# Tight Bounds on the Approximability of Almost-satisfiable Horn SAT and Exact Hitting Set 

Venkatesan Guruswami* Yuan Zhou*

Computer Science Department<br>Carnegie Mellon University<br>Pittsburgh, PA 15213.


#### Abstract

We study the approximability of two natural Boolean constraint satisfaction problems: Horn satisfiability and exact hitting set. Under the Unique Games conjecture, we prove the following optimal inapproximability and approximability results for finding an assignment satisfying as many constraints as possible given a near-satisfiable instance. 1. Given an instance of Max Horn-3SAT that admits an assignment satisfying $(1-\varepsilon)$ of its constraints for some small constant $\varepsilon>0$, it is hard to find an assignment satisfying more than $(1-1 / O(\log (1 / \varepsilon)))$ of the constraints. This matches a linear programming based algorithm due to Zwick [Zwi98], resolving the natural open question raised in that work concerning the optimality of the approximation bound. Given a $(1-\varepsilon)$ satisfiable instance of Max Horn-2SAT for some constant $\varepsilon>0$, it is possible to find a $(1-2 \varepsilon)$-satisfying assignment efficiently. This improves the algorithm given in [KSTW00] which finds a $(1-3 \varepsilon)$-satisfying assignment, and also matches the ( $1-c \varepsilon$ ) hardness for any $c<2$ derived from vertex cover (under UGC). 2. An instance of $\operatorname{Max} 1-\mathrm{in}-k$ - HS consists of a universe $U$ and a collection $\mathcal{C}$ of subsets of $U$ of size at most $k$, and the goal is to find a subset of $U$ that intersects the maximum number of sets in $\mathcal{C}$ at a unique element. We prove that Max 1 -in- $k$-HS is hard to approximate within a factor of $O(1 / \log k)$ for every fixed integer $k$. This matches (up to constant factors) an easy factor $\Omega(1 / \log k)$ approximation algorithm for the problem, and resolves a question posed in [GT05]. It is crucial for the above hardness that sets of size up to $k$ are allowed; indeed, when all sets have size $k$, there is a simple factor $1 / e$-approximation algorithm. Moreover, the hardness applies when the instance admits a solution satisfying $1-1 / k^{0.99}$ of the constraints. In contrast, for instances which are $\left(1-\frac{1}{1.01 k}\right)$-satisfiable, we give a constant factor approximation algorithm. Our hardness results are proved by constructing integrality gap instances for a semidefinite programming relaxation for the problems, and using Raghavendra's result [Rag08] to conclude that no algorithm can do better than the SDP assuming the UGC. The algorithmic results are based on rounding appropriate linear programming relaxations.


[^0]
## 1 Introduction

Schaefer proved long ago that there are only three non-trivial classes of Boolean constraint satisfaction problems (CSPs) for which satisfiability is polynomial time decidable [Sch78]. These are LIN-mod-2 (linear equations modulo 2), 2-SAT, and Horn-SAT. An instance of Horn-SAT is a CNF formula where each clause consists of at most one unnegated literal. ${ }^{1}$ The maximization versions of these problems (where the goal is to find an assignment satisfying the maximum number of constraints) are NP-hard, and in fact APX-hard, i.e., NP-hard to approximate within some constant factor bounded away from 1. An interesting special case of the maximization version is the following problem of "finding almost-satisfying assignments": Given an instance which is $(1-\varepsilon)$-satisfiable (i.e., only $\varepsilon$ fraction of constraints need to be removed to make it satisfiable for some small constant $\varepsilon)$, can one efficiently find an assignment satisfying most (say, $1-f(\varepsilon)-o(1)$ where $f(\varepsilon) \rightarrow 0$ as $\varepsilon \rightarrow 0)$ of the constraints? ${ }^{2}$

The problem of finding almost-satisfying assignments was first suggested and studied in a beautiful paper by Zwick [Zwi98]. This problem seems well-motivated, as even if a Max CSP is APX-hard in general, in certain practical situations instances might be close to being satisfiable (for example, a small fraction of constraints might have been corrupted by noise). An algorithm that is able to satisfy most of the constraints of such an instance could be very useful.

As pointed out in [KSTW00], Schaefer's reductions together with the PCP theorem imply that the previous goal is NP-hard to achieve for any Boolean CSP for which the satisfiability problem is NP-complete. Indeed, all but the above three tractable cases of Boolean CSPs have a "gap at location 1," which means that given a satisfiable instance it is NP-hard to find an assignment satisfying $\alpha$ fraction of the constraints for some constant $\alpha<1$. This result has been extended to CSPs over arbitrary domains recently [JKK09].

The natural question therefore is whether for the three tractable Boolean CSPs, LIN-mod-2, 2SAT, and Horn-SAT, one can find almost-satisfying assignments in polynomial time. Effectively, the question is whether there are "robust" satisfiability checking algorithms that can handle a small number of inconsistent constraints and still produce a near-satisfying assignment. As remarked above, the Unique Games conjecture holds if and only if there is no such robust satisfiability algorithm for the Unique Games CSP.

With respect to the feasibility of finding almost-satisfying assignments, LIN-mod-2, 2-SAT, and Horn-SAT behave rather differently from each other. For LIN-mod-2, Håstad in his breakthrough paper [Hås01] showed that for any $\varepsilon, \delta>0$, finding a solution satisfying $1 / 2+\delta$ of the equations of a $(1-\varepsilon)$-satisfiable instance is NP-hard. In fact, this result holds even when each equation depends on only 3 variables. Since just picking a random assignment satisfies $1 / 2$ the constraints in expectation, this shows, in a very strong sense, that there is no robust satisfiability algorithm for LIN-mod-2.

In sharp contrast to this extreme hardness for linear equations, Zwick [Zwi98] proved that for 2SAT and Horn-SAT one can find almost-satisfying assignments in polynomial time. For Max 2SAT,

[^1]Zwick gave a semidefinite programming (SDP) based algorithm that finds a $\left(1-O\left(\varepsilon^{1 / 3}\right)\right.$ )-satisfying assignment (i.e., an assignment satisfying a fraction $\left(1-O\left(\varepsilon^{1 / 3}\right)\right)$ of the constraints) given as input a $(1-\varepsilon)$-satisfiable instance. This algorithm was later improved to one that finds a $1-O(\sqrt{\varepsilon})$ satisfying assignment by Charikar, Makarychev, and Makarychev [CMM09]. The $1-O(\sqrt{\varepsilon})$ bound is known to be best possible under the Unique Games conjecture (UGC) [Kho02, KKMO07]. In fact, this hardness result for Max 2SAT was the first application of the UGC and one of the main initial motivations for its formulation by Khot [Kho02].

For Max Horn-SAT, Zwick gave a linear programming (LP) based algorithm to find an assignment satisfying $(1-O(\log \log (1 / \varepsilon) / \log (1 / \varepsilon)))$ of constraints of a $(1-\varepsilon)$-satisfiable instance. For Max Horn-3SAT where each clause has at most three literals, the algorithm finds a $(1-O(1 / \log (1 / \varepsilon)))$ satisfying assignment. Note that the fraction of unsatisfied constraints is exponentially worse for Max Horn-SAT compared to Max 2SAT.

Horn-SAT is a fundamental problem in logic and artificial intelligence. Zwick's robust Horn satisfiability algorithm shows the feasibility of solving instances where a small number of constraints are faulty and raises the following natural question, which was also explicitly raised in [Zwi98]. Is this $1 / \log (1 / \varepsilon)$ deficit inherent? Or could a more sophisticated algorithm, say based on an SDP relaxation instead of the LP relaxation used in [Zwi98], improve the deficit to something smaller (such as $\varepsilon^{b}$ for some constant $b$ as in the case of the SDP based algorithm for Max 2SAT)? It is known that for some absolute constant $c<1$, it is NP-hard to find a $\left(1-\varepsilon^{c}\right)$-satisfying assignment given a $(1-\varepsilon)$-satisfiable instance of Max Horn-SAT [KSTW00].

In this work, we address the above question and resolve it (conditioned on the UGC), showing the $1 / \log (1 / \varepsilon)$ deficit to be inherent. We also investigate another problem, the "exact hitting set" problem for set sizes bounded by $k$, which has a very peculiar approximation behavior [GT05]. It admits a much better approximation algorithm on satisfiable instances, as well as when sets all have size exactly (or close to) $k$. We prove that these restrictions are inherent, and relaxing these rules out a constant factor approximation algorithm (again, under the UGC). We describe our results in more detail below in Section 2.

Remark 1. For $(1-\varepsilon)$-satisfiable instances of Max 2-SAT, even the hardness of finding a $\left(1-\omega_{\varepsilon}(1) \varepsilon\right)$ satisfying assignment is not known without assuming the UGC (and the UGC implies the optimal $1-\Omega(\sqrt{\varepsilon})$ hardness bound). For Max Horn-SAT, as mentioned above, we know the NP-hardness of finding a $\left(1-\varepsilon^{c}\right)$-satisfying assignment for some absolute constant $c<1$. Under the UGC, we are able to pin down the exact asymptotic dependence on $\varepsilon$.

## 2 Our results and previous work

### 2.1 Horn-SAT

We prove the following hardness result concerning finding almost-satisfying assignments for Max Horn-SAT (in fact for the arity 3 case). In the sequel, we use the terminology "UG-hard" to mean at least as hard as refuting the Unique Games conjecture.

Theorem 1. For some absolute constant $C>0$, for every $\varepsilon>0$, given a $(1-\varepsilon)$-satisfiable instance of Max Horn-3SAT, it is UG-hard to find an assignment satisfying more than a fraction $\left(1-\frac{C}{\log (1 / \varepsilon)}\right)$
of the constraints.
Zwick gave a polynomial time algorithm that finds a $1-O\left(\frac{\log k}{\log (1 / \varepsilon)}\right)$-satisfying assignment on input a $(1-\varepsilon)$-satisfiable instance of Max Horn- $k$ SAT. Our inapproximability bound is therefore optimal up to the constant $C$, and resolves Zwick's question on whether his algorithm can be improved in the negative. (For arbitrary arity Horn-SAT, Zwick's algorithm has the slightly worse $1-O(\log \log (1 / \varepsilon) / \log (1 / \varepsilon))$ performance ratio; we do not show this to be tight.)

Theorem 1 shows that Max Horn-SAT has a very different quantitative behavior compared to Max 2SAT with respect to approximating near-satisfiable instances: the fraction of unsatisfied clauses $\Omega(1 / \log (1 / \varepsilon))$ is exponentially worse than the $O(\sqrt{\varepsilon})$ fraction that can be achieved for Max 2SAT.

A strong hardness result for Min Horn Deletion, the minimization version for Horn-SAT, was shown in [KSTW00]. It follows from their reduction that for some absolute constant $c<1$, it is NP-hard to find a $\left(1-\varepsilon^{c}\right)$-satisfying assignment given a $(1-\varepsilon)$-satisfiable instance of Max Horn-SAT. The constant $c$ would be extremely close to 1 in this result as it is related to the soundness in Raz's parallel repetition theorem. While our inapproximability bound is stronger and optimal, we are only able to show UG-hardness and not NP-hardness.

In light of our strong hardness result for Max Horn-3SAT, we also consider the approximability of the arity two case. For Max Horn-2SAT, given a $(1-\varepsilon)$-satisfiable instance, an approximation preserving reduction from vertex cover shows that it is UG-hard to find a $(1-c \varepsilon)$-satisfying assignment for $c<2$. It is also shown in [KSTW00] that one can find a $(1-3 \varepsilon)$-satisfying assignment efficiently. We improve the algorithmic bound (to the matching UG-hardness) by proving the following theorem, based on half-integrality of an LP relaxation for the problem.

Theorem 2. Given a $(1-\varepsilon)$-satisfiable instance for Max Horn-2SAT, it is possible to find a $(1-2 \varepsilon)$ satisfying assignment in polynomial time.

### 2.2 Exact hitting set

We consider the "exact hitting set" problem where the goal is to find a subset that intersects a maximum number of sets from an input family at exactly one element. Formally,

Definition 1. Let $k \geqslant 2$ be a fixed integer. An instance of Max 1-in- $k$ - HS consists of a universe $U=\left\{x_{1}, x_{2}, \ldots, x_{n}\right\}$ and a collection $\mathcal{C}$ of subsets of $U$ each of size at most $k$. The objective is to find a subset $S \subseteq U$ that maximizes the number of sets $T \in \mathcal{C}$ for which $|S \cap T|=1$. When all sets in $\mathcal{C}$ have size equal to $k$, we refer to the problem as $\operatorname{Max} 1-\mathrm{in}-\mathrm{E} k-\mathrm{HS}$.

In addition to being a natural CSP, the exact hitting set problem arises in many contexts where one has to make unique choices from certain specified subsets. The complexity of this problem was investigated in [GT05] and [DFHS08], where applications of the problem to pricing, computing ad-hoc selective families for radio broadcasting, etc. are also discussed.

Our interest in this problem stems in part from the following peculiar approximability behavior of Max 1 -in- $k$-HS, as pointed out in [GT05]. The Max $1-i n-k$-HS problem appears to be much easier to approximate on "satisfiable instances" (where a hitting set intersecting all subsets exactly once exists) or when all sets have size exactly equal to $k$ (instead of at most $k$ ). In both these cases,
there is a factor $1 / e$-approximation algorithm, and obtaining a $(1 / e+\varepsilon)$-approximation is NP-hard even when both restrictions hold simultaneously [GT05].

For Max 1-in- $k$-HS itself, the best approximation factor known to be possible in polynomial time is $\Omega(1 / \log k)$. This is based on partitioning the collection $\mathcal{C}$ into $O(\log k)$ parts based on geometric intervals $\left[2^{i}, 2^{i+1}\right.$ ) of set sizes, and running a simple randomized algorithm (that handles the case where all sets have sizes within a factor of two) on the sub-collection which has the most sets. Despite the simplicity and seeming naiveness of this algorithm, no factor $\omega(1 / \log k)$ algorithm is known for the problem. No hardness factor better than the $1 / e$ bound (which holds even for Max 1-in-Ek-HS) is known either. Improving the gap in our understanding of the approximability of Max 1 -in- $k$-HS was posed as an open question in [GT05].

For the case when $k$ is not fixed but can also grow with the universe size $n$, a factor $(\log n)^{-\Omega(1)}$ hardness was shown in [DFHS08], under the assumption NP $\nsubseteq \operatorname{TIME}\left(2^{n^{\gamma}}\right)$ for some $\gamma>0$. However, their method does not seem to be applicable to the case of bounded set size.

In this work, we prove the following tight result, establishing the difficulty of improving the simple $\Omega(1 / \log k)$-approximation algorithm. This shows that it is hard to simultaneously do well on two different "scales" of set sizes.

Theorem 3. For some absolute constant $C^{\prime}>0$, for every $\alpha>0$, given a ( $1-1 / k^{1-\alpha}$ )-satisfiable instance of $\operatorname{Max} 1-\mathrm{in}-k$ - HS , it is $U G$-hard to find a subset intersecting more than a fraction $\frac{C^{\prime}}{\alpha \log k}$ of the sets exactly once.

The gap in the above hardness result is also located at the "correct" satisfiability threshold, as we show the following complementary algorithmic result. Our algorithm in fact works for the more general Max 1-in- $k$-SAT problem where negated literals are allowed and the goal is to find an assignment for which a maximum number of clauses have exactly one literal set to true. For satisfiable instances of Max 1-in- $k$-SAT, a factor $1 / e$ approximation algorithm was given in [GT05].

Theorem 4. For every constant $B>1$, the following holds. There is a polynomial time algorithm that, given a $\left(1-\frac{1}{B k}\right)$-satisfiable instance of $\operatorname{Max} 1-\mathrm{in}-k$-SAT, finds a truth-assignment on variables satisfying exactly one literal in a fraction $\lambda$ of the clauses, where $\lambda=\left(\frac{1-1 / \sqrt{B}}{e}\right)^{2}$.

## 3 Proof method

We construct integrality gap instances for the semidefinite programming relaxation given in [Rag08], and then use Raghavendra's theorem [Rag08] to conclude that assuming the Unique Games conjecture, no algorithm can achieve an approximation ratio better than the SDP integrality gap.

For the algorithm for Max $1-\mathrm{in}-k$-HS (in fact Max $1-\mathrm{in}-k$-SAT), we use the natural LP relaxation, and show that randomized rounding gives a good approximation.

### 3.1 The canonical SDP for Boolean CSPs and UG-Hardness

For Boolean CSP instances, we write $\mathcal{C}$ as the set of constraints over variables $x_{1}, x_{2}, \cdots, x_{n} \in\{0,1\}$. The SDP relaxation from [Rag08], which we call the canonical SDP, sets up for each constraint
$C \in \mathcal{C}$ a local distribution $\pi_{C}$ on all the truth-assignments $\left\{\sigma: X_{C} \rightarrow\{0,1\}\right\}$, where $X_{C}$ is the set of variables involved in the constraint $C$. This is implemented via scalar variables $\pi_{C}(\sigma)$ which are required to be non-negative and satisfy $\sum_{\sigma: X_{C} \rightarrow\{0,1\}} \pi_{C}(\sigma)=1$. For each variable $x$, two orthogonal vectors $\boldsymbol{v}_{(x, 0)}$ and $\boldsymbol{v}_{(x, 1)}$, corresponding to the events $x=0$ and $x=1$, are set up. The SDP requires for each variable $x, \boldsymbol{v}_{(x, 0)} \cdot \boldsymbol{v}_{(x, 1)}=0$ and $\boldsymbol{v}_{(x, 0)}+\boldsymbol{v}_{(x, 1)}=\boldsymbol{I}$ where $\boldsymbol{I}$ is a global unit vector. (In the integral solution, one of the vectors $\boldsymbol{v}_{(x, 1)}, \boldsymbol{v}_{(x, 0)}$ - based on the $x$ 's Boolean value - is intended to be $\boldsymbol{I}$ and the other one to be $\mathbf{0}$.)

Then, as constraint (5), the SDP does a consistency check: for two variables $x, y$ (that need not be distinct) involved in the same constraint $C$, and for every $b_{1}, b_{2} \in\{0,1\}$, the SDP insists that the inner product $\boldsymbol{v}_{\left(x, b_{1}\right)} \cdot \boldsymbol{v}_{\left(y, b_{2}\right)}$ equals $\operatorname{Pr}_{\sigma \in \pi_{C}}\left[\left(\sigma(x)=b_{1}\right) \wedge\left(\sigma(y)=b_{2}\right)\right]$.

Maximize

$$
\begin{align*}
\mathbf{E}_{C \in \mathcal{C}}\left[\mathbf{P r}_{\sigma \in \pi_{C}}[C(\sigma)=1]\right] & & &  \tag{1}\\
\boldsymbol{v}_{\left(x_{i}, 0\right)} \cdot \boldsymbol{v}_{\left(x_{i}, 1\right)} & =0 & & \forall i \in[n]  \tag{2}\\
\boldsymbol{v}_{\left(x_{i}, 0\right)}+\boldsymbol{v}_{\left(x_{i}, 1\right)} & =\boldsymbol{I} & & \forall i \in[n]  \tag{3}\\
\|\boldsymbol{I}\|^{2} & =1 & &  \tag{4}\\
\operatorname{Pr}_{\sigma \in \pi_{C}}\left[\sigma\left(x_{i}\right)=b_{1} \wedge \sigma\left(x_{j}\right)=b_{2}\right] & =\boldsymbol{v}_{\left(x_{i}, b_{1}\right)} \cdot \boldsymbol{v}_{\left(x_{j}, b_{2}\right)} & & \forall C \in \mathcal{C}, x_{i}, x_{j} \in C, b_{1}, b_{2} \in\{0,1\} \tag{5}
\end{align*}
$$

Subject to

Note that if we discard all the vectors by removing constraints (2)~(4), and changing constraints (5) to $\operatorname{Pr}_{\sigma \in \pi_{S}}\left[\sigma\left(x_{i}\right)=b_{1} \wedge \sigma\left(x_{j}\right)=b_{2}\right]=X_{\left(x_{i}, b_{1}\right),\left(x_{j}, b_{2}\right)}$, the SDP becomes a lifted LP in Sherali-Adams system. We call this LP scheme the lifted $L P$ in this paper.

The following striking theorem (Theorem 1.1 in [Rag08]) states that once we have an integrality gap for the canonical SDP, we also get a matching UG-hardness. Below and elsewhere in the paper, a $c$ vs. $s$ gap instance is an instance with SDP optimum at least $c$ and integral optimum at most $s$.

Theorem 5. Let $1>c>s>0$. If a constraint satisfaction problem $\Lambda$ admits a $c$ vs. s integrality gap instance for the above canonical SDP, then for every constant $\eta>0$, given an instance of $\Lambda$ that admits an assignment satisfying $(c-\eta)$ of constraints, it is UG-Hard to find an assignment satisfying more than $(s+\eta)$ of constraints.

To make our construction of integrality gaps easier, we notice the following simplification of the above SDP. Suppose we are given the global unit vector $\boldsymbol{I}$ and a vector $\boldsymbol{v}_{x}$ for each variable $x$ in the CSP instance, subject to the following constraints:

$$
\begin{align*}
\left(\boldsymbol{I}-\boldsymbol{v}_{x}\right) \cdot \boldsymbol{v}_{x} & =0 \quad \forall \text { variables } x  \tag{6}\\
\operatorname{Pr}_{\sigma \in \pi_{C}}\left[\sigma\left(x_{i}\right)=1 \wedge \sigma\left(x_{j}\right)=1\right] & =\boldsymbol{v}_{x_{i}} \cdot \boldsymbol{v}_{x_{j}} \quad \forall C \in \mathcal{C}, x_{i}, x_{j} \in C . \tag{7}
\end{align*}
$$

Defining $\boldsymbol{v}_{(x, 1)}=\boldsymbol{v}_{x}$ and $\boldsymbol{v}_{(x, 0)}=\boldsymbol{I}-\boldsymbol{v}_{x}$, it is easy to check that all constraints of the above SDP are satisfied. For instance, for variables $x, y$ belonging to a constraint $C$,

$$
\begin{aligned}
\boldsymbol{v}_{(x, 0)} \cdot \boldsymbol{v}_{(y, 1)} & =\left(\boldsymbol{I}-\boldsymbol{v}_{(x, 1)}\right) \cdot \boldsymbol{v}_{(y, 1)}=\left\|\boldsymbol{v}_{(y, 1)}\right\|^{2}-\boldsymbol{v}_{(x, 1)} \cdot \boldsymbol{v}_{(y, 1)} \\
& =\mathbf{P r}_{\sigma \in \pi_{C}}[\sigma(y)=1]-\mathbf{P r}_{\sigma \in \pi_{C}}[(\sigma(x)=1) \wedge(\sigma(y)=1)] \\
& =\operatorname{Pr}_{\sigma \in \pi_{C}}[(\sigma(x)=0) \wedge(\sigma(y)=1)],
\end{aligned}
$$

and other constraints of (5) follow similarly.

Henceforth in this paper, we will work with this streamlined canonical SDP with vector variables $\boldsymbol{I},\left\{\boldsymbol{v}_{x}\right\}$, scalar variables corresponding to the local distributions $\pi_{C}$, constraints (6) and (7), and objective function (1).

### 3.2 Overview of construction of SDP gaps

Horn-3SAT. In the concluding section of [Zwi98], Zwick remarks that there is an integrality gap for the LP he uses that matches his approximation ratio. Indeed such a LP gap is not hard to construct and we start by describing one such instance. The instance begins with clause $x_{1}$, and in the intermediate $(k-1)$ clauses, the $i$-th clause $x_{1} \wedge \cdots \wedge x_{i} \rightarrow x_{i+1}$ makes $x_{i+1}$ true if all the previous clauses are satisfied. Then the last clause $\overline{x_{k}}$ generates a contradiction. Thus the optimal integral solution is at most $(1-1 / k)$. On the other hand, one possible fractional solution starts with $x_{1}=(1-\varepsilon)$ for some $\varepsilon>0$. Then for $1 \leqslant i<k$, by letting $\left(1-x_{i+1}\right)=\sum_{j=1}^{i}\left(1-x_{j}\right)$, all the intermediate $(k-1)$ clauses are perfectly "satisfied" by the LP, while the gap $\left(1-x_{i+1}\right)=2^{i-1} \varepsilon$ increase exponentially. Thus by letting $\varepsilon=1 / 2^{k-2}$, we get $x_{k}=0$ and the LP solution is at least $\left(1-1 / 2^{\Omega(k)}\right)$. The instance gives a $\left(1-2^{-\Omega(k)}\right)$ vs. $(1-1 / k)$ LP integrality gap.

Now we convert this LP gap instance into an SDP gap instance in two steps. First, we reduce the arity of the instance from $k$ to 3 . Then, we find a set of vectors for the LP solution to make it an SDP solution.

For the first step, to get an instance of Max Horn-3SAT, we introduce $y_{i}$ which is intended to be $x_{1} \wedge \cdots \wedge x_{i-1}$. For $1 \leqslant i<k$, we replace the intermediate clauses by $x_{i} \wedge y_{i} \rightarrow x_{i+1}$, and add $x_{i} \wedge y_{i} \rightarrow y_{i+1}$ to meet the intended definition of $y_{i}$. We call each of these two clauses as comprising one step (the exact instance $\mathcal{I}_{k}^{\text {Horn }}$, which is slightly different for technical reasons mentioned below, can be found in Section 4.1.1). It is easy to show that for this instance there is a solution of value $\left(1-1 / 2^{\Omega(k)}\right)$ even for the lifted LP.

Finding vectors for the SDP turns out to be more challenging. Note that if we want to perfectly satisfy all the intermediate clauses in SDP, we need to obey $\boldsymbol{v}_{x_{i}} \cdot \boldsymbol{v}_{y_{i}} \leqslant\left\|\boldsymbol{v}_{x_{i+1}}\right\|^{2}$ and $\boldsymbol{v}_{x_{i}} \cdot \boldsymbol{v}_{y_{i}} \leqslant$ $\left\|\boldsymbol{v}_{y_{i+1}}\right\|^{2}$ for $1 \leqslant i<k$. Thus to make the norms $\left\|\boldsymbol{v}_{x_{i+1}}\right\|^{2}$ and $\left\|\boldsymbol{v}_{y_{i+1}}\right\|^{2}$ decrease fast (since we want $\left\|\boldsymbol{v}_{x_{k}}\right\|^{2}=\left\|\boldsymbol{v}_{y_{k}}\right\|^{2}=0$ ), we need to make the inner product $\boldsymbol{v}_{x_{i}} \cdot \boldsymbol{v}_{y_{i}}$ decrease fast as well. But technically it is hard to make both kinds of quantities decrease at a high rate for all intermediate clauses. Our solution is to decrease the norms and inner products alternately. More specifically, we divide the intermediate clauses into blocks, each of which contains two consecutive steps. In the first step of each block, we need that the inner product is much smaller than the norms so that we can decrease the norms quickly, but we preserve the value of inner product. Thus we cannot do this step repeatedly, and we need the second step, where we decrease the inner product (while preserving the norms) in preparation to start the first step of the next block.

1-in- $k$ Hitting Set. We use a simple symmetric instance as our gap instance. Ideally, the instance includes all subsets of the universe with size at most $k$ and we put uniform weights on sets of geometrically varying sizes (see Section 5.1 for our real gap instance which is slightly different). We first show that every subset intersects at most a (weighted) fraction $O(1 / \log k)$ of the sets exactly once. Then, to prove a much better SDP solution, in contrast to Max Horn-3SAT, the main effort is in finding a good solution for lifted LP. Once we get a good solution for lifted LP, because of symmetry, the norms for all variables are defined to be the same value, and the pairwise
inner products are also defined to be the same value. Then we only need to find vectors for a highly symmetric inner-product matrix, a step which is much easier than the counterpart of Max Horn-3SAT.

For the lifted LP, for each set in the instance, we place overwhelming weight on singleton subsets (only one variable is selected to be true) in all local distributions. This guarantees a good fractional solution. If we put all the weight on singletons though, the consistency check fails even for single-variable marginal distributions, whereas we need to ensure consistency of all pairwisevariable marginal distributions. Thus, for a feasible LP solution, we need to place some small weight on other subsets in order to obtain consistent marginal distributions. Indeed, we manage to generate a valid solution by giving an appropriate probability mass to the full set and all subsets of half the size in each local distribution.

### 3.3 Overview of algorithmic results

Our algorithmic results for Max Horn-2SAT and Max 1-in- $k$-SAT (Theorems 2 and 4 respectively) are obtained by rounding fractional solutions of appropriate linear programming (LP) relaxations.

The algorithm for Max Horn-2SAT is indeed a 2-approximation algorithm for Min Horn-2SAT Deletion problem (refer to Section 4.2 for the definition of Min Horn-2SAT Deletion). We prove a half-integrality property of the optimal solution to the natural LP relaxation of the problem, which can be viewed as an generalization of half-integrality property of (the natural LP for) Vertex Cover. We take the optimal solution of the natural LP relaxation, iteratively make every variable move towards half-integral values ( 0,1 , and $1 / 2$ ), while never increasing the value of the solution. This yields an optimal half-integral solution which can then be trivially rounded to obtain an integral solution that gives a factor 2 approximation.

For almost-satisfiable instances of Max 1 -in- $k$-SAT, we prove that randomized rounding (according to the fractional value of any optimal LP solution) gives a constant factor approximation. This gives a robust version of the algorithm in [GT05] which achieved a factor $1 / e$-approximation for (perfectly) satisfiable instances.

## 4 Approximability of Max Horn-3SAT

### 4.1 SDP gap and UG hardness for Max Horn-3SAT

### 4.1.1 Instance

We consider the following Max Horn-3SAT instance $\mathcal{I}_{k}^{\text {Horn }}$ parameterized by $k \geqslant 1$.
Start point:

$$
\begin{array}{ccrl}
\text { Block } i(0 \leqslant i \leqslant k-1) & \text { Step } i .1: & x_{2 i} \wedge y_{2 i} \rightarrow x_{2 i+1}, & x_{2 i} \wedge y_{2 i} \rightarrow y_{2 i+1} \\
& \text { Step } i .2: & x_{2 i+1} \wedge y_{2 i+1} \rightarrow x_{2 i+2}, & x_{2 i+1} \wedge y_{2 i+1} \rightarrow y_{2 i+2} \\
\text { End point: } & & x_{2 k} \wedge y_{2 k} \rightarrow & x_{2 k+1}, \\
& & x_{2 k} \wedge y_{2 k} \rightarrow y_{2 k+1} \\
& & \overline{y_{2 k+1}}
\end{array}
$$

It is easy to see this instance contains $(4 k+6)$ clauses, and cannot be completely satisfied. Thus we have:

Lemma 6. Every Boolean assignment satisfies at most a fraction $1-1 /(4 k+6)$ of the clauses of $\mathcal{I}_{k}^{\text {Horn }}$.

### 4.1.2 Construction of a good SDP solution

We will work with the SDP in simplified form described at the end of Section 3.1. Recall that the SDP requires a local distribution for each clause, and uses vectors to check the consistency on every pair of variables that belong to the clause. To construct a good solution for the SDP, we want to first find a good solution in the scalar part (i.e., local distributions), and then construct vectors which meet the consistency requirement. But it is difficult to construct a lot of vectors which meet all the requirements simultaneously. Thus, we break down the whole construction task into small pieces, each of which is easy to deal with. As long as there are solutions to these small pieces, and the solutions agree with each other on some interfaces, we can coalesce the small solutions together and come up with a global solution. The following definition and claim formally help us bring down the difficulty, and focus on one local block of variables at a time.

Definition 2 (partial solution). Let $\mathcal{C}^{\prime} \subseteq \mathcal{C}$ be a subset of clauses. $f=\left\{\pi_{C}=\pi_{C}(f), \boldsymbol{v}_{x}=\right.$ $\left.\boldsymbol{v}_{x}(f), \boldsymbol{I}=\boldsymbol{I}(f) \mid \forall C \in \mathcal{C}^{\prime}, x \in C\right\}$ is said to be a partial solution on $\mathcal{C}^{\prime}$, if all constraints of the SDP restricted to the subset of variables defined in $f$ are satisfied.

Claim 7. Let $\mathcal{C}_{1}, \mathcal{C}_{2} \subseteq \mathcal{C}$ be two disjoint set of clauses. Given $f$ and $g$ are partial solution on $\mathcal{C}_{1}, \mathcal{C}_{2}$ respectively. If for all $\boldsymbol{v}_{1}, \boldsymbol{v}_{2}$ (not necessarily distinct) defined in both $f$ and $g, \boldsymbol{v}_{1}(f) \cdot \boldsymbol{v}_{2}(f)=$ $\boldsymbol{v}_{1}(\mathrm{~g}) \cdot \boldsymbol{v}_{2}(\mathrm{~g})$, then there exists a partial solution, namely $h$, for $\mathcal{C}_{1} \cup \mathcal{C}_{2}$, such that $\forall C_{1} \in \mathcal{C}_{1}, C_{2} \in \mathcal{C}_{2}$, $\pi_{C_{1}}(h)=\pi_{C_{1}}(f), \pi_{C_{2}}(h)=\pi_{C_{2}}(g)$.

Proof. Let $X$ be the set of variables $x$ for which $\boldsymbol{v}_{x}(f)$ and $\boldsymbol{v}_{x}(g)$ are both defined. Denote $V_{f}=$ $\left\{\boldsymbol{v}_{x}(f) \mid x \in X\right\} \cup\{\boldsymbol{I}(f)\}$ and $V_{g}=\left\{\boldsymbol{v}_{x}(g) \mid x \in X\right\} \cup\{\boldsymbol{I}(g)\}$. Since the dot products of every pair of vectors in $V_{f}$ exactly equals the dot product between the corresponding pair in $V_{g}$, there is a rotation (orthogonal transformation) $T$ such that $\boldsymbol{I}(f)=T \boldsymbol{I}(g)$ and for all $x \in X, \boldsymbol{v}_{x}(f)=T \boldsymbol{v}_{x}(g)$.

Now define the partial solution $g^{\prime}$ as $\pi_{C}\left(g^{\prime}\right)=\pi_{C}(g)$ for all $C \in \mathcal{C}_{2}$ and $\boldsymbol{v}_{x}\left(g^{\prime}\right)=T \boldsymbol{v}_{x}(g)$, $\boldsymbol{I}\left(g^{\prime}\right)=T \boldsymbol{I}(g)$ for all $x \in C \in \mathcal{C}_{2}$. Obviously $f$ and $g^{\prime}$ agree on all the scalar and vector variables that are defined in both $f$ and $g^{\prime}$. Letting

$$
\boldsymbol{v}_{x}(h)=\left\{\begin{array}{ll}
\boldsymbol{v}_{x}(f) & x \in C \in \mathcal{C}_{1} \\
\boldsymbol{v}_{x}\left(g^{\prime}\right) & x \in C \in \mathcal{C}_{2}
\end{array}, \pi_{C}(h)= \begin{cases}\pi_{C}(f) & C \in \mathcal{C}_{1} \\
\pi_{C}\left(g^{\prime}\right) & C \in \mathcal{C}_{2}\end{cases}\right.
$$

it is easy to see $h$ is a partial solution on $\mathcal{C}_{1} \cup \mathcal{C}_{2}$.
By the above lemma, if we establish the following lemma which constructs a good partial solution on each block (the proof of which is deferred to Section 4.1.3), it is then easy to get a good global solution.

Lemma 8. For each Block $i(0 \leqslant i \leqslant k-1)$, each $0<c \leqslant 0.2$, let $r_{c}=1.5(1+c) /(1.5+c)$, and for each $0<p \leqslant \frac{1}{(1+c) r_{c}}$, there is a partial solution $f$ which completely satisfies all the clauses in Block i (by local distributions), and with following properties,

$$
\left\|\boldsymbol{v}_{x_{2 i}}(f)\right\|^{2}=\left\|\boldsymbol{v}_{y_{2 i}}(f)\right\|^{2}=1-p
$$

$$
\begin{aligned}
\boldsymbol{v}_{x_{2 i}}(f) \cdot \boldsymbol{v}_{y_{2 i}}(f) & =1-(1+c) p \\
\left\|\boldsymbol{v}_{x_{2 i+2}}(f)\right\|^{2}=\left\|\boldsymbol{v}_{y_{2 i+2}}(f)\right\|^{2} & =1-r_{c} p \\
\boldsymbol{v}_{x_{2 i+2}}(f) \cdot \boldsymbol{v}_{y_{2 i+2}}(f) & =1-(1+c) r_{c} p
\end{aligned}
$$

Together with Claim 7, we immediately get the following corollary.
Corollary 9. For the union of Block 0 to Block $k^{\prime}\left(0 \leqslant k^{\prime} \leqslant k-1\right)$, given parameters $0<c \leqslant 0.2$ and $0<p \leqslant \frac{1}{(1+c) r_{c}^{k^{\prime}+1}}$, there is a partial solution $g$ which completely satisfies all the clauses, and with following properties,

$$
\begin{aligned}
\left\|\boldsymbol{v}_{x_{0}}(g)\right\|^{2}=\left\|\boldsymbol{v}_{y_{0}}(g)\right\|^{2} & =1-p \\
\boldsymbol{v}_{x_{0}}(g) \cdot \boldsymbol{v}_{y_{0}}(g) & =1-(1+c) p \\
\left\|\boldsymbol{v}_{x_{2 k^{\prime}+2}}(g)\right\|^{2}=\left\|\boldsymbol{v}_{y_{2 k^{\prime}+2}}(g)\right\|^{2} & =1-r_{c}^{k^{\prime}+1} p \\
\boldsymbol{v}_{x_{2 k^{\prime}+2}}(g) \cdot \boldsymbol{v}_{y_{2 k^{\prime}+2}}(g) & =1-(1+c) r_{c}^{k^{\prime}+1} p
\end{aligned}
$$

Proof. Apply induction on $k^{\prime}$. The basis case $k^{\prime}=0$ is exactly Lemma 8 . For $k^{\prime}>0$, by induction hypothesis there is a partial solution $g^{\prime}$ satisfying all the clauses of the union of Blocks 0 to $k^{\prime}-1$ with the same parameter $c, p$. By Lemma 8 , there is a partial solution $f$ satisfying all the clauses of Block $k^{\prime}$ with parameter $c, r_{c}^{k^{\prime}} p$. Since $g^{\prime}$ and $f$ agree on pairwise inner-products over the definition of $\left\{\boldsymbol{v}_{x_{2 k^{\prime}}}, \boldsymbol{v}_{y_{2 k^{\prime}}}\right\}$, by Claim 7, there is a partial solution $g$ on the union of Blocks 0 to $k^{\prime}$ completely satisfying all the clauses.

With the above pieces in place, we now come to the final SDP solution.
Lemma 10. The optimal $S D P$ solution for the instance $\mathcal{I}_{k}^{\text {Horn }}$ has value at least $1-\frac{1}{(2 k+3) 1.05^{k}}$.
Proof. By Corollary 9, for any $0<c \leqslant 0.2$, by setting $p=\frac{1}{(1+c) r_{c}^{k}}$. There is a partial solution $g$ completely satisfying all the clauses of all the blocks, with

$$
\begin{aligned}
\left\|\boldsymbol{v}_{x_{0}}(g)\right\|^{2}=\left\|\boldsymbol{v}_{y_{0}}(g)\right\|^{2} & =1-\frac{1}{(1+c) r_{c}^{k}} \\
\left\|\boldsymbol{v}_{x_{2 k}}(g)\right\|^{2}=\left\|\boldsymbol{v}_{y_{2 k}}(g)\right\|^{2} & =c /(1+c) \\
\boldsymbol{v}_{x_{2 k}}(g) \cdot \boldsymbol{v}_{y_{2 k}}(g) & =0
\end{aligned}
$$

Based on $g$, we define a local distribution on two "Start point" clauses by making $x_{0}$ (or $y_{0}$ ) equal 1 with probability $1-p$. At "End point", we define the local distribution on clause $x_{2 k} \wedge y_{2 k} \rightarrow x_{2 k+1}$ as

$$
\begin{aligned}
& \operatorname{Pr}_{\pi}\left[x_{2 k}=1 \wedge y_{2 k}=0 \wedge x_{2 k+1}=0\right]=c /(1+c) \\
& \mathbf{P r}_{\pi}\left[x_{2 k}=0 \wedge y_{2 k}=1 \wedge x_{2 k+1}=0\right]=c /(1+c) \\
& \mathbf{P r}_{\pi}\left[x_{2 k}=0 \wedge y_{2 k}=0 \wedge x_{2 k+1}=0\right]=(1-c) /(1+c)
\end{aligned}
$$

And a similar distribution for the clause $x_{2 k} \wedge y_{2 k} \rightarrow y_{2 k+1}$ can be defined (by replacing $x_{2 k+1}$ by $y_{2 k+1}$ in the equations above). The distribution on clauses $\overline{x_{2 k+1}}$ and $\overline{y_{2 k+1}}$ never picks the
corresponding variable to be 1 . By defining $\boldsymbol{v}_{x_{2 k+1}}$ and $\boldsymbol{v}_{y_{2 k+1}}$ to be zero vectors, we note that the distributions are consistent with vectors. Thus the solution we construct is valid.

On the other hand, note that all the distributions locally satisfy the clauses, except for the distributions at "Start point" satisfy the corresponding clause with probability $1-\frac{1}{(1+c) r_{c}^{k}}$, thus the SDP solution is $1-\frac{2}{(4 k+6)(1+c) r_{c}^{k}}=1 \geqslant 1-\frac{1}{(2 k+3) r_{c}^{k}}$. By setting $c=0.2$, we get $r_{c} \geqslant 1.05$. Thus the best SDP solution is better than $1-\frac{1}{(2 k+3) 1.05^{k}}$.

Combining Lemma 6 and Lemma 10, we get the following theorem.
Theorem 11. $\mathcal{I}_{k}^{\text {Horn }}$ is a $(1-\varepsilon)$ vs. $(1-\Omega(1 / \log (1 / \varepsilon)))$ gap instance of Max Horn-3SAT for the canonical SDP relaxation.

Together with Theorem 5, Theorem 11 implies our main result, Theorem 1, on Max Horn-SAT.

### 4.1.3 Proof of the Key Lemma 8

For Block $i$, denote the clauses in Step $i .1$ by $C_{1 x}$ and $C_{1 y}$, and the clauses in Step $i .2$ by $C_{2 x}$ and $C_{2 y}$. We first construct partial solutions on Step $i .1$ and Step $i .2$ separately, as follows.

Partial solution on Step $i .1$ We first define a local distribution on $C_{1 x}$ as follows, and $C_{1 y}$ in a similar way (by replacing $x_{2 i+1}$ by $y_{2 i+1}$ in following equations).

$$
\begin{aligned}
& \mathbf{P r}_{\pi_{C_{1 x}}}\left[x_{2 i}=1 \wedge y_{2 i}=1 \wedge x_{2 i+1}=1\right]=1-(1+c) p \\
& \mathbf{P r}_{\pi_{C_{1 x}}}\left[x_{2 i}=1 \wedge y_{2 i}=0 \wedge x_{2 i+1}=0\right]=c p \\
& \mathbf{P r}_{\pi_{C_{1 x}}}\left[x_{2 i}=0 \wedge y_{2 i}=1 \wedge x_{2 i+1}=0\right]=c p \\
& \mathbf{P r}_{\pi_{C_{1 x}}}\left[x_{2 i}=0 \wedge y_{2 i}=0 \wedge x_{2 i+1}=1\right]=\left(1+c-r_{c}\right) p=\frac{(1+c) c}{1.5+c} \cdot p \\
& \mathbf{P r}_{\pi_{C_{1 x}}}\left[x_{2 i}=0 \wedge y_{2 i}=0 \wedge x_{2 i+1}=0\right]=\left(r_{c}-2 c\right) p=\frac{1.5-1.5 c-2 c^{2}}{1.5+c} \cdot p
\end{aligned}
$$

Recall $r_{c}=1.5(1+c) /(1.5+c)$. Note that all the probabilities are defined to be non-negative values by the range of $c$ and $p$, and they sum up to 1 .

We observe the following inner-product matrix $A$ over $\boldsymbol{I}, \boldsymbol{v}_{x_{2 i}}, \boldsymbol{v}_{y_{2 i}}, \boldsymbol{v}_{x_{2 i+1}}, \boldsymbol{v}_{y_{2 i+1}}$ is consistent with the local distributions on $C_{1 x}$ and $C_{1 y}$.

$$
A=\left[\begin{array}{ccccc}
1 & 1-p & 1-p & 1-r_{c} p & 1-r_{c} p \\
1-p & 1-p & 1-(1+c) p & 1-(1+c) p & 1-(1+c) p \\
1-p & 1-(1+c) p & 1-p & 1-(1+c) p & 1-(1+c) p \\
1-r_{c} p & 1-(1+c) p & 1-(1+c) p & 1-r_{c} p & 1-(1+c) p \\
1-r_{c} p & 1-(1+c) p & 1-(1+c) p & 1-(1+c) p & 1-r_{c} p
\end{array}\right]
$$

By Claim 22 in Appendix A we know that $A$ is positive semidefinite, and therefore there is a set of vectors consistent with our local distributions, i.e., we get a partial solution on Step i.1.

Partial solution on Step $i .2$ We define the local distribution on $C_{2 x}$ as follows. The distribution $C_{2 y}$ is defined in a similar way (by replacing $x_{2 i+2}$ with $y_{2 i+2}$ in the following equations). Let $q=r_{c} p$ and $\varepsilon=c / 1.5$.

$$
\begin{aligned}
& \operatorname{Pr}_{\pi_{C_{2 x}}}\left[x_{2 i+1}=1 \wedge y_{2 i+1}=1 \wedge x_{2 i+2}=1\right]=1-(1+\varepsilon) q \\
& \operatorname{Pr}_{\pi_{C_{2 x}}}\left[x_{2 i+1}=1 \wedge y_{2 i+1}=0 \wedge x_{2 i+2}=0\right]=\varepsilon q \\
& \operatorname{Pr}_{\pi_{C_{2 x}}}\left[x_{2 i+1}=0 \wedge y_{2 i+1}=1 \wedge x_{2 i+2}=0\right]=\varepsilon q \\
& \operatorname{Pr}_{\pi_{C_{2 x}}}\left[x_{2 i+1}=0 \wedge y_{2 i+1}=0 \wedge x_{2 i+2}=1\right]=\varepsilon q \\
& \operatorname{Pr}_{\pi_{C_{2 x}}}\left[x_{2 i+1}=0 \wedge y_{2 i+1}=0 \wedge x_{2 i+2}=0\right]=(1-2 \varepsilon) q .
\end{aligned}
$$

Note that all the probabilities are defined to be non-negative values by the range of $c$ and $p$, and they sum up to 1 .

Then note that the following inner-product matrix $B$ over $\boldsymbol{I}, \boldsymbol{v}_{x_{2 i+1}}, \boldsymbol{v}_{y_{2 i+1}}, \boldsymbol{v}_{x_{2 i+2}}, \boldsymbol{v}_{y_{2 i+2}}$ is consistent with the local distribution.

$$
B=\left[\begin{array}{ccccc}
1 & 1-q & 1-q & 1-q & 1-q \\
1-q & 1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q \\
1-q & 1-(1+\varepsilon) q & 1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q \\
1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q & 1-q & 1-(1+1.5 \varepsilon) q \\
1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q & 1-(1+1.5 \varepsilon) p & 1-q
\end{array}\right]
$$

Again by Claim 22 in Appendix A, B is positive semidefinite, and therefore there is a set of vectors consistent with local distributions - we have constructed a partial solution on Step i.2.

Combining the two partial solutions. It is easy to check with our parameter setting, partial solutions on Step $i .1$ and Step $i .2$ agree on pairwise inner-products between their shared vectors $\boldsymbol{I}, \boldsymbol{v}_{x_{2 i+1}}, \boldsymbol{v}_{y_{2 i+1}}$. Thus, there is a partial solution on Block $i$, with

$$
\begin{aligned}
\left\|\boldsymbol{v}_{x_{2 i} i}(f)\right\|^{2}=\left\|\boldsymbol{v}_{y_{2 i}}(f)\right\|^{2} & =1-p \\
\boldsymbol{v}_{x_{2 i}}(f) \cdot \boldsymbol{v}_{y_{2 i}}(f) & =1-(1+c) p \\
\left\|\boldsymbol{v}_{x_{2 i+2}}(f)\right\|^{2}=\left\|\boldsymbol{v}_{y_{2 i+2}}(f)\right\|^{2} & =1-q=1-r_{c} p \\
\boldsymbol{v}_{x_{2 i+2}}(f) \cdot \boldsymbol{v}_{y_{2 i+2}}(f) & =1-(1+1.5 \varepsilon) q=1-(1+c) r_{c} p
\end{aligned}
$$

### 4.2 Algorithm for Min Horn-2SAT Deletion and Max Horn-2SAT

In the Min Horn-2SAT Deletion problem, we are given a Horn-2SAT instance, and the goal is to find a subset of clauses of minimum total weight whose deletion makes the instance satisfiable. A factor 3 approximation algorithm for Min Horn-2SAT Deletion is given in [KSTW00]. Here we improve the approximation ratio to 2 . By a simple reduction from vertex cover, this is optimal under the UGC. Our motivation to study Min Horn-2SAT Deletion in the context of this paper is to pin down the fraction of clauses one can satisfy in a $(1-\varepsilon)$-satisfiable instance of Horn-2SAT: we can satisfy a fraction $(1-2 \varepsilon)$ of clauses (even in the weighted case), and satisfying a ( $1-c \varepsilon$ ) fraction is hard for $c<2$ assuming that vertex cover does not admit a $c$-approximation for any constant $c<2$.

In this section, we prove the following theorem by showing half-integrality of a natural LP relaxation for the problem.

Theorem 12. There is a polynomial-time 2-approximation algorithm for Min Horn-2SAT Deletion problem.

A direct corollary of Theorem 12 is the following result for approximating near-satisfiable instances of Max Horn-2SAT.

Theorem 2 (restated). Given a $(1-\varepsilon)$-satisfiable instance for Max Horn-2SAT, it is possible to find $a(1-2 \varepsilon)$-satisfying assignment efficiently.

### 4.2.1 LP Formulation

We find it slightly more convenient to present the algorithm for dual Horn-2SAT where each clause has at most one negated literal. (So the clauses are of the form $x, \bar{x}, x \vee y$, or $x \rightarrow y$, for variables $x, y$.) Let $w_{i j}^{(D)}>0$ be the weight imposed on the disjunction constraint $x_{i} \vee x_{j}$ (for each pair of $i, j$ such that $i<j$ ), and $w_{i j}^{(I)}>0$ be the weight imposed on the implication constraint $x_{i} \rightarrow x_{j}$ (for each pair of $i, j$ such that $i \neq j$ ). For each variable $x_{i}$, let $w_{i}^{(T)}$ be the weight on $x_{i}$ being true (i.e. $x_{i}=1$ ), and $w_{i}^{(F)}$ be the weight on $x_{i}$ being false (i.e. $x_{i}=0$ ). Then we write the following LP relaxation, where each real variable $y_{i}$ corresponds to the integer variable $x_{i}$.

| Minimize | $\sum_{i \in V} w_{i}^{(T)}\left(1-y_{i}\right)+\sum_{i \in V} w_{i}^{(F)} y_{i}+\sum_{i<j} w_{i j}^{(D)} z_{i j}^{(D)}+\sum_{i \neq j} w_{i j}^{(I)} z_{i j}^{(I)}$ |
| :--- | :--- |
| Subject to | $z_{i j}^{(D)} \geqslant 1-y_{i}-y_{j} \quad \forall i<j$ |
|  | $z_{i j}^{(I)} \geqslant y_{i}-y_{j}$ |$\quad \forall i \neq j$.

Let OPT be the optimal value of the integral solution, and $\mathrm{OPT}_{\mathrm{LP}}$ be the optimal value of the LP solution. We have $\mathrm{OPT}_{\mathrm{LP}} \leqslant \mathrm{OPT}$.

### 4.2.2 Half-integrality and rounding

Given a LP solution $f=\left\{z_{i j}^{(D)}, z_{i j}^{(I)}, y_{i}\right\}$, we can assume $z_{i j}^{(D)}=\max \left\{1-y_{i}-y_{j}, 0\right\}$ and $z_{i j}^{(I)}=$ $\max \left\{y_{i}-y_{j}, 0\right\}$ to minimize $\operatorname{Val}(f)$. Thus, we only need $f=\left\{y_{i}\right\}$ to characterize a solution, and we have
$\operatorname{Val}(f)=\sum_{i \in V} w_{i}^{(T)}\left(1-y_{i}\right)+\sum_{i \in V} w_{i}^{(F)} y_{i}+\sum_{i<j} w_{i j}^{(D)} \max \left\{1-y_{i}-y_{j}, 0\right\}+\sum_{i \neq j} w_{i j}^{(I)} \max \left\{y_{i}-y_{j}, 0\right\}$.
Lemma 13. There is a polynomial-time algorithm that, given a solution $f=\left\{y_{i}\right\}$ to the above LP, converts $f$ into another solution $f^{*}=\left\{y_{i}^{*}\right\}$ such that each $y_{i}^{*}$ is half-integral, i.e. $y_{i}^{*} \in\{0,1,1 / 2\}$, and $\operatorname{Val}\left(f^{*}\right) \leqslant \operatorname{Val}(f)$.

```
Algorithm 1 Round any LP solution \(f=\left\{y_{i}\right\}\) to a half-integral solution \(f^{*}\), with \(\operatorname{Val}\left(f^{*}\right) \leqslant \operatorname{Val}(f)\)
    while \(\exists i \in V: y_{i} \notin\{0,1,1 / 2\}\) do
        choose \(k \in V\), such that \(y_{k} \notin\{0,1,1 / 2\}\) (arbitrarily)
        if \(y_{k}<1 / 2\) then
            \(p \leftarrow y_{k}\)
        else
            \(p \leftarrow 1-y_{k}\)
        end if
        \(S \leftarrow\left\{i: y_{i}=p\right\}, S^{\prime} \leftarrow\left\{i: y_{i}=1-p\right\}\)
        \(a \leftarrow \max \left\{y_{i}: y_{i}<p, 1-y_{i}: y_{i}>1-p, 0\right\}, b \leftarrow \min \left\{y_{i}: y_{i}>p, 1-y_{i}: y_{i}<1-p, 1 / 2\right\}\)
        \(f^{(a)} \leftarrow\left\{y_{i}^{(a)}=a\right\}_{i \in S} \cup\left\{y_{i}^{(a)}=1-a\right\}_{i \in S^{\prime}} \cup\left\{y_{i}^{(a)}=y_{i}\right\}_{i \in V \backslash\left(S \cup S^{\prime}\right)}\)
        \(f^{(b)} \leftarrow\left\{y_{i}^{(b)}=b\right\}_{i \in S} \cup\left\{y_{i}^{(b)}=1-b\right\}_{i \in S^{\prime}} \cup\left\{y_{i}^{(b)}=y_{i}\right\}_{i \in V \backslash\left(S \cup S^{\prime}\right)}\)
        if \(\operatorname{Val}\left(f^{(a)}\right) \leqslant \operatorname{Val}\left(f^{(b)}\right)\) then
            \(f \leftarrow f^{(a)}\)
        else
            \(f \leftarrow f^{(b)}\)
        end if
    end while
    return \(f\) (as \(f^{*}\) )
```

Proof. We run Algorithm 1 whose input is the LP formulation and one of the solutions $f=\left\{y_{i}\right\}$, and whose output is the desired $f^{*}$.

It's easy to see that Algorithm 1 always maintains a valid solution $f$ to the LP (i.e., all variables $y_{i}$ 's are within the $[0,1]$ range). Then we only need to prove the following two things to show the correctness of Algorithm 1, 1) the while loop terminates (in linear steps), 2) in each loop, $\min \left\{\operatorname{Val}\left(f^{(a)}\right), \operatorname{Val}\left(f^{(b)}\right)\right\} \leqslant \operatorname{Val}(f)$, so that $\operatorname{Val}(f)$ never increases in the whole algorithm.

To prove the first point, we consider the set $W_{f}=\left\{0<y<1 / 2: \exists i \in V\right.$, s.t. $\left.y=y_{i} \vee y=1-y_{i}\right\}$. In each loop, the algorithm picks a $p$ from $W_{f}$. At the end of the loop, we see that $p$ is wiped from $W_{f}$ while no new elements are added. Thus, after linear steps of the loop, $W_{f}$ becomes $\emptyset$ and the loop terminates.

For the second point, we define $f^{(t)}=\left\{y_{i}^{(t)}=t\right\}_{i \in S} \cup\left\{y_{i}^{(t)}=1-t\right\}_{i \in S^{\prime}} \cup\left\{y_{i}^{(t)}=y_{i}\right\}_{i \in V \backslash\left(S \cup S^{\prime}\right)}$ for $t \in[a, b]$ at Line 9 in the algorithm. Then if we can show $\operatorname{Val}\left(f^{(t)}\right)$ is a linear function within $t \in[a, b]$, together with the fact $p \in[a, b]$, we shall conclude that $\min \left\{\operatorname{Val}\left(f^{(a)}\right), \operatorname{Val}\left(f^{(b)}\right)\right\} \leqslant \operatorname{Val}\left(f^{(p)}\right)=$ $\operatorname{Val}(f)$. To prove the linearity of $\operatorname{Val}\left(f^{(t)}\right)$, we only need to show that $g_{1}(t)=\max \left\{1-y_{i}^{(t)}-y_{j}^{(t)}, 0\right\}$ and $g_{2}(t)=\max \left\{y_{i}^{(t)}-y_{j}^{(t)}, 0\right\}$ are linear with the respect to $t \in[a, b]$, for any possible $i, j$. Thus we discuss the following five cases.

- $i, j \in V \backslash\left(S \cup S^{\prime}\right)$. In this case, $g_{1}$ and $g_{2}$ are constant functions.
- $i \in V \backslash\left(S \cup S^{\prime}\right), j \in S \cup S^{\prime}$. In this case, the only "non-linear point" is at $t=1-y_{i}$ for $g_{1}$ and $t=y_{i}$ for $g_{2}$. But these two points are away from $[a, b]$.
- $i \in S \cup S^{\prime}, j \in V \backslash\left(S \cup S^{\prime}\right)$. Similar argument works as the previous case.
- $i \in S, j \in S^{\prime}$ (or $i \in S^{\prime}, j \in S$ ). In this case, $1-y_{i}^{(t)}-y_{j}^{(t)}=0$ always holds for $t \in[a, b]$ and therefore $g_{1}$ is constant function. On the other hand, since $y_{i}^{(t)} \leqslant y_{j}^{(t)}\left(\right.$ or $\left.y_{i}^{(t)} \geqslant y_{j}^{(t)}\right)$, we also have $g_{2}(t)=0\left(\right.$ or $\left.g_{2}(t)=y_{i}^{(t)}-y_{j}^{(t)}=1-2 t\right)$ being linear.
- $i, j \in S$ (or $i, j \in S^{\prime}$ ). In this case, $y_{i}^{(t)}=y_{j}^{(t)}$ always holds for $t \in[a, b]$ and therefore $g_{2}$ is constant function. On the other hand, since $y_{i}^{(t)}+y_{j}^{(t)} \leqslant 1$ (or $y_{i}^{(t)}+y_{j}^{(t)} \geqslant 1$ ), we also have $g_{1}(t)=1-y_{i}^{(t)}-y_{j}^{(t)}=1-2 t$ (or $g_{1}(t)=0$ ) being linear.

A direct corollary of Lemma 13 is the following.
Corollary 14. There is a polynomial-time algorithm to get a solution $f$ such that $\operatorname{Val}(f)=\mathrm{OPT}_{\mathrm{LP}}$ and the variables in $f$ are half-integral (i.e. being one of 0,1 , and $1 / 2$ ).

Now we are ready for the proof of Theorem 12.
Proof of Theorem 12. Apply Corollary 14 to get an optimal LP solution $f=\left\{y_{i}\right\}$ which has halfintegral values. Then define $f_{\text {int }}=\left\{x_{i}\right\}$ as follows. For each $i \in V$, let $x_{i}=1$ when $y_{i} \geqslant 1 / 2$, and $x_{i}=0$ when $y_{i}=0$. We observe that

- $x_{i} \leqslant 2 y_{i}$ and $1-x_{i} \leqslant 1-y_{i}$ for each $i \in V$.
- For each $i<j$, we have $\max \left\{1-x_{i}-x_{j}, 0\right\} \leqslant \max \left\{1-y_{i}-y_{j}, 0\right\}$ since $x_{i} \geqslant y_{i}$ and $x_{j} \geqslant y_{j}$.
- For each $i \neq j$, we see that when $\max \left\{y_{i}-y_{j}, 0\right\}=0 \Rightarrow y_{i} \leqslant y_{j}$, we always have $x_{i} \leqslant x_{j} \Rightarrow$ $\max \left\{x_{i}-x_{j}, 0\right\}=0$. On the other hand, when $\max \left\{y_{i}-y_{j}, 0\right\}>0 \Rightarrow \max \left\{y_{i}-y_{j}, 0\right\} \geqslant 1 / 2$, we have $\max \left\{x_{i}-x_{j}, 0\right\} \leqslant 1 \leqslant 2 \max \left\{y_{i}-y_{j}, 0\right\}$.

Altogether, we have

$$
\begin{aligned}
\operatorname{Val}\left(f_{\text {int }}\right) & =\sum_{i \in V} w_{i}^{(T)}\left(1-x_{i}\right)+\sum_{i \in V} w_{i}^{(F)} x_{i}+\sum_{i<j} w_{i j}^{(D)} \max \left\{1-x_{i}-x_{j}, 0\right\}+\sum_{i \neq j} w_{i j}^{(I)} \max \left\{x_{i}-x_{j}, 0\right\} \\
& \leqslant \sum_{i \in V} w_{i}^{(T)}\left(1-y_{i}\right)+\sum_{i \in V} w_{i}^{(F)} 2 y_{i}+\sum_{i<j} w_{i j}^{(D)} \max \left\{1-y_{i}-y_{j}, 0\right\}+\sum_{i \neq j} w_{i j}^{(I)} 2 \max \left\{y_{i}-y_{j}, 0\right\} \\
& \leqslant 2 \operatorname{Val}(f)=2 \mathrm{OPT}_{\mathrm{LP}} \leqslant 2 \mathrm{OPT} .
\end{aligned}
$$

## 5 Inapproximability and approximation algorithm for Max $1-\mathrm{in}-k$ HS

### 5.1 SDP gap and UG-hardness for Max 1-in- $k$-HS

In this section, we construct an SDP gap for Max 1-in- $k$-HS, and prove Theorem 3, which is restated as follows.

Theorem 3 (restated). For some absolute constant $C^{\prime}>0$, for every $\alpha>0$, given a $\left(1-1 / k^{1-\alpha}\right)$ satisfiable instance of Max 1-in- $k$-HS, it is UG-hard to find a subset intersecting more than a fraction $\frac{C^{\prime}}{\alpha \log k}$ of the sets exactly once.

We start by constructing the gap instance.
Instance. We define the (weighted) instance of Max 1 -in- $k$ - HS , denoted $\mathcal{I}^{\mathrm{EHS}}(m, n, \varepsilon)$, parameter $0<\varepsilon<1, m \geqslant 2$ and $n \geqslant \varepsilon m \cdot 2^{2\lceil m(1+\varepsilon)\rceil}$ as follows.

- The universe $U=[n]=\{1,2, \ldots, n\}$.
- Define the sets $\mathcal{C}$ by choosing $t \in m+1, m+2, \cdots,\lceil m(1+\varepsilon)\rceil$ uniform randomly, and picking a subset $S \subseteq U$ with size $2^{t}$ by random, then letting $S \in \mathcal{C}$ and the weight of $S$ be the corresponding probability.

Note that in such an instance, the size of $S_{i}$ is at most $k=2^{\lceil m(1+\varepsilon)\rceil}$.

### 5.1.1 Upper bound of optimal integral solution

In this section, we prove the following Lemma showing that the above instance does not have a good exact hitting set.

Lemma 15. There is a constant $C_{1}$ such that for all $0<\varepsilon<1, m \geqslant 2$ and $n \geqslant \varepsilon m \cdot 2^{2\lceil m(1+\varepsilon)\rceil, ~}$ the optimal solution to $\mathcal{I}^{\mathrm{EHS}}(m, n, \varepsilon)$ has value at most $C_{1} /(\varepsilon \log k)$.

We begin with the following two statements that will be useful in bounding the value of any integral solution to $\mathcal{I}^{\mathrm{EHS}}(m, n, \varepsilon)$.

Lemma 16. Suppose the hitting set $V \subseteq U$ is of size $l$. Then the probability that a size- $z(2 \leqslant z \leqslant$ $l / 2)$ set is hit exactly once by $V$, is at most $\frac{2 z}{n} \cdot l \cdot\left(\frac{1}{e}\right)^{z l / 4 n}$.

Proof.

$$
\begin{aligned}
\operatorname{Pr}_{S \in \mathcal{C}}[|S \cap V|=1| | S \mid=z] & =\frac{l\binom{n-l}{z-1}}{\binom{n}{z}} \\
& =l \cdot \frac{(n-l)!(n-z)!z!}{(n-l-z+1)!(z-1)!n!} \\
& =z l \cdot \frac{(n-l)!}{(n-l-z+1)!} \cdot \frac{(n-z)!}{n!} \\
& \leqslant z l \cdot \frac{(n-l)^{z-1}}{(n-z)^{z}} \\
& =\frac{z}{n-z} \cdot l \cdot\left(1-\frac{l-z}{n-z}\right)^{z-1} \\
& \leqslant \frac{z}{n-z} \cdot l \cdot\left(\frac{1}{e}\right)^{(z-1)(l-z) /(n-z)}
\end{aligned}
$$

$$
\begin{array}{lr}
\leqslant \frac{z}{n-z} \cdot l \cdot\left(\frac{1}{e}\right)^{z l / 4(n-z)} & (2 \leqslant z \leqslant l / 2) \\
\leqslant \frac{2 z}{n} \cdot l \cdot\left(\frac{1}{e}\right)^{z l / 4 n} & (z \leqslant l / 2 \leqslant n / 2)
\end{array}
$$

Claim 17. For all $x>0$ and $m \in \mathbb{N}^{+}$, we have

$$
\sum_{i=1}^{m} 2^{i} x e^{-2^{i} x} \leqslant 2 / \ln 2 .
$$

Proof.

$$
\begin{aligned}
\sum_{i=1}^{m} 2^{i} x e^{-2^{i} x} & \leqslant \sum_{i=-\infty}^{+\infty} 2^{i} x e^{-2^{i} x} \\
& =\sum_{i=-\infty}^{\left[\log _{2} 1 / x\right]} 2^{i} x e^{-2^{i} x}+\sum_{i=\left[\log _{2} 1 / x\right]+1}^{+\infty} 2^{i} x e^{-2^{i} x} \\
& \leqslant \int_{-\infty}^{\left[\log _{2} 1 / x\right]+1} 2^{y} x e^{-2^{y} x} d y+\int_{\left[\log _{2} 1 / x\right]}^{+\infty} 2^{y} x e^{-2^{y} x} d y \quad \text { (monotonicity) } \\
& \leqslant 2 \int_{-\infty}^{+\infty} 2^{y} x e^{-2^{y} x} d y \\
& =\frac{2}{\ln 2}
\end{aligned}
$$

We can now prove Lemma 15.
Proof of Lemma 15. Set $C_{1}=\max \{32 / \ln 2,12\}$. Given a solution $V$, let $l=|V|$. If $l \geqslant 2 \cdot 2^{[m(1+\varepsilon)]}$, then $l \geqslant 2|S|, \forall S \in \mathcal{C}$. In this case, the probability that $S \in \mathcal{C}$ is hit exactly once, is

$$
\begin{align*}
\operatorname{Pr}_{S \in \mathcal{C}}[|S \cap V|=1] & =\sum_{t=m+1}^{\lceil m(1+\varepsilon)\rceil} \operatorname{Pr}_{S \in \mathcal{C}}\left[|S|=2^{t}\right] \cdot \operatorname{Pr}_{S \in \mathcal{C}}\left[|S \cap V|=1| | S \mid=2^{t}\right] \\
& =\frac{1}{\varepsilon m} \sum_{t=m+1}^{\lceil m(1+\varepsilon)\rceil} \operatorname{Pr}_{S \in \mathcal{C}}\left[|S \cap V|=1| | S \mid=2^{t}\right] \\
& \leqslant \frac{1}{\varepsilon m} \sum_{t=m+1}^{\lceil m(1+\varepsilon)\rceil} \frac{2 \cdot 2^{t}}{n} \cdot l \cdot\left(\frac{1}{e}\right)^{2^{t} / / 4 n} \quad \text { (by Lemma 16) } \\
& \leqslant \frac{1}{\varepsilon m} \cdot \frac{16}{\ln 2}  \tag{byClaim17}\\
& \leqslant C_{1} /(\varepsilon \log k) .
\end{align*} \quad \text { (by Claim 17 } \quad \text {. } \quad \text {. }
$$

On the other hand, if $l<2 \cdot 2^{\lceil m(1+\varepsilon)\rceil}$, then

$$
\begin{aligned}
\operatorname{Pr}_{S \in \mathcal{C}}[|S \cap V|=1] & \leqslant \operatorname{Pr}_{S \in \mathcal{C}}[|S \cap V| \geqslant 1] \\
& \leqslant \operatorname{Pr}_{S \in \mathcal{C}}\left[|S \cap V| \geqslant 1| | S \mid=2^{[m(1+\varepsilon)]}\right] \\
& =1-\binom{n-l}{2^{\lceil m(1+\varepsilon)\rceil}} /\binom{n}{2^{\lceil m(1+\varepsilon)\rceil}} \\
& =1-\frac{(n-l)!}{\left(n-l-2^{\lceil m(1+\varepsilon)\rceil)!}\right.} \cdot \frac{\left(n-2^{\lceil m(1+\varepsilon)\rceil}\right)!}{n!} \\
& \leqslant 1-\left(\frac{n-l-2^{\lceil m(1+\varepsilon)\rceil}}{n}\right)^{l} \\
& \leqslant 1-\left(\frac{n-3 \cdot 2^{\lceil m(1+\varepsilon)\rceil}}{n}\right)^{2 \cdot 2^{[m(1+\varepsilon)\rceil}} \quad\left(l<2 \cdot 2^{\lceil m(1+\varepsilon)\rceil}\right) \\
& \leqslant \frac{6 \cdot 2^{2\lceil m(1+\varepsilon)\rceil}}{n} \quad\left(\forall 0 \leqslant x \leqslant 1, y \geqslant 0,(1-x)^{y} \geqslant 1-x y\right) \\
& \leqslant 6 / \varepsilon m \leqslant C_{1} /(\varepsilon \log k) .
\end{aligned}
$$

And this proves the lemma.

### 5.1.2 Construction of good SDP solution

We prove that the canonical SDP has a solution with value close to 1 .
Lemma 18. For the Max 1-in- $k$ - HS instance $\mathcal{I}^{\mathrm{EHS}}(m, n, \varepsilon)$, the optimal solution to the canonical SDP has value at least $1-4 / 2^{m} \geqslant 1-4 / k^{1-\varepsilon}$.

To prove Lemma 18, recall the canonical SDP for Max 1 -in- $k$-HS as follows.

$$
\begin{array}{rlll}
\text { Maximize } & \mathbf{E}_{S \in \mathcal{C}}\left[\mathbf{P r}_{\sigma \in \pi_{S}}\left[\left|\sigma^{-1}(1)\right|=1\right]\right] & & \\
\text { Subject to } & \boldsymbol{v}_{i} \cdot \boldsymbol{I} & =\left\|\boldsymbol{v}_{i}\right\|^{2} & \\
\|\boldsymbol{I}\|^{2} & =1, j \in U \\
& & \forall i \in U \\
& \operatorname{Pr}_{\sigma \in \pi_{S}}[\sigma(i)=1] & =\left\|\boldsymbol{v}_{i}\right\|^{2} & \\
& \forall S \in \mathcal{C}, i \in S \\
\operatorname{Pr}_{\sigma \in \pi_{S}}[\sigma(i)=1 \wedge \sigma(j)=1] & =\boldsymbol{v}_{i} \cdot \boldsymbol{v}_{j} & & \forall S \in \mathcal{C}, i \neq j \in S
\end{array}
$$

Now, we exhibit an SDP solution for the instance $\mathcal{I}^{\mathrm{EHS}}(m, n, \varepsilon)$ that has value close to 1 . We first construct the scalars, and then the vectors.

Constructing the solution - scalars. Let $M=2^{m}, p=2 / M, q=1 / M . p$ and $q$ will be the marginal probability for single element pairs. and For each $S \in \mathcal{C}$, and each $\sigma: S \rightarrow\{0,1\}$, define the local distribution $\pi_{S}$ as follows:

Given $M<|S|$ for all $S \in \mathcal{C}$, it is easy to check $\pi_{S}$ is always non-negative. And it can be checked that $\sum_{\sigma \subseteq S} \pi_{S}(\sigma)=1$. Thus, $\pi_{S}$ is a valid probability distribution.

Then we calculate the following values which are related to the SDP.

- For all $i \in S \in \mathcal{C}$,

$$
\begin{aligned}
\operatorname{Pr}_{\sigma \in \pi_{S}}[\sigma(i)=1]= & 1-\frac{|S|-1}{|S|} \cdot \frac{|S|}{|S|-2} \cdot\left(\frac{|S|}{|S|-1}-\frac{3|S|-2}{|S|-1} \cdot p+2 q\right) \\
& -\frac{1}{2} \cdot \frac{4}{|S|-2} \cdot((|S|-1)(p-q)-(1-p)) \\
= & p .
\end{aligned}
$$

- For all $i \neq j \in S \in \mathcal{C}$,

$$
\begin{aligned}
\operatorname{Pr}_{\sigma \in \pi_{S}}[\sigma(i)=1 \wedge \sigma(j)=1]= & \left(\frac{1}{|S|-1}-\frac{|S|}{|S|-1} \cdot p+2 q\right) \\
& +\left(1-\frac{|S| / 2-1}{2(|S|-1)}\right) \cdot \frac{4}{|S|-2} \cdot((|S|-1)(p-q)-(1-p)) \\
= & q .
\end{aligned}
$$

Constructing the solution - vectors. Now we need to show there exists a set of vectors passing the consistency check on local distributions we defined above. In fact, we show there exists set of vectors satisfying even stricter requirements, where the inner-product between every pair of vectors is defined, as follows,

$$
\begin{aligned}
\left\|\boldsymbol{v}_{i}\right\|^{2} & =p & & \forall i \in U \\
\boldsymbol{v}_{i} \cdot \boldsymbol{v}_{j} & =q & & \forall i \neq j \in U \\
\boldsymbol{v}_{i} \cdot \boldsymbol{I} & =\left\|\boldsymbol{v}_{i}\right\|^{2} & & \forall i \in U \\
\|\boldsymbol{I}\|^{2} & =1 & &
\end{aligned}
$$

Thus we only need to show the corresponding inner-product matrix is positive semidefinite. The matrix is in the form of

$$
A=\left[\begin{array}{cc}
1 & p \boldsymbol{b}^{T} \\
p \boldsymbol{b} & (p-q) I+q J
\end{array}\right]
$$

where $\boldsymbol{b}$ is $n \times 1$ all-one vector, $J$ is the $n \times n$ all-one matrix, and $I$ is the identity matrix.
Given $\boldsymbol{x}=\left(x_{0}, x_{1}, \cdots, x_{n}\right) \in \mathbb{R}^{n}$,

$$
\begin{aligned}
\boldsymbol{x}^{T} A \boldsymbol{x} & =\left(x_{0}, x_{1}, \cdots, x_{n}\right)\left[\begin{array}{cc}
1 & p \boldsymbol{b}^{T} \\
p \boldsymbol{b} & (p-q) I+q J
\end{array}\right]\left(x_{0}, x_{1}, \cdots, x_{n}\right)^{T} \\
& =x_{0}^{2}+2 p x_{0}\left(\sum_{i=1}^{n} x_{i}\right)+q\left(\sum_{i=1}^{n} x_{i}\right)^{2}+(p-q) \sum_{i=1}^{n} x_{i}^{2}
\end{aligned}
$$

Note that this quadratic form is always non-negative when $p \geqslant q$ and $4 p^{2}-4 q \leqslant 0 \Leftrightarrow q \geqslant p^{2}$. Our $p=2 / M$ and $q=1 / M$ satisfies these conditions. Therefore the inner-product matrix is positive semidefinite and the vectors exist.

Now we can prove Lemma 18, which says the optimal SDP solution has value close to 1 .
Proof of Lemma 18. The value of the solution we exhibited above is

$$
\begin{aligned}
\mathbf{E}_{S \in \mathcal{C}}\left[\mathbf{P r}_{\sigma \in \pi_{S}}\left[\left|\sigma^{-1}(1)\right|=1\right]\right] & =\mathbf{E}_{S \in \mathcal{C}}\left[\sum_{\sigma: S \rightarrow\{0,1\},|\sigma|^{-1}(1)=1} \pi_{S}(\sigma)\right] \\
& =\mathbf{E}_{S \in \mathcal{C}}\left[\frac{|S|}{|S|-2} \cdot\left(\frac{|S|}{|S|-1}-\frac{3|S|-2}{|S|-1} \cdot p+2 q\right)\right] \\
& \geqslant \mathbf{E}_{S \in \mathcal{C}}\left[\frac{|S|}{|S|-1}-\frac{3|S|-2}{|S|-1} \cdot p+2 q\right] \\
& =\mathbf{E}_{S \in \mathcal{C}}[1-3 p+2 q+(1-p) /(|S|-1)] \\
& \geqslant \mathbf{E}_{S \in \mathcal{C}}[1-3 p+2 q] \\
& =1-3 p+2 q=1-4 / M .
\end{aligned}
$$

Together with Theorem 5, Lemmas 15 and 18 imply Theorem 3.

### 5.2 A robust algorithm for almost-satisfiable Max 1-in- $k$-SAT

In this section, we prove the following theorem.
Theorem 4 (restated). For every constant $B>1$, the following holds. There is a polynomial time algorithm that given a $\left(1-\frac{1}{B k}\right)$-satisfiable instance of $\operatorname{Max} 1-\mathrm{in}-k$-SAT, finds a truth-assignment on variables satisfying exactly one term for a fraction $\lambda$ of the clauses, where $\lambda=\left(\frac{1-1 / \sqrt{B}}{e}\right)^{2}$.

The algorithm is based on rounding an LP relaxation for the problem, and gives a robust version of the algorithm in [GT05] which achieved a factor $1 / e$-approximation for (perfectly) satisfiable instances.

Given a truth-assignment $\sigma$ and a clause $C$, we denote $\sigma \cap C$ by the set of terms in $C$ satisfied by $\sigma$. Our algorithm first solves the following LP relaxation of the problem.

$$
\begin{array}{rrr}
\text { Maximize } & \mathbf{E}_{C \in \mathcal{C}}\left[\mathbf{P r}_{\sigma \in \pi_{C}}[|\sigma \cap C|=1]\right] \\
\text { Subject to } & \operatorname{Pr}_{\sigma \in \pi_{C}}[\sigma(i)=1]=x_{i} & \forall C \in \mathcal{C}, i \in C
\end{array}
$$

Given a solution $\left\{\pi_{C}\right\}$ and $\left\{x_{i}\right\}$ to the LP, we generate an assignment $\tau$ by for each $i \in U$ letting $\tau\left(x_{i}\right)=1$ with probability $x_{i}$. Then we prove the following lemma which directly implies Theorem 4.

Lemma 19. For every constant $B>1$, when $\mathrm{OPT}_{\mathrm{LP}}>1-\frac{1}{B k}$, we have

$$
\mathbf{E}_{\tau}\left[\operatorname{Pr}_{C \in \mathcal{C}}[|\tau \cap C|=1]\right] \geqslant\left(\frac{1-1 / \sqrt{B}}{e}\right)^{2}
$$

Proof. Given $\mathbf{E}_{C \in \mathcal{C}}\left[\operatorname{Pr}_{\sigma \in \pi_{C}}[|\sigma \cap C|=1]\right]>1-\frac{1}{B k}$, by an averaging argument, we know that for at least $(1-1 / \sqrt{B})$ fraction of $C \in \mathcal{C}$ are "good", i.e., for these $C$ clauses, we have

$$
\operatorname{Pr}_{\sigma \in \pi_{C}}[|\sigma \cap C|=1]>1-\frac{1}{\sqrt{B} k} .
$$

For each good $C \in \mathcal{C}$, and for each term $t \in C$, let $p(t)=x_{i}$ if $t=x_{i}$, or $p(t)=1-x_{i}$ if $t=\overline{x_{i}}$, i.e. $p(t)$ is the probability that $t$ is satisfied by $\tau$. Then we know that

$$
\sum_{t \in C} p(t)=\mathbf{E}_{\sigma \in \pi_{C}}[|\sigma \cap C|] \geqslant \operatorname{Pr}_{\sigma \in \pi_{C}}[|\sigma \cap C|=1] \geqslant 1-\frac{1}{\sqrt{B} k} .
$$

On the other hand,

$$
\begin{equation*}
\sum_{t \in C} p(t)=\mathbf{E}_{\sigma \in \pi_{C}}[|\sigma \cap C|] \leqslant \operatorname{Pr}_{\sigma \in \pi_{C}}[|\sigma \cap C|=1]+\left(1-\operatorname{Pr}_{\sigma \in \pi_{C}}[|\sigma \cap C|=1]\right)|C| \leqslant 1+1 / \sqrt{B} \tag{8}
\end{equation*}
$$

We now lower bound the probability that $\tau$ satisfies $C$, using the Lemma 20 proved at the end of the section. We discuss the following two cases to establish the lower bound.

Case 1. If all the terms in $C$ depend on distinct variables, then

$$
\begin{equation*}
\operatorname{Pr}_{\tau}[|\tau \cap C|=1]=\sum_{t \in C} p(t) \prod_{t^{\prime} \in C, t \neq t^{\prime}}\left(1-p\left(t^{\prime}\right)\right) . \tag{9}
\end{equation*}
$$

For good $C$ we know that $\sum_{t \in C} p(t) \in\left[1-\frac{1}{\sqrt{B} k}, 1+1 / \sqrt{B}\right] \subseteq[1-1 / \sqrt{B}, 1+1 / \sqrt{B}]$, By Lemma 20 given right after this proof, we know that $(9) \geqslant(1-1 / \sqrt{B}) / e^{2}$.

Case 2. If some terms in $C$ depend on the same variable, i.e. $\exists i: x_{i}, \overline{x_{i}} \in C$, then by (8) we know that $\sum_{t \in C \backslash\left\{x_{i}, \overline{x_{i}}\right\}} \leqslant 1 / \sqrt{B}<1$. Thus terms in $C \backslash\left\{x_{i}, \overline{x_{i}}\right\}$ depend on distinct variables, and also by Lemma 20, we know that

$$
\operatorname{Pr}_{\tau}[|\tau \cap C|=1]=1 \cdot \prod_{t \in C \backslash\left\{x_{i}, \overline{x_{i}}\right\}}(1-p(t)) \geqslant(1-1 / \sqrt{B}) / e^{2} .
$$

Combining the two cases above, we get

$$
\begin{aligned}
\mathbf{E}_{\tau}\left[\operatorname{Pr}_{C \in \mathcal{C}}[|\tau \cap C|=1]\right] & =\mathbf{E}_{C \in \mathcal{C}}\left[\mathbf{P r}_{\tau}[|\tau \cap C|=1]\right] \\
& \geqslant(1-1 / \sqrt{B}) \mathbf{E}_{C \in \mathcal{C}}\left[\operatorname{Pr}_{\tau}[|\tau \cap C|=1] \mid C \text { is good }\right] \\
& \geqslant\left(\frac{1-1 / \sqrt{B}}{e}\right)^{2} .
\end{aligned}
$$

It remains to prove the following inequality which was used in the above proof.
Lemma 20. Given $x_{1}, x_{2}, \cdots, x_{n} \in[0,1]$, and $1-\varepsilon \leqslant \sum_{i} x_{i} \leqslant 1+\varepsilon$ where $\varepsilon<1$ then

$$
\sum_{i} x_{i} \prod_{j \neq i}\left(1-x_{j}\right) \geqslant \frac{1-\varepsilon}{e^{2}}
$$

Proof. We use the following claim to prove this lemma.
Claim 21. For $n \geqslant 2$, given a set of $n$ numbers $\left\{x_{i}\right\}$ as described in the lemma, the objective function $\sum_{i} x_{i} \prod_{j \neq i}\left(1-x_{j}\right)$ is minimized when

- All the $x_{i}$ 's are the same, or
- $\exists i: x_{i}=0$ or $\exists i: x_{i}=1$.

Proof. Suppose the first condition doesn't hold, we prove the second one holds. Without loss of generality assume that $x_{1} \neq x_{2}$. Then rewrite the objective function as

$$
\begin{aligned}
& \sum_{i} x_{i} \prod_{j \neq i}\left(1-x_{j}\right) \\
= & \left(1-x_{1}\right)\left(1-x_{2}\right)\left(\sum_{i \geqslant 3} x_{i} \prod_{j \neq i, j \geqslant 3}\left(1-x_{j}\right)\right)+\left(x_{1}\left(1-x_{2}\right)+x_{2}\left(1-x_{1}\right)\right) \prod_{j \geqslant 3}\left(1-x_{j}\right)
\end{aligned}
$$

Let $C_{1}=\sum_{i \geqslant 3} x_{i} \prod_{j \neq i, j \geqslant 3}\left(1-x_{j}\right)$ and $C_{2}=\prod_{j \geqslant 3}\left(1-x_{j}\right)$, we have

$$
\sum_{i} x_{i} \prod_{j \neq i}\left(1-x_{j}\right)=C_{1}+\left(C_{2}-C_{1}\right)\left(x_{1}+x_{2}\right)+\left(C_{1}-2 C_{2}\right) x_{1} x_{2}
$$

Note that when fixing the sum $x_{1}+x_{2}$, we can change individual values of $x_{1}$ and $x_{2}$ within $[0,1]$ while still $\left\{x_{i}\right\}$ still being a valid solution. By the perturbing, only the term $\left(C_{1}-2 C_{2}\right) x_{1} x_{2}$ in the objective function might have value changed. Since $x_{1} \neq x_{2}$, we know that $C_{1}-2 C_{2}>0$ or making $x_{1}^{\prime}=x_{2}^{\prime}=\left(x_{1}+x_{2}\right) / 2$ gets no larger objective function value. When $C_{1}-2 C_{2}>0, x_{1}$ and $x_{2}$ should be "apart from" each other, thus one of $x_{1}$ and $x_{2}$ must touch their bound, i.e., 0 or 1 .

Now we use this claim and induction on $n$ to prove the lemma. The lemma trivially holds in the base case when $n=1$. When $n=k>1$, supposing the lemma holds for all $n<k$, we discuss the three cases proposed by Claim 21 (splitting the second case in the claim into two).

- When all $x_{i}$ 's are the same, we know that $x_{i}=S / n$ where $S=\sum_{i} x_{i}$. Then

$$
\sum_{i} x_{i} \prod_{j \neq i}\left(1-x_{j}\right)=S\left(1-\frac{S}{n}\right)^{n-1} \geqslant S e^{-S} \geqslant \frac{1-\varepsilon}{e^{2}}
$$

- When $\exists i: x_{i}=0$, with out loss of generality, suppose $x_{1}=0$. Then this reduces to the same problem with $(n-1)$ variables and the induction hypothesis gives us a $(1-\varepsilon) / e^{2}$ lower bound.
- When $\exists i: x_{i}=1$, again with out loss of generality, suppose $x_{1}=1$. Now the objective function becomes $\prod_{i \geqslant 2}\left(1-x_{i}\right)$ while $\sum_{i \geqslant 2} x_{i}$ is at most $\varepsilon$. It is easy to see the product is lower bounded by $(1-\varepsilon)$. (All but one of $x_{i}$ are 0 .)


## 6 Concluding remarks on finding almost-satisfying assignments for CSPs

In the world of "CSP dichotomy" (see [HN08] for a recent survey), the tractability of LIN-mod-2, 2SAT, and Horn-SAT is explained by the existence of non-trivial polymorphisms which combine many satisfying assignments to produce a new satisfying assignment. The Boolean functions which are polymorphisms for LIN-mod-2, 2-SAT, and Horn-SAT are xor (of odd size), majority, and minimum respectively. The existence of algorithms to find almost-satisfying assignments to 2-SAT and HornSAT can be attributed to the "noise stability" of the majority and minimum functions. The xor function of many variables, on the other hand, is highly sensitive to noise. This distinction seems to underly the difficulty of solving near-satisfiable instances of LIN-mod-2 and Håstad's tight hardness result for the problem.

For Boolean CSPs, we understand the complexity of finding almost-satisfying assignments for all the cases where deciding satisfiability is tractable: it is possible in polynomial time for 2-SAT and Horn-SAT, and NP-hard for LIN-mod-2. Further, under the UGC, the exact approximation threshold as a function of the gap $\varepsilon$ to perfect satisfiability is also pinned down for both 2-SAT and Horn-SAT. What about CSPs over larger domains? For any CSP $\Pi$ that can "express linear equations" (this notion is formalized in the CSP dichotomy literature, but we can work with the intuitive meaning for this discussion), Håstad's strong inapproximability result for near-satisfiable linear equations over abelian groups [Hås01] implies hardness of finding an almost satisfying assignment for ( $1-\varepsilon$ )satisfiable instances of $\Pi$. A recent breakthrough [BK09] established that every other tractable CSP (i.e., a polynomial time decidable CSP that cannot express linear equations) must be of "bounded width," which means that a natural local propagation algorithm correctly decides satisfiability of every instance of that CSP.

We end this paper with the appealing conjecture that every bounded width CSP admits a robust satisfiability algorithm that can find a $(1-g(\varepsilon))$-satisfying assignment given a $(1-\varepsilon)$-satisfiable instance for some function $g()$ such that $g(\varepsilon) \rightarrow 0$ as $\varepsilon \rightarrow 0$. (We should clarify that in this context we always treat the domain size $k$ of the CSP as fixed and let $\varepsilon \rightarrow 0$, so $g(\varepsilon)$ can have an arbitrary dependence on $k$. Note that Unique Games itself, which is a bounded width CSP, admits a robust satisfiability algorithm that satisfies a fraction $1-O(\sqrt{\varepsilon \log k})$ of constraints in a $(1-\varepsilon)$-satisfiable instance [CMM06].) By the preceding discussion, this conjecture would imply that bounded width characterizes the existence of robust satisfiability algorithms for CSPs.

## Acknowledgments

We thank Andrei Krokhin for a useful discussion on the relationship between the Unique Games conjecture and our "bounded width implies robust satisfiability" conjecture.

## References

[BK09] Libor Barto and Marcin Kozik. Constraint satisfaction problems of bounded width. In Proceedings of the 50th IEEE Symposium on Foundations of Computer Science, pages

595-603, October 2009. 23
[CMM06] Moses Charikar, Konstantin Makarychev, and Yury Makarychev. Near-optimal algorithms for unique games. In Proceedings of the 38th Annual ACM Symposium on Theory of Computing, pages 205-214, 2006. 23
[CMM09] Moses Charikar, Konstantin Makarychev, and Yury Makarychev. Near-optimal algorithms for maximum constraint satisfaction problems. ACM Transactions on Algorithms, 5(3), 2009. 3
[CST01] Pierluigi Crescenzi, Riccardo Silvestri, and Luca Trevisan. On weighted vs unweighted versions of combinatorial optimization problems. Information and Computation, 167(1):10-26, 2001. 2
[DFHS08] Erik D. Demaine, Uriel Feige, MohammadTaghi Hajiaghayi, and Mohammad R. Salavatipour. Combination can be hard: Approximability of the unique coverage problem. SIAM J. Comput., 38(4):1464-1483, 2008. 4, 5
[GT05] Venkatesan Guruswami and Luca Trevisan. The complexity of making unique choices: Approximating 1-in- k sat. In Proceedings of the 8th International Workshop on Approximation Algorithms for Combinatorial Optimization Problems (APPROX), pages 99-110, 2005. 1, 3, 4, 5, 8, 20
[Hås01] Johan Håstad. Some optimal inapproximability results. Journal of the ACM, 48(4):798859, 2001. 2, 23
[HN08] Pavol Hell and Jaroslav Nesetril. Colouring, constraint satisfaction, and complexity. Computer Science Review, 2(3):143-163, 2008. 23
[JKK09] Peter Jonsson, Andrei A. Krokhin, and Fredrik Kuivinen. Hard constraint satisfaction problems have hard gaps at location 1. Theor. Comput. Sci., 410(38-40):3856-3874, 2009. 2
[Kho02] Subhash Khot. On the power of unique 2-prover 1-round games. In Proceedings of the 34th Annual ACM Symposium on Theory of Computing, pages 767-775, 2002. 3
[KKMO07] Subhash Khot, Guy Kindler, Elchanan Mossel, and Ryan O'Donnell. Optimal inapproximability results for MAX-CUT and other 2-variable CSPs? SIAM J. Comput., 37(1):319-357, 2007. 3
[KSTW00] Sanjeev Khanna, Madhu Sudan, Luca Trevisan, and David P. Williamson. The approximability of constraint satisfaction problems. SIAM J. Comput., 30(6):1863-1920, 2000. 1, 2, 3, 4, 12
[Rag08] Prasad Raghavendra. Optimal algorithms and inapproximability results for every CSP? In Proceedings of the 40th ACM Symposium on Theory of Computing, pages 245-254, 2008. 1, 5, 6
[Sch78] T. J. Schaefer. The complexity of satisfiability problems. In Proceedings of the 10th ACM Symppsium on Theory of Computing, pages 216-226, 1978. 2
[Zwi98] Uri Zwick. Finding almost-satisfying assignments. In Proceedings of the 30th Annual ACM Symposium on the Theory of Computing, pages 551-560, 1998. 1, 2, 3, 7

## A Two positive semidefinite matrices

We now establish the positive semidefiniteness of the matrices encountered in Section 4.1.3.
Claim 22. Given $0<c \leqslant 0.2,0<p \leqslant \frac{1}{1+c} r_{c}, q=r_{c} p, \varepsilon=c / 1.5$, the following two matrices are positive semidefinite.

$$
\begin{aligned}
& A=\left[\begin{array}{ccccc}
1 & 1-p & 1-p & 1-r_{c} p & 1-r_{c} p \\
1-p & 1-p & 1-(1+c) p & 1-(1+c) p & 1-(1+c) p \\
1-p & 1-(1+c) p & 1-p & 1-(1+c) p & 1-(1+c) p \\
1-r_{c} p & 1-(1+c) p & 1-(1+c) p & 1-r_{c} p & 1-(1+c) p \\
1-r_{c} p & 1-(1+c) p & 1-(1+c) p & 1-(1+c) p & 1-r_{c} p
\end{array}\right], \\
& B=\left[\begin{array}{ccccc}
1 & 1-q & 1-q & 1-q & 1-q \\
1-q & 1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q \\
1-q & 1-(1+\varepsilon) q & 1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q \\
1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q & 1-q & 1-(1+1.5 \varepsilon) q \\
1-q & 1-(1+\varepsilon) q & 1-(1+\varepsilon) q & 1-(1+1.5 \varepsilon) p & 1-q
\end{array}\right] .
\end{aligned}
$$

Proof. Let $J$ be the all 1 matrix, $E_{1}$ be the matrix with 1 in entry $(1,1)$ as the only one non-zero entry. We also define $E_{i, j}, F_{i, j}$ and $G_{i, j}$ as matrices with only four non-zero entries located in the intersections of Column $i, j$ and Row $i, j$. The sub-matrices of $E_{i, j}, F_{i, j}$ and $G_{i, j}$ on Column $i, j$ and Row $i, j$ are defined as

$$
\left(\text { for } E_{i, j} \text { ) }\left[\begin{array}{ll}
1 & 1 \\
1 & 1
\end{array}\right],\left(\text { for } F_{i, j}\right)\left[\begin{array}{cc}
2 & 1 \\
1 & 0.5
\end{array}\right] \text { and (for } G_{i, j}\right)\left[\begin{array}{cc}
1 & -1 \\
-1 & 1
\end{array}\right] \text {. }
$$

Clearly, all of $J, E_{1}, E_{i, j}, F_{i, j}$ and $G_{i, j}$ are positive semidefinite matrices.
Then we can write $A$ as

$$
\begin{aligned}
A & =(1-(1+c) p) J+c p\left(E_{1,2}+E_{1,3}\right)+\left(1+c-r_{c}\right) p\left(E_{1,4}+E_{1,5}\right)+\left(2 r_{c}-1-3 c\right) p E_{1} \\
& =(1-(1+c) p) J+c p\left(E_{1,2}+E_{1,3}\right)+\frac{(1+c) c}{1.5+c} \cdot p\left(E_{1,4}+E_{1,5}\right)+\frac{1.5-2.5 c-3 c^{2}}{1.5+c} \cdot p E_{1},
\end{aligned}
$$

Note that all the coefficient before matrices are non-negative within the range of $c$. Since $A$ can be written as the sum of several positive semidefinite matrices, A is positive semidefinite.

For matrix $B$, note that

$$
B=(1-(1+\varepsilon) q) J+\varepsilon q\left(E_{1,2}+E_{1,3}+F_{1,4}+F_{1,5}\right)+0.5 \varepsilon q G_{4,5}+(1-5 \varepsilon) E_{1},
$$

Clearly, as long as $5 \varepsilon=5 c / 1.5<1, B$ can be expressed as sum of positive semidefinite matrices, and hence $B$ is positive semidefinite.


[^0]:    *Supported in part by a Packard Fellowship, NSF CCF 0953155, and US-Israel BSF grant 2008293. Email: guruswami@cmu.edu, yuanzhou@cs.cmu.edu.

[^1]:    ${ }^{1}$ The dual variant dual-Horn-SAT which is an instance of SAT where each clause has at most one negated literal is also polynomial time solvable.
    ${ }^{2}$ Throughout the paper, constraints could have weights, and by a "fraction $\alpha$ of constraints" we mean any subset of constraints whose total weight is a fraction $\alpha$ of the sum of the weights of all constraints. For CSPs with no unary constraints, the approximability of the weighted and unweighted versions are known to be the same [CST01].

