} p$, and thus

$$
\left\|\mathbf{D}_{1}-\mathbf{D}_{1}^{\prime}\right\|_{\square} \leq\left|d_{1}-d_{1}^{\prime}\right| m n+2 C \gamma^{\prime} m n p \leq 226 C \gamma^{\prime} m n p
$$

Now assume inductively that for $1 \leq i \leq j-1$ we have

$$
\begin{equation*}
\left\|\mathbf{D}_{i}-\mathbf{D}_{i}^{\prime}\right\|_{\square} \leq 2^{8 j} C \gamma^{\prime} m n p \tag{9}
\end{equation*}
$$

Then, for $j>1$ Lemma 4 implies that

$$
\left|\mathbf{A}_{j-1}\left(R_{0}, T_{j}\right)\right| \leq\left|\mathbf{A}\left(R_{0}, T_{j}\right)\right|+\sum_{i=1}^{j-1}\left|d_{i}\right|\left|R_{0}\right| n \leq 2^{j} C \gamma^{\prime} m n p
$$

Similarly, $\left|\mathbf{A}_{j-1}\left(S_{j}, C_{0}\right)\right| \leq 2^{j} C \gamma^{\prime} m n p$. Thus the first term in (8) is at most $2^{j+5} C \gamma^{\prime} p$.
Moreover, once more by Lemma 4 we have

$$
\left|\mathbf{A}_{j-1}\left(S_{j}, T_{j}\right)\right| \leq\left|\mathbf{A}\left(S_{j}, T_{j}\right)\right|+m n \sum_{i=1}^{j-1}\left|d_{i}\right| \leq 2^{j} C m n p
$$

Consequently, the second term in (8) is at most $3 \cdot 2^{j+6} C \gamma^{\prime} p$. Finally, by induction we obtain

$$
\begin{aligned}
\left|\mathbf{A}_{j-1}^{\prime}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)-\mathbf{A}_{j-1}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)\right| & \leq \sum_{i=1}^{j-1}\left|\mathbf{D}_{i}^{\prime}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)-\mathbf{D}_{i}\left(S_{j}^{\prime}, T_{j}^{\prime}\right)\right| \leq \sum_{i=1}^{j-1}\left\|\mathbf{D}_{i}^{\prime}-\mathbf{D}_{i}\right\|_{\square} \\
& \leq C \gamma^{\prime} m n p \sum_{i=1}^{j-1} 2^{8 i}<2^{8 j-7} C \gamma^{\prime} m n p
\end{aligned}
$$

Hence, the third term in (8) is at most $2^{8 j-3} C \gamma^{\prime} p$. Plugging these bounds into (8) we get $\left|d_{j}^{\prime}-d_{j}\right| \leq 2^{8 j-1} C \gamma^{\prime} p$. Finally, (5) and Lemma 4 yield

$$
\left\|\mathbf{D}_{j}-\mathbf{D}_{j}^{\prime}\right\|_{\square} \leq\left|d_{j}-d_{j}^{\prime}\right| m n+\left(\left|R_{0}\right| n+m\left|C_{0}\right|\right)\left|d_{j}\right| \leq 2^{8 j-1} C \gamma^{\prime} m n p+2^{j+1} C \gamma^{\prime} m n p \leq 2^{8 j} C \gamma^{\prime} m n p
$$

as desired.
Combining Lemma 3 with Corollary 4, we conclude that the two matrices $\mathbf{A}_{j}=\mathbf{A}-\sum_{i=1}^{j} \mathbf{D}_{j}$ and $\mathbf{A}_{j}^{\prime}=$ $\mathbf{A}^{\prime}-\sum_{i=1}^{j} \mathbf{D}_{j}^{\prime}$ are close in cut norm.
Corollary 5 For all $1 \leq j \leq \kappa$ we have $\left\|\mathbf{A}_{j}^{\prime}-\mathbf{A}_{j}\right\|_{\square} \leq 2^{8 j+1} C \gamma^{\prime} m n p$.

Proof of Lemma 2. We are going to show that

$$
\begin{equation*}
\left\|\mathbf{A}_{j+1}^{\prime}\right\|_{F}^{2} \leq\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2}-\alpha_{0}^{2} \varepsilon^{2} m n p^{2} / 256 \tag{10}
\end{equation*}
$$

for any $1 \leq j<\kappa$; this bound immediately implies the assertion. Remember that $\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right| \geq \alpha_{0}\left\|\mathbf{A}_{j}\right\|_{\square} / 4$ by the construction of $S_{j+1}, T_{j+1}$ in Step 3 of ApxMatrix. Therefore, combining Corollary 5 with Lemmas 3 and 4, we obtain

$$
\begin{align*}
\left|\mathbf{A}_{j}^{\prime}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)\right| \geq & \left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)\right|-\left|\mathbf{A}_{j}\left(S_{j+1}, T_{j+1}\right)-\mathbf{A}_{j}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)\right| \\
& \quad-\left|\mathbf{A}_{j}^{\prime}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)-\mathbf{A}_{j}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)\right| \\
\geq & \alpha_{0} \varepsilon m n p / 4-\left(2 C \gamma^{\prime} m n p+\sum_{i=1}^{j} 2^{i} C \gamma^{\prime} m n p\right)-2^{8 j+1} C \gamma^{\prime} m n p \\
\geq & \alpha_{0} \varepsilon m n p / 8 \tag{11}
\end{align*}
$$

where the last inequality follows from our choice of $\gamma$ and the fact that $\gamma^{\prime}=2^{\kappa} \gamma$. Further, as

$$
d_{j+1}^{\prime}=\frac{\mathbf{A}_{j}^{\prime}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)}{\left|S_{j+1}^{\prime}\right|\left|T_{j+1}^{\prime}\right|}
$$

by construction, (11) implies that

$$
\begin{aligned}
\left\|\mathbf{A}_{j}^{\prime}\right\|_{F}^{2}-\left\|\mathbf{A}_{j+1}^{\prime}\right\|_{F}^{2} & =\sum_{(s, t) \in S_{j+1}^{\prime} \times T_{j+1}^{\prime}} \mathbf{A}_{j}^{\prime}(s, t)^{2}-\left(\mathbf{A}_{j}^{\prime}(s, t)-d_{j+1}^{\prime}\right)^{2} \\
& ={d_{j+1}^{\prime}}^{2} \mathbf{A}^{\prime}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)=\frac{\mathbf{A}_{j}^{\prime}\left(S_{j+1}^{\prime}, T_{j+1}^{\prime}\right)^{2}}{\left|S_{j+1}^{\prime}\right|\left|T_{j+1}^{\prime}\right|} \geq \frac{\left(\alpha_{0} \varepsilon m n p\right)^{2}}{256 m n}
\end{aligned}
$$

whence (10) follows.

### 3.2 Proof of Corollary 1

Let $0<\varepsilon<\frac{1}{2}$ and $C>1$, and let $\kappa$, $\gamma$ be as in (2). Given a $(C, \gamma)$-bounded matrix $\mathbf{A}$ of size $m \times n$ and the numbers $C, \varepsilon$, PartMatrix calls ApxMatrix $(\mathbf{A}, C, \varepsilon)$ to obtain cut matrices

$$
\mathbf{D}_{i}=\operatorname{CUT}\left(d_{i}, S_{i}, T_{i}\right) \quad(i=1, \ldots, s)
$$

for some $1 \leq s \leq \kappa$. Then, it computes the coarsest partition $\mathcal{S}$ of $[m]$ such that each class $S_{i}$ is a union of classes of $\mathcal{S}$ $(1 \leq i \leq s)$. Similarly, $\mathcal{T}$ is the coarsest partition of $[n]$ such that each class $T_{i}$ is a union of classes of $\mathcal{T}(1 \leq i \leq s)$.

This construction ensures that $|\mathcal{S}|,|\mathcal{T}| \leq 2^{s} \leq 2^{\kappa}$. Hence, the running time of PartMatrix is $\left(\kappa+2^{\kappa}\right) \Pi(m n)$ for a fixed polynomial $\Pi$. Thus, to complete the proof of Corollary 1 we just need to show that $\left\|\mathbf{A}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \leq$ $2 \varepsilon\|\mathbf{A}\|_{\square}$. Since the matrix $\mathbf{D}=\sum_{i=1}^{s} \mathbf{D}_{i}$ satisfies $\|\mathbf{A}-\mathbf{D}\|_{\square} \leq \varepsilon\|\mathbf{A}\|_{\square}$ by Theorem 1, it suffices to prove that

$$
\begin{equation*}
\left\|\mathbf{D}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \leq\|\mathbf{D}-\mathbf{A}\|_{\square} \tag{12}
\end{equation*}
$$

To prove (12), we use the same argument as in [16, Section 5]. Let $X \subset[m]$ and $Y \subset[n]$ be sets such that $\left|\left(\mathbf{D}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right)(X, Y)\right|=\left\|\mathbf{D}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square}$. On each square $S \times T$ with $S \in \mathcal{S}$ and $T \in \mathcal{T}$ the matrix $\mathbf{D}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}$ is constant. We may therefore assume that $X$ is a union of classes of $\mathcal{S}$ and $Y$ is a union of classes of $\mathcal{T}$. Furthermore, as $A(S, T)=A_{\mathcal{S} \times \mathcal{T}}(S, T)$ for any $S \in \mathcal{S}$ and $T \in \mathcal{T}$ by the definition of $A_{\mathcal{S} \times \mathcal{T}}(S, T)$, we conclude that

$$
\begin{aligned}
\left\|\mathbf{D}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} & =\left|\mathbf{D}(X, Y)-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}(X, Y)\right|=\left|\mathbf{D}(X, Y)-\sum_{S \in \mathcal{S}: S \subset X} \sum_{T \in \mathcal{T}: T \subset Y} \mathbf{A}_{\mathcal{S} \times \mathcal{T}}(S, T)\right| \\
& =\left|\mathbf{D}(X, Y)-\sum_{S \in \mathcal{S}: S \subset X} \sum_{T \in \mathcal{T}: T \subset Y} \mathbf{A}(S, T)\right|=|\mathbf{D}(X, Y)-\mathbf{A}(X, Y)| \leq\|\mathbf{D}-\mathbf{A}\|_{\square}
\end{aligned}
$$

thereby proving (12).

### 3.3 Proof of Corollary 2

Let $C>1$ and $0<\varepsilon<\frac{1}{2}$, let $\kappa, \gamma$ be as in (2) and suppose that $G=(V, E)$ is a $(C, \gamma)$-bounded graph on $n$ vertices $V=\{1, \ldots, n\}$ with adjacency matrix $\mathbf{A}$. The algorithm WeakPartition $(G, C, \varepsilon)$ calls PartMatrix $(\mathbf{A}, C, \varepsilon)$ to obtain two partitions $\mathcal{S}, \mathcal{T}$ of $V$ such that $\left\|\mathbf{A}_{\mathcal{S} \times \mathcal{T}}-\mathbf{A}\right\|_{\square} \leq 2 \varepsilon\|A\|_{\square}$. By Corollary 1 both $\mathcal{S}, \mathcal{T}$ have at most $2^{\kappa}$ classes. Then, the algorithm constructs the coarsest partition $\mathcal{V}$ of $V$ that is a refinement of both $\mathcal{S}$ and $\mathcal{T}$. Clearly, $|\mathcal{V}| \leq 2^{2 \kappa}$, and the running time of the algorithm is at most $2^{2 \kappa} \Pi(n)$ for some fixed polynomial $\Pi$.

To complete the proof, we need to show that $\left\|\mathbf{A}_{\mathcal{V} \times \mathcal{V}}-\mathbf{A}\right\|_{\square} \leq 4 \varepsilon\|\mathbf{A}\|_{\square}$. Since $\left\|\mathbf{A}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \leq 2 \varepsilon\|\mathbf{A}\|_{\square}$ by Corollary 1, we just need to prove that $\left\|\mathbf{A}_{\mathcal{V} \times \mathcal{V}}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \leq\left\|\mathbf{A}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square}$. To show this, we use a similar argument to that given in the proof of Corollary 1 . Namely, let $X, Y \subset V$ be such that $\left|\left(\mathbf{A}_{\mathcal{V} \times \mathcal{V}}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right)(X, Y)\right|=$ $\left\|\mathbf{A}_{\mathcal{V} \times \mathcal{V}}-\mathbf{A}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square}$. Since both $\mathbf{A}_{\mathcal{V} \times \mathcal{V}}$ and $\mathbf{A}_{\mathcal{S} \times \mathcal{T}}$ are constant on each square $S \times T$ with $S, T \in \mathcal{V}$, we may assume that $X, Y$ are unions of classes of $\mathcal{V}$. Therefore, the definition of $\mathbf{A}_{\mathcal{V} \times \mathcal{V}}$ entails that

$$
\left.\begin{aligned}
\left\|\mathbf{A}_{\mathcal{S} \times \mathcal{T}}-\mathbf{A}_{\mathcal{V} \times \mathcal{V}}\right\|_{\square} & =\left|\left(\mathbf{A}_{\mathcal{S} \times \mathcal{T}}-\mathbf{A}_{\mathcal{V} \times \mathcal{V}}\right)(X, Y)\right| \\
& =\left\lvert\, \begin{array}{|l|}
\mathbf{A}_{\mathcal{S} \times \mathcal{T}}(X, Y)-\sum_{S, T \in \mathcal{V}: S \times T \subset X \times Y} \\
\\
\end{array}\right. \\
& \mathbf{A}_{\mathcal{V} \times \mathcal{V}}(S, T) \mid \\
& \leq\left|\left(\mathbf{A}_{\mathcal{S} \times \mathcal{T}}-\mathbf{A}\right)(X, Y)\right| \leq\left\|\mathbf{A}_{\mathcal{S} \times \mathcal{T}}-\mathbf{A}\right\|_{\square}
\end{aligned} \right\rvert\,
$$

as claimed.

## 4 Approximating and partitioning $k$-dimensional tensors

### 4.1 Proof of Theorem 2

Let $0<\varepsilon<\frac{1}{2}$ and $C>1$. Let $\kappa=64(\zeta+100)(C / \varepsilon)^{2}$, where $\zeta$ is the constant from Theorem 1, and set $\gamma=2^{-3 \kappa}$. Throughout this section we assume that $\mathbf{A}: R_{1} \times \cdots \times R_{k} \rightarrow\{0,1\}$ is a $(C, \gamma)$-bounded tensor. Let $k_{1}=\lfloor k / 2\rfloor$.

## Algorithm 8 ApxTensor $(\mathbf{A}, C, \varepsilon)$

Input: A tensor A : $R_{1} \times \cdots \times R_{k} \rightarrow\{0,1\}$, numbers $C, \varepsilon>0$.
Output: A sequence of cut tensors.

1. Set up the matrix $\mathbf{B}=\mathbf{B}(\mathbf{A})$ as in (3) and let $p$ be the density of $\mathbf{B}$.
2. Call PartMatrix $(\mathbf{B}, C, \varepsilon / 8)$ to obtain partitions $\mathcal{S}$ of $R_{1} \times \cdots \times R_{k_{1}}$ and $\mathcal{T}$ of $R_{k_{1}+1} \times \cdots \times R_{k}$ (cf. Corollary 1).
3. Let $\hat{\mathbf{A}}: R_{1} \times \cdots \times R_{k} \rightarrow[0,1]$ be the tensor defined by

$$
\hat{\mathbf{A}}\left(i_{1}, \ldots, i_{k}\right)=\min \left\{1, \frac{\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\left(\left(i_{1}, \ldots, i_{k_{1}}\right),\left(i_{k_{1}+1}, \ldots, i_{k}\right)\right)}{C p}\right\}
$$

where $\mathbf{B}_{\mathcal{S} \times \mathcal{T}}$ is the approximation of $\mathbf{B}$ corresponding to the parition $\mathcal{S} \times \mathcal{T}$ (cf. Corollary 1 ).
4. Call FKTensor $(\hat{\mathbf{A}}, \varepsilon /(2 C))$ to obtain cut tensors $\mathbf{D}_{1}, \ldots, \mathbf{D}_{s}$.

Output the cut tensors $C p \cdot \mathbf{D}_{1}, \ldots, C p \cdot \mathbf{D}_{s}$.
The basic idea behind the algorithm ApxTensor for Theorem 2 is to transform the given sparse tensor $\mathbf{A}$ into a dense tensor $\hat{\mathbf{A}}$ and to apply the algorithm FKTensor from Theorem 5 to the latter. To obtain $\hat{\mathbf{A}}$, ApxTensor sets up the $\left|R_{1} \times \cdots \times R_{k_{1}}\right|$ by $\left|R_{k_{1}+1} \times \cdots \times R_{k}\right|$ matrix $\mathbf{B}(\mathbf{A})$ as in (3). As this matrix is $(C, \gamma)$-bounded by assumption, we can apply PartMatrix to obtain a cut norm approximation $\mathbf{B}_{\mathcal{S} \times \mathcal{T}}$ that is constant on rectangles $S \times T$ with $S \in \mathcal{S}, T \in \mathcal{T}$ and whose entries are in $[0,1]$. Then, $\hat{\mathbf{A}}$ is (basically) obtained by dividing $\mathbf{B}_{\mathcal{S} \times \mathcal{T}}$ by $C p$.

Finally, ApxTensor applies FKTensor to $\hat{\mathbf{A}}$ to obtain cut tensors $\mathbf{D}_{1}, \ldots, \mathbf{D}_{s}$, which of course need to get scaled by a factor $C p$ to get the desired approximation of $\mathbf{A}$.

The key step in the analysis is to show that $C p \hat{\mathbf{A}}$ is close to $\mathbf{A}$.
Lemma 5 We have $\|\mathbf{A}-C p \hat{\mathbf{A}}\|_{\square}<\varepsilon\|\mathbf{A}\|_{\square} / 2$.
Proof Let $m=\prod_{1 \leq i \leq k_{1}}\left|R_{i}\right|$ and $n=\prod_{k_{1}<i \leq k}\left|R_{i}\right|$. Moreover, let $\hat{\mathbf{B}}$ be the matrix defined by

$$
\hat{\mathbf{B}}\left(\left(i_{1}, \ldots, i_{k_{1}}\right),\left(i_{k_{1}+1}, \ldots, i_{k}\right)\right)=C p \hat{\mathbf{A}}\left(i_{1}, \ldots, i_{k}\right)
$$

Then

$$
\begin{align*}
\|\mathbf{A}-C p \hat{\mathbf{A}}\|_{\square} & \leq\|\mathbf{B}-\hat{\mathbf{B}}\|_{\square}  \tag{13}\\
& \leq\left\|\mathbf{B}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square}+\left\|\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \\
& \leq \frac{\varepsilon}{8}\|\mathbf{B}\|_{\square}+\left\|\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square}  \tag{14}\\
& =\frac{\varepsilon}{8}\|\mathbf{A}\|_{\square}+\left\|\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \tag{15}
\end{align*}
$$

Here (13) follows from the fact that the LHS is effectively the maximum over a subset of choices of the RHS. Equation (14) follows from our choice of $\mathcal{S}, \mathcal{T}$ and (15) follows because $\|\mathbf{A}\|_{\square}=\|\mathbf{B}\|_{\square}=$ the number of 1's in both matrices/tensors.

To bound $\left\|\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square}$, observe that

$$
\begin{equation*}
0 \leq \hat{\mathbf{B}}(i, j) \leq \mathbf{B}_{\mathcal{S} \times \mathcal{T}}(i, j) \text { for all } i \in \prod_{1 \leq a \leq k_{1}} R_{a}, j \in \prod_{k_{1}<b \leq k} R_{b} \tag{16}
\end{equation*}
$$

Moreover, let $S_{0}$ be the union of all classes $S \in \mathcal{S}$ such that $|S|<\gamma m$. Similarly, let $T_{0}$ be the union of all $T \in \mathcal{T}$ of size $|T|<\gamma n$. Then

$$
\begin{equation*}
\left|S_{0}\right| \leq 2^{\kappa} \gamma m<\varepsilon m / 100, \quad\left|T_{0}\right| \leq 2^{\kappa} \gamma n<\varepsilon n / 100 \tag{17}
\end{equation*}
$$

due to the upper bound on the number of classes in $\mathcal{S}, \mathcal{T}$ from Corollary 1 and our choice of $\gamma$. Further, we claim that

$$
\begin{equation*}
\hat{\mathbf{B}}(i, j)=\mathbf{B}_{\mathcal{S} \times \mathcal{T}}(i, j) \text { for all } i \in \prod_{1 \leq a \leq k_{1}} R_{a} \backslash S_{0}, j \in \prod_{k_{1}<b \leq k} R_{b} \backslash T_{0} \tag{18}
\end{equation*}
$$

To see this, consider $i \notin S_{0}, j \notin T_{0}$, and let $S \in \mathcal{S}, T \in \mathcal{T}$ be the classes such that $i \in S, j \in T$. Then by the construction of $S_{0}, T_{0}$ we have $|S| \geq \gamma m,|T| \geq \gamma n$. Therefore, the fact that $\mathbf{B}$ is $(C, \gamma)$-bounded implies that $\mathbf{B}(S, T) \leq C \cdot|S \times T| \cdot p$. Hence, $\mathbf{B}_{\mathcal{S} \times \mathcal{T}}(i, j)=\mathbf{B}(S, T) /|S \times T| \leq C p$. The definiton of $\hat{\mathbf{A}}$ in Step 3 of ApxTensor yields (18).

Finally, for any two sets $X \subset \prod_{1 \leq a i \leq k_{1}} R_{a}, Y \subset \prod_{k_{1}<b \leq k} R_{b}$ we obtain

$$
\begin{align*}
& \left.\mid\left(\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right)\right)(X, Y) \mid \\
& \leq\left|\left(\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right)\left(X \backslash S_{0}, Y \backslash T_{0}\right)\right|+\left|\left(\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right)\left(X \cap S_{0}, Y\right)\right|+\left|\left(\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right)\left(X, Y \cap T_{0}\right)\right| \\
& \stackrel{(18)}{=}\left|\left(\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right)\left(X \cap S_{0}, Y\right)\right|+\left|\left(\hat{\mathbf{B}}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right)\left(X, Y \cap T_{0}\right)\right| \\
& \stackrel{(16)}{\leq}\left|\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\left(X \cap S_{0}, Y\right)\right|+\left|\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\left(X, Y \cap T_{0}\right)\right| \\
& \leq\left|\mathbf{B}\left(X \cap S_{0}, Y\right)\right|+\left|\mathbf{B}\left(X, Y \cap T_{0}\right)\right|+2\left\|\mathbf{B}-\mathbf{B}_{\mathcal{S} \times \mathcal{T}}\right\|_{\square} \\
& \leq\left|\mathbf{B}\left(X \cap S_{0}, Y\right)\right|+\left|\mathbf{B}\left(X, Y \cap T_{0}\right)\right|+\frac{\varepsilon}{4}\|\mathbf{B}\|_{\square} . \tag{19}
\end{align*}
$$

As B is $(C, \gamma)$-bounded and $\left|X \cap S_{0}\right| \leq\left|S_{0}\right| \leq \varepsilon m / 100,\left|Y \cap T_{0}\right| \leq\left|T_{0}\right| \leq \varepsilon n / 100$ by (17), we have (see (7))

$$
\left|\mathbf{B}\left(X \cap S_{0}, Y\right)\right|,\left|\mathbf{B}\left(X, Y \cap T_{0}\right)\right| \leq \varepsilon\|B\|_{\square} / 100
$$

Consequently, (19) yields $|(\hat{\mathbf{B}}-\mathbf{B})(X, Y)| \leq 27 \varepsilon\|\mathbf{B}\|_{\square} / 100$. Since this is true for any $X, Y$, we conclude that $\|\hat{\mathbf{B}}-\mathbf{B}\|_{\square}<\varepsilon\|\mathbf{B}\|_{\square} / 3$. Plugging this bound into (15) completes the proof.
Proof of Theorem 2. The desired bound on the running time follows directly from Corollary 1 and Theorem 5. Moreover, since FKTensor is applied with $\delta=\varepsilon /(2 C)$, Theorem 5 implies that the number of cut tensors is at most $s \leq(\Gamma C / \varepsilon)^{2(k-1)}$ for a certain constant $\Gamma>0$. Furthermore, Lemma 5 and Theorem 5 yield

$$
\begin{aligned}
\left\|\mathbf{A}-C p \sum_{i=1}^{s} \mathbf{D}_{i}\right\|_{\square} & \leq\|\mathbf{A}-C p \hat{\mathbf{A}}\|_{\square}+C p \cdot\left\|\hat{\mathbf{A}}-\sum_{i=1}^{s} \mathbf{D}_{i}\right\|_{\square}<\frac{\varepsilon}{2}\|\mathbf{A}\|_{\square}+C p \cdot\left\|\hat{\mathbf{A}}-\sum_{i=1}^{s} \mathbf{D}_{i}\right\|_{\square} \\
& \leq \frac{\varepsilon}{2}\|\mathbf{A}\|_{\square}+C p \cdot \frac{\varepsilon}{2 C} \prod_{i=1}^{k}\left|R_{i}\right|=\varepsilon\|\mathbf{A}\|_{\square}
\end{aligned}
$$

as desired.

### 4.2 Proof of Corollary 3

Let $0<\varepsilon<\frac{1}{2}$ and $C>1$. Moreover, let $\Gamma$ be a sufficiently large constant. Suppose that $\mathbf{A}: R_{1} \times \cdots \times R_{k} \rightarrow\{0,1\}$ is a $(C, \gamma)$-bounded $0 / 1$ tensor. The algorithm PartTensor calls ApxTensor to obtain cut tensors

$$
\mathbf{D}_{i}=\operatorname{CUT}\left(d_{i}, S_{1 i}, \ldots, S_{k i}\right) \quad(1 \leq i \leq s)
$$

such that $\left\|\mathbf{A}-\sum_{i=1}^{s} \mathbf{D}_{i}\right\|_{\square} \leq \varepsilon\|A\|_{\square} / 2$. Then, for each $1 \leq j \leq k$ ApxTensor outputs the coarsest partition $\mathcal{R}_{j}$ of $R_{j}$ such that each of the sets $S_{j i}$ is a union of classes of $\mathcal{S}_{j}(i=1, \ldots, s)$.

Each partition $\mathcal{R}_{j}$ has at most $2^{s}$ classes. Hence, the upper bound on $s$ from Theorem 2 entails the upper bound on $\left|\mathcal{R}_{j}\right|$ stated in Corollary 3. Moreover, bound on the running time follows from Theorem 2 as well. Hence, we finally need to show that

$$
\left\|\mathbf{A}-\mathbf{A}_{\mathcal{R}_{1} \times \cdots \times \mathcal{R}_{k}}\right\|_{\square} \leq \varepsilon\|A\|_{\square}
$$

To simplify the notation we let $\mathbf{D}=\sum_{i=1}^{s} \mathbf{D}_{i}, \mathbf{B}=\mathbf{A}_{\mathcal{R}_{1} \times \cdots \times \mathcal{R}_{k}}$. We know that

$$
\begin{equation*}
\|\mathbf{D}-\mathbf{A}\|_{\square} \leq \varepsilon\|A\|_{\square} / 2 \tag{20}
\end{equation*}
$$

To complete the proof, we are going to show that $\|\mathbf{D}-\mathbf{B}\|_{\square} \leq \varepsilon\|A\|_{\square} / 2$ as well. Thus, let $X_{1} \subset R_{1}, \ldots, X_{k} \subset R_{k}$ be sets such that

$$
\|\mathbf{D}-\mathbf{B}\|_{\square}=\left|(\mathbf{D}-\mathbf{B})\left(X_{1}, \ldots, X_{k}\right)\right|
$$

Since both $\mathbf{D}$ and $\mathbf{B}$ are constant on any rectangle $S_{1} \times \cdots \times S_{k}$ with $S_{i} \in \mathcal{R}_{i}$, we may assume that $X_{i}$ is a union of classes of $\mathcal{R}_{i}$ for all $1 \leq i \leq k$. Furthermore, if $S_{i} \in \mathcal{R}_{i}$ for $1 \leq i \leq k$, then $\mathbf{B}\left(S_{1}, \ldots, S_{k}\right)=\mathbf{A}\left(S_{1}, \ldots, S_{k}\right)$. Therefore,

$$
\begin{aligned}
\|\mathbf{D}-\mathbf{B}\|_{\square} & =\left|(\mathbf{D}-\mathbf{B})\left(X_{1}, \ldots, X_{k}\right)\right|=\left|\mathbf{D}\left(X_{1}, \ldots, X_{k}\right)-\sum_{S_{1} \in \mathcal{R}_{1}: S_{1} \subset X_{1}} \ldots \sum_{S_{k} \in \mathcal{R}_{k}: S_{k} \subset X_{k}} \mathbf{B}\left(S_{1}, \ldots, S_{k}\right)\right| \\
& =\left|\mathbf{D}\left(X_{1}, \ldots, X_{k}\right)-\sum_{S_{1} \in \mathcal{R}_{1}: S_{1} \subset X_{1}} \ldots \sum_{S_{k} \in \mathcal{R}_{k}: S_{k} \subset X_{k}} \mathbf{A}\left(S_{1}, \ldots, S_{k}\right)\right| \\
& =\left|(\mathbf{D}-\mathbf{A})\left(X_{1}, \ldots, X_{k}\right)\right| \leq\|\mathbf{D}-\mathbf{A}\|_{\square} \leq \varepsilon\|\mathbf{A}\|_{\square} / 2,
\end{aligned}
$$

by (20). Hence, $\|\mathbf{A}-\mathbf{B}\|_{\square} \leq\|\mathbf{A}-\mathbf{D}\|_{\square}+\|\mathbf{D}-\mathbf{B}\|_{\square} \leq \varepsilon\|A\|_{\square}$, as desired.

## 5 Approximating MAX CSP problems

Throughout this section we keep the notation from Section 1.5. Given $0<\varepsilon<\frac{1}{2}, C>1$, we set $\gamma=\exp \left(-\Gamma(C / \varepsilon)^{2}\right)$, where $\Gamma$ is the constant from Theorem 2. Moreover, we assume that $\mathcal{F}$ is a $(C, \gamma)$-bounded $k$-CSP instance on $n$ variables $V=\{1, \ldots, n\}$, where $n>n_{0}$ for some sufficiently large number $n_{0}=n_{0}(C, \varepsilon, k)$. Let $m=|\mathcal{F}|$ be the number of constraints.

### 5.1 The algorithm ApxCSP

## Algorithm $9 \operatorname{ApxCSP}(\mathcal{F}, C, \varepsilon)$

Input: A $k$-CSP instance $\mathcal{F}$ over $V=\left\{x_{1}, \ldots, x_{n}\right\}$, numbers $C, \varepsilon>0$.
Output: An assignment $\hat{\sigma}: V \rightarrow\{0,1\}$.

1. Set up the tensors $\mathbf{A}_{\mathcal{F}}^{\psi}$ for all $\psi \in \Psi$.

Let $\alpha=\varepsilon 2^{-2^{k}-2 k-2}$.
Call ApxTensor $\left(\mathbf{A}_{\mathcal{F}}^{\dot{\psi}}, C, \alpha\right)$ for each $\psi \in \Psi$ to obtain tensors

$$
\mathbf{B}^{\psi}=\sum_{i=1}^{s} \mathbf{D}_{i}^{\psi}, \text { where } \mathbf{D}_{i}^{\psi}=\operatorname{CUT}\left(d_{i}^{\psi}, S_{i 1}^{\psi}, \ldots, S_{i k}^{\psi}\right) .
$$

Let $\mathcal{P}$ be the coarsest partition of $V$ such that each set $S_{i h}^{\psi}$ is a union of classes of $\mathcal{P}$ ( $1 \leq i \leq s, 1 \leq h \leq k$, $\psi \in \Psi)$.
2. Let $\delta=C^{-1} \Gamma^{-k} s^{-1} 2^{-2^{k}-4 k-4}$ and $\nu=\lceil\delta n\rceil$.

Compute an optimal solution $\left(\hat{\tau}_{i h}^{\psi}(1), \hat{\tau}_{i h}^{\psi}(0), \hat{z}_{P}\right)_{i \in[s], h \in[k], \psi \in \Psi, P \in \mathcal{P}}$ to the following optimization problem.

$$
\begin{aligned}
\text { OPT }^{\prime \prime}=\max & \sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{y \in\{0,1\}^{k}} d_{i}^{\psi} \psi(y) \nu^{k} \prod_{h=1}^{k} \tau_{i h}^{\psi}\left(y_{h}\right) \\
\text { s.t. } & 0 \leq \tau_{i h}^{\psi}(1) \leq\left\lfloor\left|S_{i h}^{\psi}\right| / \nu\right\rfloor \quad \text { is an integer for all } 1 \leq i \leq s, 1 \leq h \leq k, \psi \in \Psi, \\
& \tau_{i h}^{\psi}(0)=\nu^{-1} S_{i h}^{\psi}-\tau_{i h}^{\psi}(1) \quad \text { for all } 1 \leq i \leq s, 1 \leq h \leq k, \psi \in \Psi \\
& \tau_{i h}^{\psi}(1) \nu \leq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} z_{P} \leq\left(\tau_{i h}^{\psi}(1)+1\right) \nu \quad \text { for all } 1 \leq i \leq s, 1 \leq h \leq k, \psi \in \Psi \\
& 0 \leq z_{P} \leq|P| \quad \text { for all } P \in \mathcal{P} .
\end{aligned}
$$

(The numbers $z_{P}$ are not required to be integers.)
Output an assignment $\hat{\sigma}: V \rightarrow\{0,1\}$ such that $\left|\left|\hat{\sigma}^{-1}(1) \cap P\right|-\hat{z}_{P}\right| \leq 1$ for all $P \in \mathcal{P}$.
The first step of ApxCSP relies on the procedure ApxTensor from Theorem 2. Since we assume that all the tensors $\mathbf{A}_{\mathcal{F}}^{\psi}$ are $(C, \gamma)$-bounded, we can apply ApxTensor to each of them to obtain an approximation $\mathbf{B}^{\psi}$ consisting of a bounded number of cut tensors $\mathbf{D}_{i}^{\psi}$. The basic idea is to approximate the MAX CSP problem, i.e., the optimization problem

$$
\mathrm{OPT}=\max _{\sigma \in\{0,1\}^{V}} \sum_{\phi \in \mathcal{F}} \phi(\sigma)=\max _{\sigma \in\{0,1\}^{V}} \sum_{\psi \in \Psi} \sum_{\left(z_{1}, \ldots, z_{k}\right) \in V^{k}} \mathbf{A}^{\psi}\left(z_{1}, \ldots, z_{k}\right) \psi\left(\sigma\left(z_{1}\right), \ldots, \sigma\left(z_{k}\right)\right)
$$

by the optimization problem

$$
\mathrm{OPT}^{\prime}=\max _{\sigma \in\{0,1\}^{V}} \sum_{\psi \in \Psi} \sum_{\left(z_{1}, \ldots, z_{k}\right) \in V^{k}} \mathbf{B}^{\psi}\left(z_{1}, \ldots, z_{k}\right) \psi\left(\sigma\left(z_{1}\right), \ldots, \sigma\left(z_{k}\right)\right)
$$

The following lemma, whose proof we defer to Section 5.2, shows that any assignment $\sigma$ that approximates $\mathrm{OPT}^{\prime}$ well also provides a good approximation for OPT.

Lemma 6 Let $\sigma \in\{0,1\}^{V}$ be such that

$$
\begin{equation*}
\sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\sigma(z)) \geq O P T^{\prime}-2^{-k-1} \varepsilon m \tag{21}
\end{equation*}
$$

Then $\sum_{\phi \in \mathcal{F}} \phi(\sigma) \geq(1-\varepsilon) O P T(\mathcal{F})$.
It is worth pointing out that $\mathrm{OPT}^{\prime}$ can be solved exactly in "polynomial" time. This is because the tensors $\mathbf{B}^{\psi}$ consist of only a bounded (w.r.t. $n$ ) number of cut tensors. More precisely, the partition $\mathcal{P}$ constructed in Step 1 of the algorithm has the following property: if $S_{1}, \ldots, S_{k} \in \mathcal{P}$, then all the tensors $\mathbf{B}^{\psi}, \psi \in \Psi$, are constant on the rectangle $S_{1} \times \cdots \times S_{k}$. Therefore, as far as OPT $^{\prime}$ is concerned, the individual variables in each set $S \in \mathcal{P}$ are completely indistinguishable. More precisely, consider an assignment $\sigma: V \rightarrow\{0,1\}$ and let

$$
\begin{align*}
\mathcal{Z}_{P} & =|\{v \in P: \sigma(v)=1\}| \text { for each } P \in \mathcal{P}  \tag{22}\\
\mathcal{T}_{i h}^{\psi}(1) & =\sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} \mathcal{Z}_{P}, \mathcal{T}_{i h}^{\psi}(0)=\left|S_{i h}^{\psi}\right|-\mathcal{T}_{i h}^{\psi}(1) \quad \text { for } 1 \leq i \leq s, 1 \leq h \leq k, \psi \in \Psi \tag{23}
\end{align*}
$$

In words, $\mathcal{T}_{i h}^{\psi}(y)$ is the number of variables in $S_{i h}^{\psi}$ that attain the value $y$ under $\sigma(y=0,1)$. Let us further define

$$
\sigma(y)=\left(\sigma\left(y_{1}\right), \ldots, \sigma\left(y_{k}\right)\right) \quad \text { for } y \in V^{k}
$$

Then

$$
\sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\sigma(z))=\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{z \in \prod_{h=1}^{k} S_{i h}^{\psi}} d_{i}^{\psi} \psi(\sigma(z))=\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{y \in\{0,1\}^{k}} d_{i}^{\psi} \psi(y) \prod_{h=1}^{k} \mathcal{T}_{i h}^{\psi}\left(y_{h}\right)
$$

Hence, to solve $\mathrm{OPT}^{\prime}$ optimally, we could just try all possible tuples $\left(\mathcal{Z}_{P}\right)_{P \in \mathcal{P}}$ such that $0 \leq \mathcal{Z}_{P} \leq|P|$ is an integer. Since the number of such tuples is at most $n^{|\mathcal{P}|}$ and the number $|\mathcal{P}|$ of classes is independent of $n$, this yields a polynomial time algorithm for any fixed $\varepsilon, k, C$.

To speed things up, we use an idea developed in [16] for dense MAX CSP problems; this will eventually lead to the problem $\mathrm{OPT}^{\prime \prime}$ detailed in Step 2 of ApxCSP. The basic idea is the following. Instead of optimizing over all possible $\left(\mathcal{Z}_{P}\right)_{P \in \mathcal{P}}$, we could just enumerate all tuples $\left(\mathcal{T}_{i h}^{\psi}(1)\right)_{i, h, \psi}$ with $0 \leq \mathcal{T}_{i h}^{\psi}(1) \leq\left|S_{i h}^{\psi}\right|$. The issue is that not all such tuples correspond to an assignment $\sigma: V \rightarrow\{0,1\}$ as in (22) and (23). Hence, for each tuple $\left(\mathcal{T}_{i h}^{\psi}(1)\right)_{i, h, \psi}$ we will have to check feasbility, i.e., if there is a tuple $\left(\mathcal{Z}_{P}\right)_{P}$ such that (23) holds. Since we are just aiming to solve $\mathrm{OPT}^{\prime}$ approximately, we can drop the requirement that all $\mathcal{Z}_{P}$ must be integeral. Thus, checking (23) turns into a linear programming problem. In effect, we can reduce the running time from $\exp (|\mathcal{P}| \cdot \ln n)$ to $\exp \left(s k 2^{2^{k}} \cdot \ln n\right)$. (Remember that in general $|\mathcal{P}|$ is exponential in $s k 2^{2^{k}}$.)

Finally, to remove the $\ln n$ factor, we chop each set $S_{i h}^{\psi}$ into chunks of size $\nu=\lceil\delta n\rceil$, where $\delta>0$ is bounded by a function of $C, \varepsilon, k$ only. Hence, instead of optimizing over the number $0 \leq \mathcal{T}_{i h}^{\psi}(1) \leq\left|S_{i h}^{\psi}\right|$ of variables to be set to 1 in each $S_{i h}^{\psi}$, we optimize over the number $0 \leq \tau_{i h}^{\psi}(1)=\left\lfloor\mathcal{T}_{i h}^{\psi}(1) / \nu\right\rfloor \leq\left\lfloor S_{i h}^{\psi} / \nu\right\rfloor$ of chunks set to 1 . This is sufficient because we just need to solve $\mathrm{OPT}^{\prime}$ within an additive $\varepsilon 2^{-k-1} m$ (cf. (21)). Of course, for each $\tau_{i h}^{\psi}(1)$ the number of possible values is at most $1+\delta^{-1}$, i.e., independent of $n$. To check feasibility, we then have to verify that there are $0 \leq z_{P} \leq 1(P \in \mathcal{P})$ such that $\tau_{i h}^{\psi}(1) \nu \leq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} z_{P} \leq\left(\tau_{i h}^{\psi}(1)+1\right) \nu$ for all $i, h, \psi$, which is again an LP problem. This leaves us with the optimization problem $\mathrm{OPT}^{\prime \prime}$ quoted in Step 2 of ApxCSP. After finding an optimal solution to $\mathrm{OPT}^{\prime \prime}$, the algorithm sets up the assignment $\hat{\sigma}$ that mirrors the resulting $z_{P}$ values. We defer the proof of the following proposition to Section 5.3

Proposition 2 The assigment $\hat{\sigma}$ satisfies (21).
Proof of Theorem 3. The fact that the assignment $\hat{\sigma}$ computed by ApxCSP satisfies at least $(1-\varepsilon)$ OPT constraints follows from Lemma 6 and Proposition 2. Thus, we finally need to analyze the running time. By Theorem 2 the running time of Step 1 is at most

$$
2^{2^{k}} 2^{(C / \alpha)^{2}}(C / \alpha)^{3 k} \Pi^{\prime}\left(n^{k}\right)
$$

for some polynomial $\Pi^{\prime}$. Moreover, for each $\psi \in \Psi$ the resulting decomposition of $\mathbf{A}_{\mathcal{F}}^{\psi}$ consists of $s \leq(\Gamma C / \alpha)^{2(k-1)}$ cut matrices for some constant $\Gamma>0$. Hence,

$$
|\mathcal{P}| \leq \Lambda=2^{k 2^{2^{k}} s}
$$

Step 2 solves $\mathrm{OPT}^{\prime \prime}$ by enumerating all possible values for the integer variables $\tau_{i h}^{\psi}(1)$. The number of these integer variables is $s k 2^{2^{k}}$. Futhermore, for each $\tau_{i h}^{\psi}(1)$ there are at most $1+\left|S_{i h}^{\psi}\right| / \nu \leq 1+\delta^{-1} \leq \delta^{-2}$ values to consider. Therefore, the total number of possibilities that Step 2 enumerates over is at most

$$
\Lambda^{\prime}=\exp \left(-2 s k 2^{2^{k}} \ln \delta\right)
$$

For each of these choices we need to check the feasibility of the system of linear inequalities

$$
\tau_{i h}^{\psi}(1) \nu \leq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} z_{P} \leq\left(\tau_{i h}^{\psi}(1)+1\right) \nu, \quad 0 \leq z_{P} \leq|P| \quad(1 \leq i \leq s, 1 \leq h \leq k, \psi \in \Psi, P \in \mathcal{P})
$$

This can be performed in time polynomial in the encoding length of these linear equations. There are $|\mathcal{P}|+k 2^{2^{k}} s \leq 2 \Lambda$ constraints and $|\mathcal{P}| \leq \Lambda$ variables. The encoding length of the numbers involved is at most $\ln (n / \delta)$. Therefore, the running time is

$$
\text { poly }(|\Lambda| \ln (n / \delta))
$$

Consequently, the total running time is at most

$$
\Pi\left(\exp \left(-s k 2^{2^{k}} \ln \delta\right) \cdot n^{k}\right) \leq \Pi\left(\exp \left(k \Gamma^{k} 2^{2^{k}}(C / \varepsilon)^{2 k} \ln (C / \varepsilon)\right) \cdot n^{k}\right)
$$

for some fixed polynomial $\Pi$ and a constant $\Gamma>0$.

### 5.2 Proof of Lemma 6

We shall prove below that for any $\tau \in\{0,1\}^{V}$

$$
\begin{equation*}
\left|\sum_{\phi \in \mathcal{F}} \phi(\tau)-\sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\tau(z))\right| \leq 2^{-k-2} \varepsilon m \tag{24}
\end{equation*}
$$

This implies the assertion as follows. Since (24) holds for any $\tau \in\{0,1\}^{V}$, we have $\mathrm{OPT}^{\prime} \geq \mathrm{OPT}(\mathcal{F})-2^{-k-2} \varepsilon m$. Hence, if $\sigma \in\{0,1\}^{V}$ satisfies (21), then (24) yields

$$
\begin{equation*}
\sum_{\phi \in \mathcal{F}} \phi(\sigma) \geq \sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\sigma(z))-2^{-k-2} \varepsilon m \geq \mathrm{OPT}^{\prime}-\frac{3}{4} \cdot 2^{-k} \varepsilon m \geq \mathrm{OPT}(\mathcal{F})-2^{-k} \varepsilon m \tag{25}
\end{equation*}
$$

Finally, as for a random assignment $\tau \in\{0,1\}^{V}$ we have

$$
\mathrm{E}\left[\sum_{\phi \in \mathcal{F}} \phi(\tau)\right] \geq 2^{-k} m
$$

we conclude that $\operatorname{OPT}(\mathcal{F}) \geq 2^{-k} m$. Hence, the assertion follows from (25).
To prove (24), we fix $\tau \in\{0,1\}^{V}$ and let $D=\sum_{\psi \in \Psi}\left|\sum_{z \in V^{k}}\left(\mathbf{A}^{\psi}-\mathbf{B}^{\psi}\right)(z) \psi(\tau(z))\right|$, so that

$$
\begin{equation*}
\left|\sum_{\phi \in \mathcal{F}} \phi(\tau)-\sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\tau(z))\right|=\left|\sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{A}^{\psi}(z) \psi(\tau(z))-\sum_{\psi \in \Psi} \sum_{z \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\tau(z))\right| \leq D \tag{26}
\end{equation*}
$$

In order to estimate $D$, let $T_{\psi}=\psi^{-1}(1) \subset\{0,1\}^{k}$ for any $\psi \in \Psi$. Moreover, for any $t=\left(t_{1}, \ldots, t_{k}\right) \in\{0,1\}^{k}$ define $R(t)=\prod_{i=1}^{k} \sigma^{-1}\left(t_{i}\right) \subset V^{k}$. Then

$$
\left|\sum_{z \in R(t)}\left(\mathbf{A}^{\psi}-\mathbf{B}^{\psi}\right)(z)\right| \leq\left\|\mathbf{A}^{\psi}-\mathbf{B}^{\psi}\right\|_{\square}
$$

Therefore, the bound $\left\|\mathbf{A}^{\psi}-\mathbf{B}^{\psi}\right\|_{\square} \leq \alpha m$ entails

$$
D=\sum_{\psi \in \Psi}\left|\sum_{t \in T_{\psi}} \sum_{z \in R(t)}\left(\mathbf{A}^{\psi}-\mathbf{B}^{\psi}\right)(z)\right| \leq \sum_{\psi \in \Psi} \sum_{t \in T_{\psi}}\left\|\mathbf{A}^{\psi}-\mathbf{B}^{\psi}\right\|_{\square} \leq 2^{2^{k}+k} \alpha m \leq 2^{-k-2} \varepsilon m .
$$

Thus, (24) follows from (26).

### 5.3 Proof of Proposition 2

Lemma 7 We have $O P T^{\prime}-O P T^{\prime \prime} \leq 2^{-k-2} \varepsilon m$.
Proof Given an assignment $\sigma \in\{0,1\}^{V}$, we let $\theta_{i h}^{\psi}(y)=\left|S_{i h}^{\psi} \cap \sigma^{-1}(y)\right|$ for all $\psi, i, h$ and $y=0,1$. Then

$$
\begin{align*}
\sum_{\psi \in \Psi} \sum_{v \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\sigma(z)) & =\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{z \in V^{k}} \mathbf{D}_{i}^{\psi}(z) \psi(\sigma(z)) \\
& =\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{z \in \prod_{h=1}^{k} S_{i h}^{\psi}} d_{i}^{\psi} \psi(\sigma(z))=\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{y \in\{0,1\}^{k}} d_{i}^{\psi} \psi(y) \prod_{h=1}^{k} \theta_{i h}^{\psi}\left(y_{h}\right) \tag{27}
\end{align*}
$$

We obtain a feasible solution to $\mathrm{OPT}^{\prime \prime}$ by letting $z_{P}=\left|\sigma^{-1}(1) \cap P\right|$ for all $P \in \mathcal{P}, \tau_{i h}^{\psi}(1)=\left\lceil\theta_{i h}^{\psi}(1) / \nu\right\rceil$, and $\tau_{i h}^{\psi}(0)=\nu^{-1}\left|S_{i h}^{\psi}\right|-\tau_{i h}^{\psi}(1)(1 \leq i \leq s, 1 \leq h \leq k, \psi \in \Psi)$. To complete the proof, we shall compare the objective function value attained by this solution with (27). To this end, observe that $\left|\theta_{i h}^{\psi}(y)-\tau_{i h}^{\psi}(y) \nu\right| \leq \nu$ for all $i, h, \psi, y$. Therefore,

$$
\begin{equation*}
\left|\prod_{h=1}^{k} \theta_{i h}^{\psi}\left(y_{h}\right)-\prod_{h=1}^{k} \tau_{i h}^{\psi}\left(y_{h}\right) \nu\right| \leq 2^{k} \nu n^{k-1} \quad \text { for all } i, y, \psi \tag{28}
\end{equation*}
$$

Since by Theorem 2 we have $\left|d_{i}^{\psi}\right| \leq C p \Gamma^{k}$ for all $i, \psi$, (28) yields

$$
\begin{align*}
\left|\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{y \in\{0,1\}^{k}} d_{i}^{\psi} \psi(y)\left[\prod_{h=1}^{k} \theta_{i h}^{\psi}\left(y_{h}\right)-\nu^{k} \prod_{h=1}^{k} \tau_{i h}^{\psi}\left(y_{h}\right)\right]\right| & \leq C p \Gamma^{k} \cdot s 2^{2^{k}+2 k} \nu n^{k-1} \\
& \leq C p \Gamma^{k} \cdot \delta s 2^{2^{k}+2 k+1} n^{k} \leq 2^{-k-2} \varepsilon m \tag{29}
\end{align*}
$$

by our choice of $\delta$. Finally, combining (27) and (29), we conclude that $\mathrm{OPT}^{\prime \prime} \geq \mathrm{OPT}^{\prime}-2^{-k-2} \varepsilon m$, as desired.
Proof of Proposition 2. Letting $\theta_{i h}^{\psi}(y)=\left|S_{i h}^{\psi} \cap \sigma^{-1}(y)\right|$ for all $\psi, i, h$ and $y=0,1$, we have

$$
\sum_{\psi \in \Psi} \sum_{v \in V^{k}} \mathbf{B}^{\psi}(z) \psi(\sigma(z))=\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{y \in\{0,1\}^{k}} d_{i}^{\psi} \psi(y) \prod_{h=1}^{k} \theta_{i h}^{\psi}\left(y_{h}\right)
$$

Furthermore, since $\left|z_{P}-\left|\sigma^{-1}(1) \cap P\right|\right| \leq 1$, we have

$$
\tau_{i h}^{\psi}(1) \nu \leq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} z_{P} \leq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} 1+\left|\sigma^{-1}(1) \cap P\right| \leq\left|\sigma^{-1}(1) \cap S_{i h}^{\psi}\right|+s \leq\left|S_{i h}^{\psi}\right|+\nu
$$

the last inequality follows from our assumption that $n>n_{0}=2 s / \delta$. Similarly,

$$
\left(1+\tau_{i h}^{\psi}(1)\right) \nu \geq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}} z_{P} \geq \sum_{P \in \mathcal{P}: P \subset S_{i h}^{\psi}}\left|\sigma^{-1}(1) \cap P\right|-1 \geq\left|\sigma^{-1}(1) \cap S_{i h}^{\psi}\right|-s \geq\left|S_{i h}^{\psi}\right|-\nu
$$

Hence, $\left|\tau_{i h}^{\psi}(y) \nu-\theta_{i h}^{\psi}(y)\right| \leq 2 \nu$. Consequently,

$$
\begin{equation*}
\left|\prod_{h=1}^{k} \theta_{i h}^{\psi}\left(y_{h}\right)-\prod_{h=1}^{k} \tau_{i h}^{\psi}\left(y_{h}\right) \nu\right| \leq 2^{2 k} \nu n^{k-1} \quad \text { for all } i, y, \psi \tag{30}
\end{equation*}
$$

As $\left|d_{i}^{\psi}\right| \leq C p \Gamma^{k}$ for all $i, \psi$ by Theorem 2, (30) yields

$$
\begin{equation*}
\left|\sum_{\psi \in \Psi} \sum_{i=1}^{s} \sum_{y \in\{0,1\}^{k}} d_{i}^{\psi} \psi(y)\left[\prod_{h=1}^{k} \theta_{i h}^{\psi}\left(y_{h}\right)-\nu^{k} \prod_{h=1}^{k} \tau_{i h}^{\psi}\left(y_{h}\right)\right]\right| \leq C p \Gamma^{k} \cdot s 2^{2^{k}+3 k} \nu n^{k-1} \leq 2^{-k-2} \varepsilon m \tag{31}
\end{equation*}
$$

by the definition of $\delta$. Thus, (30) and (31) yield the assertion

## 6 Examples

We present a few examples of bounded problem instances of MAX CUT and (MAX) $k$-SAT. The present techniques provide a unified approach to problems that were previously studied by individually tailored methods. The first two examples demonstrate how our results can be used to generalize average case analyses of algorithms. In the third instance we show how our techniques complement a prior result on "planted" random 3-SAT.

### 6.1 MAX CUT

Let $0 \leq p=p(n) \leq 1$ be a sequence of edge probabilities, and let $G_{n, p}$ be a random graph on $n$ vertices $V=$ $\{1, \ldots, n\}$ obtained by including each of the $\binom{n}{2}$ possible edges with probability $p$ independently. We say that $G_{n, p}$ has some property $\mathcal{E}$ with high probability ("w.h.p.") if the probability that $\mathcal{E}$ holds tends to 1 as $n \rightarrow \infty$. For any graph $G$ we let $\mathcal{I}(G)$ denote the set of all subgraphs $H$ of $G$ such that $|E(H)| \geq 0.01|E(G)|$. Furthermore, for a fixed $\varepsilon>0$ we say that an algorithm $\mathcal{A}$ approximates MAX CUT within $1-\varepsilon$ on $G_{n, p}$-bounded graphs if the following two conditions are satisfied:

1. For any input graph $G$ the algorithm $\mathcal{A}$ either outputs a cut that is within a $1-\varepsilon$ factor of the maximum cut, or just outputs "fail". In the first case we say that the algorithm succeeds, in the second case it fails.
2. If $G=G_{n, p}$ is a random graph, then with high probability $\mathcal{A}$ succeeds for all graphs in $\mathcal{I}(G)$.

Thus, the algorithm never outputs a solution that is off by more than $1-\varepsilon$, and for almost all outcomes $G=G_{n, p}$ it succeeds on all subgraphs $G_{*} \subset G$ that contain at least $1 \%$ of the edges of $G$. One can think of $G_{*}$ being constructed by a malicious adversary, starting from the random graph $G$.

Theorem 10 Suppose that $n p \geq c_{0}(\varepsilon)$ for a number $c_{0}(\varepsilon)$ that only depends on $\varepsilon>0$. The polynomial time algorithm ApxCSP from Theorem 3 approximates MAX CUT within $1-\varepsilon$ on $G_{n, p}$-bounded graphs.

Proof MAX CUT fits into the general CSP framework discussed in Section 1.5 as follows. The set of variables is the vertex set of the input graph $G_{*}=\left(V, E_{*}\right)$. Moreover, each edge $e=\{v, w\} \in E_{*}$ yields the (binary) contraint

$$
\sigma \in\{0,1\}^{V} \mapsto \begin{cases}1 & \text { if } \sigma(v) \neq \sigma(w) \\ 0 & \text { otherwise }\end{cases}
$$

Thus, the objective function value of the resulting CSP $\mathcal{F}$ is just the number of crossing edges of a maximum cut of $G_{*}$.

Let $\varepsilon>0$, and let $\gamma$ be as in Theorem 3 with $C=360$. We claim that if $n p \geq c_{0}(\varepsilon)$ for a sufficiently large $c_{0}(\varepsilon)>0$, then whp $G=G_{n, p}$ has the property that for any $G_{*} \in \mathcal{I}(G)$ the CSP instance $\mathcal{F}$ is $(360, \gamma)$-bounded. By the construction of $\mathcal{F}$, it is sufficient to show that the adjacency matrix $A=A\left(G_{*}\right)$ is $(180, \gamma)$-bounded. To see this, consider any two sets $S, T \subset V$ of sizes $|S|,|T| \geq \gamma n$. Then

$$
A(S, T) \leq 2 e_{G_{*}}(S, T) \leq 2 e_{G}(S, T)
$$

Since $G=G(n, p)$ is a random graph, we have $\mathrm{E}\left(2 e_{G}(S, T)\right) \leq 2|S \times T| p$. Moreover, as $e_{G}(S, T)$ is binomially distributed, Chernoff bounds entail that

$$
\mathrm{P}\left[2 e_{G}(S, T)>3|S \times T| p\right] \leq \exp (-0.01|S \times T| p) \leq \exp \left(-0.01 \cdot \gamma^{2} n^{2} p\right) \leq \exp \left(-0.01 \gamma^{2} c_{0}(\varepsilon) \cdot n\right)
$$

Hence, if $c_{0}(\varepsilon)$ is sufficiently large, then $A(S, T) \leq 3|S \times T| p$ with probability at least $1-\exp (-2 n)$. Since there are at most $2^{n}$ ways to choose $S, T$, the union bound entails that whp for all pairs of sets $S, T$ of size at least $\gamma n$ we have

$$
\begin{equation*}
A(S, T) \leq 3|S \times T| p \tag{32}
\end{equation*}
$$

Finally, let $q$ be the density of $A$. Since the number of edges of $G(n, p)$ is $(1+o(1))\binom{n}{2} p$ whp (by Chernoff bounds), and since $G_{*} \in \mathcal{I}(G)$, we have $0.009 p \leq q$ whp. Hence, (32) entails that $A(S, T) \leq 180|S \times T| q$ for all $S, T$ of size at least $\gamma n \mathbf{w h p}$, i.e., $A$ is $(180, \gamma)$-bounded.

Theorem 10 readily yields a result on the "planted model" for MAX CUT. In this model a random graph $G=$ $G_{n, p, q}$ is generated by partitioning the vertex set $V=\{1, \ldots, n\}$ randomly into two parts $V_{1}, V_{2}$, inserting each possible $V_{1}-V_{2}$-edge with probability $p$, and each possible edge inside $V_{1}, V_{2}$ with probability $q<p$ independently. Improving upon prior work by Boppana [6], Coja-Oghlan [7] showed that a MAX CUT of $G_{n, p, q}$ can be computed in polynomial time whp, provided that $n(p-q) \geq \zeta \sqrt{n p \ln (n p)}$ for a certain constant $\zeta>0$ (actually [6, 7] are stated in terms of MIN BISECTION, but things carries over to MAX CUT easily). Since the random graph $G_{n, p, q}$ can be obtained by first choosing $G_{n, p}$, then choosing a random partition ( $V_{1}, V_{2}$ ), and finally removing random edges inside of $V_{1}, V_{2}$, Theorem 10 encompasses this model. In fact, Theorem 10 comprises various generalizations of the "planted cut" model (e.g., instead of planting a single cut, we could plant an arbitrary number of cuts, etc.).

### 6.2 MAX $k$-SAT

Let $V=\left\{x_{1}, \ldots, x_{n}\right\}$ be a set of $n$ propositional variables, and let $F_{k}(n, p)$ signify a $k$-SAT formula obtained by including each of the $(2 n)^{k}$ possible $k$-clauses with probability $0 \leq p \leq 1$ independently (hence, we think of each clause as an order $k$-tuple of literals). Let $m=(2 n)^{k} p$ denote the expected number of clauses. We say that $F_{k}(n, p)$ has some property $\mathcal{E}$ with high probability if the probability that $\mathcal{E}$ holds tends to one as $n \rightarrow \infty$. Moreover, for any $k$-SAT formula $F$ we let $\mathcal{I}(F)$ denote the set of all sub-formulas $F_{*}$ of $F$ that contain at least $0.01 m$ clauses. Furthermore, for a fixed $\varepsilon>0$ we say that an algorithm $\mathcal{A}$ approximates $M A X k$-SAT within $1-\varepsilon$ on $F_{k}(n, p)$-bounded formulas if the following two conditions are satisfied:

1. For any input $F$ the algorithm $\mathcal{A}$ either outputs an assignment such that the number of satisfied clauses is within a $1-\varepsilon$ factor of the optimum for MAX $k$-SAT or just outputs "fail".
2. If $F=F_{k}(n, p)$, then whp $\mathcal{A}$ succeeds on all formulas in $\mathcal{I}(F)$.

Theorem 11 Suppose that $k \geq 2$ is fixed and that $c_{0}(\varepsilon) n^{\lceil k / 2\rceil} \leq m=o\left(n^{k}\right)$ for a number $c_{0}(\varepsilon)$ that only depends on $\varepsilon$. The polynomial time algorithm ApxCSP from Theorem 3 approximates MAX $k$-SAT within $1-\varepsilon$ on $F_{k}(n, p)$ bounded formulas.

Proof Let $\mathcal{F}=F_{k}(n, p)$ be a random $k$-SAT formula. Then the problem of finding an assignment that maximizes the number of simultaneously satisfied clauses can be stated as a MAX CSP problem in the sense of Section 1.5 as follows. Each clause $l_{1} \vee \cdots \vee l_{k}$ of $\mathcal{F}$ yields Boolean function as follows. Let $s_{i}=1$ if $l_{i}$ is just a variable $y_{i}$, and $s_{i}=-1$ if $l_{i}$ is the negation of a variable $y_{i}$. Then the clause yields the function

$$
\sigma \in\{0,1\}^{V} \mapsto \max _{i=1, \ldots, k} \frac{1+2 \sigma\left(y_{i}\right) s_{i}-s_{i}}{2} \in\{0,1\}
$$

Hence, for at most $2^{k}$ functions $\psi \in \Psi$ the tensor $\mathbf{A}_{\mathcal{F}}^{\psi}$ is non-zero, and each of these $2^{k}$ functions corresponds to one way of choosing the signs $\left(s_{1}, \ldots, s_{k}\right)$.

Furthermore, the tensors $\mathbf{A}_{\mathcal{F}}^{\psi}$ corresponding to a sequence $\left(s_{1}, \ldots, s_{k}\right)$ of signs are random. More precisely, for any tuple of indices $1 \leq i_{1}<\cdots<i_{k} \leq n$ the entry of $\mathbf{A}_{\mathcal{F}}^{\psi}\left(i_{1}, \ldots, i_{k}\right)$ is one iff one of the $k$ ! clauses corresponding to any permutation of the variables $x_{i_{1}} \vee \cdots \vee x_{i_{k}}$ occurs in $\mathcal{F}$; since we are assuming that $m=o\left(n^{k}\right)$, the probability that more than one of these clauses occurs is $o(1)$. Hence, the entries of $\mathbf{A}_{\mathcal{F}}^{\psi}$ are mutually independent random variables. Therefore, similarly as in the proof of Theorem 10 Chernoff bounds show that $\mathbf{A}_{\mathcal{F}}^{\psi}$ is $(1000, \gamma)$-bounded for any fixed $\gamma>0$ whp.

In particular, Theorem 11 applies to plainly random formula $F_{k}(n, p)$, in which case the algorithm yields a lower and an upper bound on the number of simultaneously satisfiable clauses. If $k \geq 3$, then for $m \geq c_{0}(\varepsilon) n^{\lceil k / 2\rceil}$ the optimal assignment of $F_{k}(n, p)$ satisfies a $1-2^{-k}+o(1)$ fraction of the clauses whp (by a standard first moment argument). Hence, w.h.p. the polynomial time algorithm ApxSAT yields a proof that there is no assignment satisfying more than a $1-2^{-k}+\varepsilon$ fraction of all clauses. The problem of deriving such a proof in polynomial time is known as the "strong refutation problem" for random $k$-SAT (cf. Feige [9]), and a number of authors have tailored algorithms specifically for this problem [8, 10, 15]. For even values of $k$, Theorem 10 matches the best known result [8].

### 6.3 Planted 3-SAT

Througout this section we let $\delta>0$ be a sufficiently small and $\zeta>0$ a sufficiently large constant; their precise values will be specified implicitly in due course.

Consider the following model of random 3-SAT. Let $V=\left\{x_{1}, \ldots, x_{n}\right\}$ be a set of Boolean variables, and let $L=\left\{x_{1}, \bar{x}_{1}, \ldots, x_{n}, \bar{x}_{n}\right\}$ be the set of literals. Let $\vec{p}=\left(p_{1}, p_{2}, p_{3}\right)$ be a triple of numbers between 0 and 1 . Then the random formula $F(n, \vec{p})$ is the outcome of the following experiment.

- Choose an assignment $\sigma: V \rightarrow\{0,1\}$ uniformly at random.
- For any triple $\left(l_{1}, l_{2}, l_{3}\right) \in\left\{x_{1}, \bar{x}_{1}, \ldots, x_{n}, \bar{x}_{n}\right\}$ of literals such that $i=\left|\left\{j \in\{1,2,3\}: \sigma\left(l_{j}\right)=1\right\}\right| \geq 1$ include the clause $l_{1} \vee l_{2} \vee l_{3}$ with probaility $p_{i}$ independently.

In words, $F(n, \vec{p})$ has a "planted" assignment $\sigma$, and each possible clause containing $i \geq 1$ satisfied literals under $\sigma$ gets included with probability $p_{i}$ independently. Considering the clauses as ordered triples of literals, we see that the expected total number of clauses is $n^{3} p_{3}+3 n^{2} p_{2}+3 n^{2} p_{1}$. The following result concerning this model is due to Flaxman [14].

Theorem 12 There is a polynomial time algorithm SpecSAT that satisfies the following. Assume that either $p_{2} \geq$ $\delta\left(p_{1}+p_{3}\right)$ or $p_{1} \geq(1+\delta) p_{3}$ or $p_{1} \leq(1-\delta) p_{3}$. Moreover, assume that $n^{2}\left(p_{1}+3 p_{2}+3 p_{3}\right) \geq \zeta$. Then SpecSAT applied to a random formula $\mathcal{F}=F(n, \vec{p})$ finds a satisfying assingment whp.

SpecSAT exploits spectral properties of the "projection graph" $G(\mathcal{F})$ of a random formula $\mathcal{F}=F(n, \vec{p})$. The vertex set of the projection graph is the set of literals, and two literals $l, l^{\prime}$ are adjacent iff they occur together in a clause of $\mathcal{F}$. Hence, each clause corresponds to a triangle in $G(\mathcal{F})$. If the triple $\vec{p}$ satisfies the assumptions of Theorem 12, then the assignment $\sigma$ yields a partition of $G(F)$ with one of the following properties; let $T$ be the set of literals set to true under $\sigma$ and $F=L \backslash T$.

1. The number of $T-F$-edges is at least $\left(\frac{1}{2}+\delta^{\prime}\right)|E(G(\mathcal{F}))|$.
2. The number of edges within the set $T$ is at least $\left(\frac{1}{4}+\delta^{\prime}\right)|E(G(\mathcal{F}))|$.
3. The number of edges within the set $F$ is at least $\left(\frac{1}{4}+\delta^{\prime}\right)|E(G(\mathcal{F}))|$.

Here $\delta^{\prime}>0$ is a number that depends only on $\delta$. In each of the three cases, the partition $T \cup F$ of the vertex set $L$ is reflected in spectral properties of $G(\mathcal{F})$. The algorithm SpecSAT exploits this spectral information to recover (a very good approximation to) the partition $(T, F)$ and hence a satisfying assignment.

However, if $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $(1-\delta) p_{3} \leq p_{1} \leq(1+\delta) p_{3}$, then the partition $(T, F)$ does not stand out anymore. In fact, if $p_{2}=0$ and $p_{1}=p_{3}$, then $G(\mathcal{F})$ is a quasi-random graph (i.e., the global edge distribution is identical to
that of a uniformly random graph with the same number of edges). Hence, in this case it is not possible to recover the partition $(T, F)$ from $G(\mathcal{F})$. Nonetheless, in the case $p_{2}=0$ and $p_{1}=p_{3}$ it is easy to find a satisfying assignment, because then $\mathcal{F}$ is a random 3-XOR formula and thus a satisfying assignment can be found by Gaussian elimination. Of course, this trick only applies if $p_{2}$ is identically zero; if $p_{2}>0$ but $p_{2}<\delta\left(p_{1}+p_{3}\right)$, then the resulting problem is a perturbed 3-XOR formula, in which case Gaussian elimination fails. Our contribution here is an algorithm for solving $F(n, \vec{p})$ that also applies to this case.

Theorem 13 There is a polynomial time algorithm Find3SAT that satisfies the following. Suppose that $n^{2}\left(p_{1}+\right.$ $\left.3 p_{2}+3 p_{3}\right) \geq \sqrt{n} \ln ^{10} n$. Then applied to a random formula $\mathcal{F}=F(n, \vec{p})$ Find3SAT yields a satisfying assignment whp.

Note that Theorem 13 requires the expected number of clauses to be at least $n^{3 / 2} \ln ^{10} n$, whereas Theorem 12 just requires $\zeta n$ clauses for some constant $\zeta>0$. The reason for this is that random "perturbed" 3-XOR formulas seem more difficult to deal with than other types of random formulas. Indeed, perturbed 3-XOR formulas play a distinguished role in the context of refuting the existence of a satisfying assignment for a random 3-SAT formula $F_{3}(n, p)$. Here $p$ is chosen so that $F_{3}(n, p)$ is unsatisfiable whp and the goal is to certify in polynomial time that no satisfying assignment exists (cf. Section 6.2). Given a random formula $\mathcal{F}=F_{3}(n, m)$ with $m \geq \zeta n$ clauses for some large enough constant $\zeta$, it is easy to certify in polynomial time that if $\mathcal{F}$ has a satisfying assignment $\tau$, then $\tau$ satisfies all but $\delta m$ clauses in a 3-XOR fashion (i.e., either all or exactly one literal is satisfied) [9]. But in order to refute the existence of a satisfying assignment of this type (and thus to certify that the formula has no satisfying assignment at all), the best current polynomial time algorithm requires $m \geq \zeta n^{3 / 2}$ [11]. In fact, techniques that allow to improve the bound in Theorem 13 to $n^{3 / 2-\Omega(1)}$ may very well yield improved refutation algorithms (and vice versa).

## Algorithm 14 Find3SAT $(\mathcal{F})$

Input: A 3-SAT formula $\mathcal{F}$ over the variables $V=\left\{x_{1}, \ldots, x_{n}\right\}$ with $m$ clauses.
Output: An assignment $\tau: V \rightarrow\{0,1\}$.

1. If $\operatorname{SpecSAT}(\mathcal{F})$ finds a satisfying assignment of $\mathcal{F}$, output this assignment and terminate.
2. Let $\mathcal{R}$ be the 4-SAT formula obtained from $\mathcal{F}$ as follows.

If $l_{1}, l_{2}, l_{3}, l_{4}$ are four literals such that there is a variable $z$ such that the clauses $l_{1} \vee l_{2} \vee z$ and
$l_{3} \vee l_{4} \vee \bar{z}$ occur in $\mathcal{F}$, then include the clause $l_{1} \vee l_{3} \vee l_{2} \vee l_{4}$ into $\mathcal{R}$.
Call $\operatorname{ApxCSP}(\mathcal{R}, 100, \delta)$ and let $\tau^{\prime}$ be the resulting assignment. Let $\tau^{\prime \prime}$ be the inverse of $\tau^{\prime}$. Let $\tau$ be the assignment among $\tau^{\prime}, \tau^{\prime \prime}$ that satisfies the larger number of clauses of $\mathcal{F}$.
3. Repeat the following $\lceil\ln n\rceil$ times.
4. For any literal $\lambda$ set to false under $\tau$ compute the number $\nu_{\lambda}$ of clauses $\lambda \vee l \vee l^{\prime}$ such that $l, l^{\prime}$ are
literals set to true under $\tau$. Let $\Lambda=\left\{\lambda: 8 n \nu_{\lambda}>m\right\}$. Modify $\tau$ so that all literals in $\Lambda$ get set to true.
5. Output $\tau$ if it is a satisfying assignment. Otherwise, output "fail".

In its first step Find3SAT calls SpecSAT, which finds a satisfying assignment whp unless $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $(1-\delta) p_{3} \leq p_{1} \leq(1+\delta) p_{3}$. If SpecSAT fails, Find3SAT sets up the 4 -SAT formula $\mathcal{R}$ via the resolution principle. Hence, $\mathcal{R}$ is a satisfiable $4-$ SAT formula. Then, Find3SAT applies ApxCSP to $\mathcal{R}$. The following lemma shows that calling ApxCSP is feasible.

Lemma 8 Suppose that $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $(1-\delta) p_{3} \leq p_{1} \leq(1+\delta) p_{3}$. For any fixed number $\gamma>0$ the formula $\mathcal{R}$ is $(100, \gamma)$-bounded whp

Hence, $\operatorname{ApxCSP}$ outputs an assignment $\tau$ that satisfies at least a $1-\delta$ fraction of all clauses of $\mathcal{R}$.
Lemma 9 Suppose that $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $(1-\delta) p_{3} \leq p_{1} \leq(1+\delta) p_{3}$. The assignment $\tau$ computed in Step 2 is within Hamming distance at most $0.01 n$ of the planted assignment $\sigma$ whp

Hence, Step 2 yields an assignment that is "close" to the planted assignment whp Then, Steps 3 and 4 perform a local improvement operation.

Lemma 10 Suppose that $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $(1-\delta) p_{3} \leq p_{1} \leq(1+\delta) p_{3}$. After $i$ iterations of Step 4 the assignment $\tau$ is at Hamming distance at most $5^{-i} n$ from the planted assignment $\sigma$.

Finally, Theorem 13 is an immediate consequence of Theorem 12 and Lemmas 8-10.

### 6.3.1 Proof of Lemma 8.

Let $\mathcal{F}$ be a 3-SAT formula over $V$. We set up a 4 -tensor $\mathbf{A}(\mathcal{F}): L^{4} \rightarrow\{0,1\}$ with entries

$$
\mathbf{A}\left(l_{1}, l_{2}, l_{3}, l_{4}\right)= \begin{cases}1 & \text { if the clause } l_{1} \vee l_{2} \vee l_{3} \vee l_{4} \text { occurs in } \mathcal{R} \\ 0 & \text { otherwise } .\end{cases}
$$

In order to show that $\mathcal{F}$ is $(100, \gamma)$-bounded whp, it suffices to show that
$\mathbf{A}(\mathcal{F}(n, \vec{p}))$ is $(100, \gamma)$-bounded whp for any fixed $\gamma>0$.
Let $q=p_{1}+p_{2}+p_{3}$, set $\vec{q}=(q, q, q)$, and let $\mathcal{F}^{*}=F(n, \vec{q})$. We can think of $\mathcal{F}^{*}$ as a random formula obtained by first choosing $\mathcal{F}=F(n, \vec{p})$, and then adding each possible clause with $i$ satisfied literals that is not present in $\mathcal{F}$ with probability $\left(q-p_{i}\right) /\left(1-p_{i}\right)$ independently. Since the expected number of clauses in $\mathcal{F}(n, \vec{q})$ is at most three times the expected number of clauses in $F(n, \vec{p})$, the following implies (33).
$\mathbf{A}(F(n, \vec{q}))$ is $(33, \gamma)$-bounded whp for any fixed $\gamma>0$.
To show (34), we employ the following result from [8, Lemma 3.3]. We let $\vec{J}$ signify a matrix with all entries equal to one.

Lemma 11 Let $\mathcal{F}=F(n, \vec{q})$, Let $\mathbf{B}(\mathcal{F})$ be the $(2 n)^{2} \times(2 n)^{2}$ matrix constructed from $\mathbf{A}(\mathcal{F})$ as in (3). Let $Q=n q^{2}$. Then $\|Q \vec{J}-\mathbf{B}(\mathcal{F})\|=o\left(n^{2} Q\right)$.
Proof of (34). Let $\mathbf{B}=\mathbf{B}(F(n, \vec{q}))$. The expected number $\mathrm{E}\|\mathbf{B}\|_{\square}$ of ones in $\mathbf{B}$ is $(2 n)^{4} Q$. For $\mathbf{B}\left(\left(l_{1}, l_{2}\right),\left(l_{3}, l_{4}\right)\right)$ equals one iff there is a variable $z$ such that both $l_{1} \vee l_{3} \vee z$ and $l_{2} \vee l_{4} \vee \bar{z}$ occur in $F(n, \vec{q})$, and since the probability of this event is $q^{2}$, and there are $n$ ways to choose $z$, we have $\mathrm{P}\left[\mathbf{B}\left(\left(l_{1}, l_{2}\right),\left(l_{3}, l_{4}\right)\right)=1\right]=n q^{2}=Q$. Hence, Chernoff bounds entail that $\|\mathbf{B}\|_{\square} \sim(2 n)^{4} Q$ whp. Consequently, the density $\hat{Q}=(2 n)^{-4}\|\mathbf{B}\|_{\square}$ satisfies $\hat{Q} \sim Q$ whp.

Let $S, T \subset L \times L$ be sets of size at least $\gamma n^{2}$. Let $\overrightarrow{1}_{S} \in\{0,1\}^{L \times L}$ be the indicator of $S$, and let $\overrightarrow{1}_{T}$ be the indicator of $T$. Then by Lemma 11

$$
\begin{aligned}
&|Q| S \times T|-\mathbf{B}(S, T)|=\left|\left\langle(Q \vec{J}-\mathbf{B}) \overrightarrow{1}_{T}, \overrightarrow{1}_{S}\right\rangle\right| \leq\|Q \vec{J}-\mathbf{B}\| \cdot\left\|\overrightarrow{1}_{T}\right\| \cdot\left\|\overrightarrow{1}_{S}\right\| \\
&=o\left(n^{2} Q\right) \cdot \sqrt{|S \times T|}=o(Q|S \times T|)
\end{aligned}
$$

where the last step follows from the assumption $|S|,|T| \geq \gamma n^{2}$. Hence, $\mathbf{B}(S, T) \sim Q|S \times T| \mathbf{w h p}$. As $Q \sim \hat{Q}$, this implies that $\mathbf{B}(S, T) \sim \hat{Q}|S \times T|$, and thus $\mathbf{B}(S, T) \leq 1.01 \hat{Q}|S \times T|$ whp. Consequently, $\mathbf{B}$ is $(1.01, \gamma)$-bounded whp, whence (34) follows (with room to spare).

### 6.3.2 Proof of Lemma 9.

Let $\alpha=0.01$.
Lemma 12 W.h.p. either $\tau^{\prime}$ or $\tau^{\prime \prime}$ is within Hamming distance $\leq \alpha$ n of $\sigma$.
Proof Without loss of generality we may assume that $\sigma(x)=1$ for all $x \in V$. We will use a first moment argument to show that whp any assignment $\tau^{\prime}$ that satisfies a $(1-\delta)$-fraction of the clauses of $\mathcal{R}$ has the property that either $\tau^{\prime}$ or its inverse $\tau^{\prime \prime}$ is at Hamming distance at most $\alpha n$ from $\sigma$. Then the assertion follows from Theorem 3.

Thus, consider any assignment $\tau^{\prime}$ such that both $\tau^{\prime}$ and $\tau^{\prime \prime}$ are at Hamming distance more than $\alpha n$ from $\sigma$. Then there are at least $\alpha n$ literals $l$ such that $\sigma(l)=1$ and $\tau^{\prime}(l)=0$ and at least $\alpha n$ literals $l^{\prime}$ such that $\sigma\left(l^{\prime}\right)=0$ and
$\tau^{\prime}\left(l^{\prime}\right)=1$. For a variable $z$ we let $\ell_{z}$ be the set of all clauses $l_{1} \vee l_{2} \vee z$ in $\mathcal{F}=F(n, \vec{p})$ such that $\sigma\left(l_{1}\right)=\sigma\left(l_{2}\right)=1$ and $\tau^{\prime}\left(l_{1}\right)=\tau^{\prime}\left(l_{2}\right)=0\left(l_{1}, l_{2} \in L\right)$. Moreover, $\bar{\ell}_{z}$ denotes the set of all clauses $l_{3} \vee l_{4} \vee \bar{z}$ in $\mathcal{F}$ such that $\sigma\left(l_{3}\right)=1$ and $\sigma\left(l_{4}\right)=\tau^{\prime}\left(l_{3}\right)=\tau^{\prime}\left(l_{4}\right)=0\left(l_{3}, l_{4} \in L\right)$. Then $\left(\left|\ell_{z}\right|,\left|\bar{\ell}_{z}\right|\right)_{z \in Z}$ is a family of mutually independent binomial random variables. Moreover, for all $z \in V$

$$
\begin{equation*}
\mathrm{E}\left(\left|\ell_{z}\right|\right) \geq \alpha^{2} n^{2} p_{3} \geq \sqrt{n} \ln ^{9} n, \quad \mathrm{E}\left(\left|\bar{\ell}_{z}\right|\right) \geq \alpha^{2} n^{2} p_{1} \geq \sqrt{n} \ln ^{9} n \tag{35}
\end{equation*}
$$

Let us call $z$ bad for $\tau^{\prime}$ if either $\ell_{z}<\mathrm{E}\left(\left|\ell_{z}\right|\right) / 2$ or $\bar{\ell}_{z}<\mathrm{E}\left(\left|\bar{\ell}_{z}\right|\right) / 2$. Then (35) implies in combination with Chernoff bounds that $z$ is bad with probability at most $\exp (-\sqrt{n})$. Since the numbers $\left|\ell_{z}\right|,\left|\bar{\ell}_{z}\right|$ are independent for all $z$, this entails that with probability at least $1-4^{-n}$ there are at most $n / 2$ bad variables. Hence, by the union bound there is no assignment $\tau^{\prime}$ such that both $\tau^{\prime}$ and $\tau^{\prime \prime}$ are at Hamming distance more than $\alpha n$ from $\sigma$ and $\tau^{\prime}$ has more than $n / 2$ bad variables.

Thus, suppose that there are less than $n / 2$ bad variables for $\tau^{\prime}$. Consider a $z \in V$ that is not bad. Then every pair of clauses $\left(l_{1}, l_{2}, z\right) \in \ell_{z},\left(l_{3}, l_{4}, z\right) \in \bar{\ell}_{z}$ yields a clause $\left(l_{1}, l_{3}, l_{2}, l_{4}\right)$ of $\mathcal{R}$ that $\tau^{\prime}$ does not satisfy. Consequently, $\tau^{\prime}$ fails to satisfy at least

$$
\begin{equation*}
\sum_{z \in V} \ell_{z} \bar{\ell}_{z} \geq \frac{n}{2} \cdot \frac{\alpha^{4}}{4} n^{4} p_{1} p_{3}=\alpha^{4} n^{5} p_{1} p_{3} / 8 \tag{36}
\end{equation*}
$$

clauses. On the other hand, the expected number of clauses in $\mathcal{R}$ is at most $2 n^{5} p_{1} p_{3}$ by our assumption that $p_{2}<$ $\delta\left(p_{1}+p_{3}\right)$ and $(1-\delta) p_{3} \leq p_{1} \leq(1+\delta) p_{3}$. Hence, by Chernoff bounds the total number of clauses in $\mathcal{R}$ is at most $3 n^{5} p_{1} p_{3}$ whp. Therefore, the assertion follows from the lower bound (36) on the number of clauses that are unsatisfied under $\tau^{\prime}$, provided that $\delta$ is sufficiently small.

To complete the proof of Lemma 9, we establish the following two facts.
Fact 15 W.h.p. any assignment $\chi: V \rightarrow\{0,1\}$ at Hamming distance at least $(1-\alpha) n$ from the planted assignment $\sigma$ fails to satisfy at least $(1-\alpha)^{4} n^{3} p_{3}$ clauses of $\mathcal{F}=F(n, \vec{p})$.

Proof We use a first moment argument. If $\chi$ is at Hamming distance at least $(1-\alpha) n$ from $\sigma$, then there are at least $(1-\alpha) n$ literals $l$ such that $\sigma(l)=1-\chi(l)=1$. Hence, the expected number of clauses $l_{1} \vee l_{2} \vee l_{3}$ occurring in $F(n, \vec{p})$ such that $\sigma\left(l_{i}\right)=1-\chi\left(l_{i}\right)=1$ for $i=1,2,3$ is at least $(1-\alpha)^{3} n^{3} p_{3} \geq(1-\alpha)^{3} n^{3 / 2} \ln ^{9} n$. Clearly, each of these clauses is unsatisfied under $\chi$. Moreover, the number of such clauses is binomially distributed, whence Chernoff bounds entail that with probability at least $1-4^{-n} \chi$ fails to satisfy at least $(1-\alpha)^{4} n^{3 / 2} \ln ^{9} n$ clauses. As there are only $2^{n}$ assignments $V \rightarrow\{0,1\}$ in total, the assertion follows from the union bound.

Fact 16 W.h.p. any assignment $\chi: V \rightarrow\{0,1\}$ at Hamming distance at most $\alpha n$ from the planted assignment $\sigma$ fails to satisfy at most $6 \alpha\left(n^{3} p_{3}+3 n^{3} p_{1}+3 n^{3} p_{2}\right)$ clauses of $\mathcal{F}=F(n, \vec{p})$.

Proof Any clause of $\mathcal{F}=F(n, \vec{p})$ that is unsatisfied under $\chi$ contains a literal $l$ such that $\sigma(l) \neq \chi(l)$. The probability that a randomly chosen literal has this property is at most $\alpha$. Since each clause contains three literals, the expected number of clauses that contain a literal $l$ such that $\sigma(l) \neq \chi(l)$ is at most $3 \alpha\left(n^{3} p_{3}+3 n^{3} p_{1}+3 n^{3} p_{2}\right)$. As the number of such clauses is binomially distributed, with probability at least $1-4^{-n}$ there are at most $6 \alpha\left(n^{3} p_{3}+\right.$ $3 n^{3} p_{1}+3 n^{3} p_{2}$ ) of them. Hence, the assertion follows from the union bound.

Finally, Lemma 9 is an immediate consequence of Lemma 12 and Facts 15 and 16.

### 6.3.3 Proof of Lemma 10.

We shall establish the following fact via a first moment argument. Let $\alpha=0.01$.
Fact 17 W.h.p. any assignment $\tau: V \rightarrow\{0,1\}$ that is at Hamming distance $\Delta \leq \alpha n$ from $\sigma$ satisfies the following.

1. The number of literals $\lambda$ such that $\tau(\lambda)=1-\sigma(\lambda)=0$ and $\nu_{\lambda} \leq m /(8 n)$ is less than $\Delta / 10$.
2. The number of literals $\lambda$ such that $\tau(\lambda)=\sigma(\lambda)=0$ and $\nu_{\lambda}>m /(8 n)$ is less than $\Delta / 10$.

Since the initial assignment $\tau$ has Hamming distance at most $\alpha n$ from $\sigma$ by Lemma 9, Fact 17 shows that the Hamming distance of $\tau$ and $\sigma$ decreases by a factor of 5 in each iteration. Hence, after at most $\ln n$ iterations we have $\sigma=\tau$, as desired.

Thus, the remaining task is to establish Fact 17. Let $0<\Delta \leq \alpha n$ and let $\tau: V \rightarrow\{0,1\}$ be any assignment such that $\tau$ and $\sigma$ have Hamming distance $\Delta$. With respect to the first item, suppose that $\tau(\lambda)=0$ and $\sigma(\lambda)=1$. There are at least $(1-\alpha)^{2} n^{2}$ pairs $\left(l, l^{\prime}\right)$ of literals such that $\sigma(l)=\sigma\left(l^{\prime}\right)=\tau(l)=\tau\left(l^{\prime}\right)=1$, and for each such pair the clause $\lambda \vee l \vee l^{\prime}$ is present in the random formula $\mathcal{F}=F(n, \vec{p})$ with probability $p_{3}$ independently. Hence, letting $\nu_{\lambda}^{\prime}$ signify the number of such clauses, we have $\nu_{\lambda} \geq \nu_{\lambda}^{\prime}$ and $\nu_{\lambda}^{\prime}$ is binomially distributed with mean at least $(1-\alpha)^{2} n^{2} p_{3} \geq \sqrt{n} \ln ^{9} n$. Hence, by Chernoff bounds

$$
\mathrm{P}\left[\nu_{\lambda}<0.99 \cdot(1-\alpha)^{2} n^{2} p_{3}\right] \leq \mathrm{P}\left[\nu_{z}^{\prime}<0.99 \cdot(1-\alpha)^{2} n^{2} p_{3}\right] \leq \exp (-\sqrt{n})
$$

Since the random variables $\nu_{\lambda}$ are mutually independent for all $\lambda$, the number of literals $\lambda \in \tau^{-1}(0) \cap \sigma^{-1}(1)$ such that $\nu_{\lambda}<0.99 \cdot(1-\alpha)^{2} n^{2} p_{3}$ is dominated by a binomially distributed random variable with mean $\exp (-\sqrt{n}) \Delta$. Hence, with probability at least $1-n^{-3 \Delta}$ the number of $\lambda$ with $\tau(\lambda)=1-\sigma(\lambda)=0$ and $\nu_{\lambda}<0.99 \cdot(1-\alpha)^{2} n^{2} p_{3}$ is less than $\Delta / 10$. Since there are $\binom{n}{\Delta}$ assignments $\tau$ at Hamming distance $\Delta$ from $\sigma$, we thus conclude that with probability at least $1-\binom{n}{\Delta} n^{-3 \Delta} \geq 1-n^{-2}$ all of them have the property that there are at most $\Delta / 10$ such $\lambda$. Finally, as $m$ is binomially distributed with mean $n^{3} p_{3}+3 n^{3} p_{2}+3 n^{3} p_{1}$, Chernoff bounds yield that $m \sim n^{3} p_{3}+3 n^{3} p_{2}+3 n^{3} p_{1}$ whp. If this is so, then any literal $\lambda$ with $\tau(\lambda))=1-\sigma(\lambda)=0$ that satisfies $\nu_{\lambda} \geq 0.99 \cdot(1-\alpha)^{2} n^{2} p_{3}$ also satisfies $\nu_{\lambda} \geq m /(8 n)$, because we are assuming that $p_{2} \leq \delta\left(p_{1}+p_{2}\right)$ and $p_{1} \leq(1+\delta) p_{3}$. This complete the proof of the first item.

Regarding the second item, we consider a literal $\lambda$ such that $\tau(\lambda)=\sigma(\lambda)=0$. If $l, l^{\prime} \in \tau^{-1}(1)$ are literals such that $\lambda \vee l \vee l^{\prime}$ occurs as a clause in $\mathcal{F}=F(n, \vec{p})$, then either $\sigma(l)=\sigma\left(l^{\prime}\right)=1$ and $\lambda \vee l \vee l^{\prime}$ has exactly two satisfied literals under $\sigma$, or $\sigma$ and $\tau$ differ on exactly one of $l, l^{\prime}$. Hence, the expected number of such clauses is at most

$$
\mathrm{E}\left(\nu_{\lambda}\right) \leq n^{2} p_{2}+2 \alpha n^{2} p_{1} \leq 3 \alpha n^{2} p_{1},
$$

because we are assuming that $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $p_{1} \geq(1-\delta) p_{3}$. Since $3 \alpha n^{2} p_{1} \geq \sqrt{n} \ln ^{9} n$ and $\nu_{\lambda}$ is binomially distributed, Chernoff bounds imply that

$$
\mathrm{P}\left[\nu_{\lambda}>6 \alpha n^{2} p_{1}\right] \leq \exp (-\sqrt{n})
$$

Furthermore, the numbers $\nu_{\lambda}$ are mutually independent. Hence, the number of all $\lambda$ with $\tau(\lambda)=\sigma(\lambda)=0$ such that $\nu_{\lambda}>6 \alpha n^{2} p_{1}$ is binomially distributed with mean $\leq n \exp (-\sqrt{n})$. Therefore, Chernoff bounds entail that with probability at least $1-n^{-3 \Delta}$ there are at most $\Delta / 10$ such $\lambda$. Since the total number of assignments $\tau$ at Hamming distance $\Delta$ from $\sigma$ is $\binom{n}{\Delta}$, we conclude that with probability at least $1-n^{-2}$ for all of them there are at most $\Delta / 10$ literals $\lambda$ with $\tau(\lambda)=\sigma(\lambda)=0$ such that $\nu_{\lambda}>6 \alpha n^{2} p_{1}$. Finally, as $m \sim n^{3} p_{3}+3 n^{3} p_{2}+3 n^{3} p_{1}$ whp and because we are assuming that $p_{2} \leq \delta\left(p_{1}+p_{3}\right)$ and $(1+\delta) p_{3} \geq p_{1} \geq(1-\delta) p_{3}$, we have $m /(8 n)>6 \alpha n^{2} p_{1}$ whp. Hence, there are at most $\Delta / 10$ literals $\lambda$ with $\tau(\lambda)=\sigma(\lambda)=0$ and $\nu_{\lambda}>m /(8 n)$ whp.

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