# MEASURE-DRIVEN ALGORITHM DESIGN AND ANALYSIS: 

## A NEW APPROACH FOR SOLVING NP-HARD PROBLEMS

A Dissertation<br>by<br>YANG LIU

Submitted to the Office of Graduate Studies of Texas A\&M University in partial fulfillment of the requirements for the degree of DOCTOR OF PHILOSOPHY

August 2009

Major Subject: Computer Science

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ABSTRACT<br>Measure-Driven Algorithm Design and Analysis:<br>A New Approach for Solving NP-hard Problems. (August 2009)<br>Yang Liu, B.S., Zhejiang University;<br>M.S., Rose-Hulman Institute of Technology<br>Chair of Advisory Committee: Dr. Jianer Chen

NP-hard problems have numerous applications in various fields such as networks, computer systems, circuit design, etc. However, no efficient algorithms have been found for NP-hard problems. It has been commonly believed that no efficient algorithms for NP-hard problems exist, i.e., that $\mathrm{P} \neq \mathrm{NP}$. Recently, it has been observed that there are parameters much smaller than input sizes in many instances of NP-hard problems in the real world. In the last twenty years, researchers have been interested in developing efficient algorithms, i.e., fixed-parameter tractable algorithms, for those instances with small parameters. Fixed-parameter tractable algorithms can practically find exact solutions to problem instances with small parameters, though those problems are considered intractable in traditional computational theory.

In this dissertation, we propose a new approach of algorithm design and analysis: discovering better measures for problems. In particular we use two measures instead of the traditional single measure - input size to design algorithms and analyze their time complexity. For several classical NP-hard problems, we present improved algorithms designed and analyzed with this new approach,

First we show that the new approach is extremely powerful for designing fixedparameter tractable algorithms by presenting improved fixed-parameter tractable algorithms for the 3D-matching and 3D-PACKING problems, the multiway cut
problem, the FEEDBACK VERTEX SET problems on both directed and undirected graph and the max-LEAF problems on both directed and undirected graphs. Most of our algorithms are practical for problem instances with small parameters.

Moreover, we show that this new approach is also good for designing exact algorithms (with no parameters) for NP-hard problems by presenting an improved exact algorithm for the well-known SATISFIABILITY problem.

Our results demonstrate the power of this new approach to algorithm design and analysis for NP-hard problems. In the end, we discuss possible future directions on this new approach and other approaches to algorithm design and analysis.

To my family

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## TABLE OF CONTENTS

## CHAPTER <br> Page

I INTRODUCTION ..... 1
A. Why Study Exact Algorithms ..... 1
B. Why Study Fpt-algorithms ..... 3
C. Branch-and-Search Tree ..... 4
D. Outline of This Dissertation ..... 8
II 3D-MATCHING AND 3-SET PACKING ..... 10
A. Introduction ..... 11
B. Preliminaries and Reformulations ..... 14
C. Improved Packing Algorithms ..... 16
D. Matching Algorithms Further Improved ..... 22
E. Final Remarks ..... 27
III MULTIWAY CUT ..... 29
A. Introduction ..... 29
B. Minimum V-cuts between Two Terminal Sets ..... 32
C. The Main Algorithm ..... 36
D. Final Remarks ..... 47
IV UNDIRECTED FEEDBACK VERTEX SET ..... 48
A. Introduction ..... 48
B. Feedback Vertex Set in Unweighted Graphs ..... 52
C. Feedback Vertex Set in Weighted Graphs ..... 60
V DIRECTED FEEDBACK VERTEX SET ..... 76
A. Introduction ..... 76
B. Preliminaries ..... 80
C. Solving the Skew separator Problem ..... 84
D. Solving the Dag-bipartition fVs Problem ..... 95
E. Solving the DFvs Problem ..... 99
F. Final Remarks ..... 103
VI MAX-LEAF ..... 108
A. Introduction ..... 108
B. Preliminaries ..... 110
C. Extending an Out-tree ..... 113

1. Properties for Extending an Out-tree ..... 114
2. Extending an Out-tree ..... 117
D. The Main Algorithm and Complexity Analysis ..... 124
E. Final Remarks ..... 130
VII SATISFIABILITY ..... 131
A. Introduction ..... 131
B. Preliminaries ..... 133
C. Reduction Rules ..... 136
D. Main Algorithm ..... 141
E. Analysis of the Main Algorithm ..... 142
3. Analysis for Degree-4 Formulas ..... 150
a. $d_{0}=1$ ..... 151
b. $d_{0}=2$ ..... 151
4. Analysis for Degree-5 Formulas ..... 156
a. $d_{0}=1$ ..... 156
b. $d_{0}=2$ ..... 157
5. Analysis for Formulas of Degree Larger Than 5 ..... 163
6. Branching Vector for the Main Algorithm ..... 164
F. Final Remarks ..... 165
VIII SUMMARY AND FUTURE RESEARCH ..... 168
A. Dissertation Summary ..... 168
B. Future Work ..... 171
7. Further Study of the New Approach ..... 171
8. Randomized and Algebraic Algorithms ..... 172
9. Kernelization ..... 173
REFERENCES ..... 174
VITA ..... 187

## LIST OF TABLES

TABLE Page
I Comparison of algorithms for 3-D matching ..... 14
II History of parameterized algorithms for the UNWEIGHTED FEED-BACK VERTEX SET problem . . . . . . . . . . . . . . . . . . . . . . . 50
III History of exact algorithms for the SATISFIABILITY problem ..... 132

## LIST OF FIGURES

FIGURE Page
1 Dynamic programming for 3-SET PACKING ..... 20
2
Dynamic programming for 3-D matching ..... 23
3 Decomposition of separators ..... 39
4
An algorithm for the PARAMETERIZED NODE MULTIWAY CUT problem ..... 43
5
Algorithm for the UNWEIGHTED FEEDBACK VERTEX SET problem. ..... 54
6
Algorithm for the WEIGHTED FEEDBACK VERTEX SET problem ..... 75
7 Sets in the proof of Theorem C.4. ..... 87
8 An algorithm for the SKEW SEPARATOR problem. ..... 106
9 An algorithm for the DAG-BIPARTITION FVS problem. ..... 107
10
Algorithm for the SPAN $k$-OUT-TREE problem ..... 121
11
Algorithm for the EXTENDING MAX-LEAF problem ..... 125
12 The reduction algorithm ..... 166Algorithm for the satisfiability problem167

## CHAPTER I

## INTRODUCTION

There are abundant applications of NP-hard problems. These applications have attracted intensive studies of NP-hard problems. There are many methods to tackle NP-hard problems: exact algorithms, approximation algorithms, probabilistic algorithms, heuristics, and parameterized algorithms. In this dissertation, we are interested in design and analysis of exact algorithms and parameterized algorithms for NP-hard problems. In particular, we are interested in fpt-algorithms which are a special class of parameterized algorithms. The rest of this chapter is organized as follows: in section A, we discuss why we study exact algorithms. In section B, we discuss why we study fixed-parameter tractable algorithms. In section C, we introduce the branch-and-search tree technique which will be widely applied later in this dissertation. In section D, we give an outline of this dissertation.

## A. Why Study Exact Algorithms

First, exact algorithms can find optimal solutions, which minimize the cost or maximize the benefits in applications. This is the major reason why we prefer exact algorithms to approximation algorithms and heuristics. Moreover, only exact algorithms can find meaningful and correct answers for decision problems. Solutions from approximation algorithms make no sense at all, and heuristics can be wrong for some inputs. For example, the Satisfiability problem should be answered with either "satisfiable" or "not satisfiable". Approximation algorithms can not help at all. Therefore, we need to study exact algorithms for NP-hard problems.

Second, exact algorithms also can be used as a subroutine by other heuristics. For example, optimal solutions to small instances are found and stored in a table. Then heuristics break a large instance into small instances, and use the solutions in the table to find a solution (not necessary optimal) for the large input instance. This technique helps to find better solutions [26] for the RECTILINEAR STEINER MINIMAL tree problem. So studying exact algorithms can help to develop heuristics with better solutions.

Moreover, faster exact algorithms for NP-hard problem may increase significantly the size of instances which can be handled, since those algorithms are expected to be of exponential running time according to the widely accepted belief that $\mathrm{P} \neq \mathrm{NP}$. For example, if an exact algorithm of time $O\left(2^{n / 2}\right)$ is developed and the previous best algorithm is of time $O\left(2^{n}\right)$, then the size of instances which can be handled by the new exact algorithm is twice that of instances which can be handled by the previous best algorithm. On the other hand, the speedup of processors can not increase linearly the size of instances which can be handled as exact algorithms. With the amazing progress of technology, the processors are almost 9, 000 times faster in 2004 than the processors in 1978 [60]. But the size of instances which can be handled by the 9,000 times faster processors are only $\log _{2} 9,000$ larger than that of instances which can be handled by the old processors, for an $O\left(2^{n}\right)$ algorithm. The increase in instance size is less than 13. Therefore, it is still desirable to study exact algorithms, despite the incredible speedup of processors.

Current progress in exact algorithms for NP-hard problems appeals to further studies of exact algorithms. A notable example is the Rectilinear Steiner Trees problem. There is an exact algorithm which finds an optimal Rectilinear Steiner Tree in 38 CPU hours for instances of size up to 1000 [100]. Another exciting example is the traveling salesman problem. This problem now can be solved for inputs of
size 13,509 , though 48 workstations had run for around 10 years [4].

## B. Why Study Fpt-algorithms

Though studying exact algorithms is necessary for NP-hard problems, it is notoriously difficult to develop fast algorithms such that we can solve instances with considerable sizes. However, it has been observed that there are parameters in many instances of NP-hard problems, which are small compared to their input sizes. Given such an instance, let $n$ be its input size and $k$ be its parameter. Then $k$ may be much smaller than $n$. Moreover, the running time of algorithms for that instance is a function of both $n$ and $k$. Parameterized algorithms are those algorithms whose running time is a function of $n$ and $k$. In last decade, many researchers have studied parameterized algorithms extensively.

Not all parameterized algorithms are practical. For example, an algorithm of running time $O\left(n^{k}\right)$ is not practical even for $k=10$ when $n$ is moderately large. In particular, people are interested in efficient parameterized algorithms whose running time are $O\left(f(k) n^{c}\right)$ where (1) $c$ is a constant independent of input size $n$ and the parameter $k$, and (2) $f(k)$ is independent of $n$. Those efficient parameterized algorithms are fixed parameter tractable (fpt-) algorithms. A problem is fixed parameter tractable if there exists an fpt-algorithm for the problem.

For some instances of NP-hard problems in real world, there are natural parameters that are much smaller than the sizes of inputs. We give two examples.

The first example is the problem Type Checking in ML [38]. One of the tasks of a compiler is to check the compatibility of type declarations. This problem is exponential-time complete [59]. There is an fpt-algorithm of running time $O\left(2^{k} n\right)$ where $n$ is the size of the program and $k$ is the maximum nesting depth of the type
declarations [77]. Normally $k$ is not greater than 6 . So this algorithm works well in practice.

Another example is a recent fpt-algorithm for the Individual Haplotyping problem [103]. The fpt-algorithm is of time complexity $O\left(n k_{2} 2^{k_{2}}+m \log m+m k_{1}\right)$ where $m$ is the number of fragments (the number of $0 / 1 / 2$ sequences), $n$ is the number of SNP sites (the length of each $0 / 1 / 2$ sequence), $k_{1}$ is the maximum number of SNP sites that a fragment covers (the maximum number of $0 / 1$ 's in a sequences), and $k_{2}$ is the maximum number of the fragments covering an SNP site (the number of sequences which have $0 / 1$ in a particular position in the sequence). The parameter $k_{1}$ is bounded by $n$. But normally $k_{1}$ is smaller than 10 . The parameter $k_{2}$ is usually not more than 19 [63]. Thus this algorithm is practical. Moreover, it is more accurate in haplotype reconstruction than other known algorithms.

From these examples, we can see that fpt-algorithms are practical for instances with small parameters. So we need to study fpt-algorithms for NP-hard problems. We will investigate a new approach for designing fpt-algorithms in this dissertation.

## C. Branch-and-Search Tree

There are many algorithmic approaches to design exact algorithms (without parameters) and fpt-algorithms. Algorithmic approaches for exact algorithms (without parameters) include dynamic programming, pruning the search tree, preprocessing the data and local search [101], measure and conquer [49, 50]. Algorithmic approaches for fpt-algorithms include kernelization [39], greedy localization [19, 42], iterative compression [90, 32], coloring coding [3], divide and conquer [25], divide-and-color [70], and branch-and-search tree [20, 73]. Among these approaches, branch-and-search tree is the most widely used approach, and can be applied to design both exact algorithms
(without parameters) and fpt-algorithms.
Given an instance $I$ of a problem, the branch-and-search tree approach can be applied to design a recursive algorithm $\mathcal{A}$ to solve $I$. If the instance $I$ satisfies some conditions, $\mathcal{A}$ solves the problem for the instance $I$ directly. Otherwise, $\mathcal{A}$ reduces the instance $I$ into some smaller instances $I_{1}, \cdots, I_{p}$, calls $\mathcal{A}$ recursively on each of $I_{1}, \cdots, I_{p}$ to find solutions to $I_{1}, \cdots, I_{p}$, and then finds solutions to $I$ from solutions to $I_{1}, \cdots, I_{p}$.

To study the time complexity of such an algorithm, we consider the algorithm $\mathcal{A}$ as a tree $\mathcal{T}$. Each node represents an instance. Let instance $(N)$ be the instance represented by node $N$. If $N$ is the root of $\mathcal{T}$, then instance $(N)$ is the original instance $I$. If $N$ is a leaf of $\mathcal{T}$, then $\operatorname{instance}(N)$ is an instance that can be solved directly without recursive call of $\mathcal{A}$. A node $N$ and its children $N_{1}, \cdots, N_{q}$ represent a call of $\mathcal{A}$ on instance $(N)$ during which instance $(N)$ is reduced to instance $\left(N_{1}\right), \cdots$, instance $\left(N_{q}\right)$.

Now we can distribute the running time of $\mathcal{A}$ on instance $I$ in the tree $\mathcal{T}$ as follows: the time for a leaf $L$ is the time to solve directly instance $(L)$. The time for an internal node $N$ with children $N_{1}, \cdots, N_{q}$ is the time to reduce instance $(N)$ to instance $\left(N_{1}\right), \cdots$, instance $\left(N_{q}\right)$ plus the time to find solution to instance $(N)$ from solutions to instance $\left(N_{1}\right), \cdots$, instance $\left(N_{q}\right)$. It is clear that the running time of $\mathcal{A}$ on $I$ is the summation of the time for nodes in the tree $\mathcal{T}$.

Normally it is required that the time for each node be polynomial in the input size $n$ of $I$, and the height of $\mathcal{T}$ be also polynomial in $n$. Let $T_{\text {path }}$ be the summation of the time for all nodes in a path from the root to a leaf in $\mathcal{T}$. Then $T_{\text {path }}$ is polynomial in $n$. Let $T_{\text {max }}$ be the maximum among all possible $T_{\text {path }}$. The running time of $\mathcal{A}$ is then bounded by the product of $T_{\max }$ and the number of leaves in $\mathcal{T}$.

In general, it is easier to give an upper bound of the number of leaves in $\mathcal{T}$ than
to calculate the precise number of leaves in $\mathcal{T}$. To bound the number of leaves in $\mathcal{T}$, we assign a measure to every node in $\mathcal{T}$. Let $m$ be the measure of the root in $\mathcal{T}$, and $m_{1}, \cdots, m_{q}$ be the measures corresponding to the children of the root in $\mathcal{T}$. A characteristic polynomial for $m, m_{1}, \cdots, c_{q}$ is $\sum x^{d_{i}}=1$ where $d_{i}=m_{i}-m$ for $i=1, \cdots, q$. We also say that the root has a characteristic polynomial $\sum x^{d_{i}}=1$. There is only one positive root for any characteristic polynomial. Let $\alpha$ be the root of a characteristic polynomial. It can be proved that $f(m) \leq \alpha^{m}$ when all internal nodes have the same characteristic polynomial. If there are multiple characteristic polynomials $\mathcal{P}_{1}, \cdots, \mathcal{P}_{q}$, let $\alpha_{i}$ be the positive root of the characteristic polynomial $\mathcal{P}_{i}$, and $\alpha=\left(\max \left\{\alpha_{1}, \cdots, \alpha_{q}\right\}\right)$. Then we have $f(m) \leq \alpha^{m}$. More explanations and proofs of this approach can be found in [73].

From the arguments above, the total running time of algorithm $\mathcal{A}$ is bounded by $O\left(\alpha^{m} \operatorname{poly}(n)\right)$. For simplicity, we use notation $O^{*}\left(\alpha^{m}\right)^{1}$ to ignore the part of poly $(n)$ in $O\left(\alpha^{m}\right.$ poly $\left.(n)\right)$.

Now we give examples to illustrate the approach of branch-and-search tree. Consider the minimum vertex cover problem, which is NP-hard. Given a graph $G$, a vertex cover $C$ of $G$ is a subset of vertices such that every edge has at least one endpoint in $C$. A minimum vertex cover is a vertex cover of the minimum number of vertices among all vertex covers. The minimum vertex cover problem is to find a minimum vertex cover for the input graph $G$. We design an exact algorithm for the MINIMUM VERTEX COVER problem with the branch-and-search tree approach.

Normally, we choose the number of vertices $n$ to be the measure. Let $C$ be a minimum vertex cover. If there are no edges in the input graph $G$, then return $\phi$ as a minimum vertex cover. Otherwise, pick an edge $x y$ and consider two cases: either

[^0]$x$ or $y$ is in the minimum vertex cover $C$. If $x$ is in $C$, let $G_{x}$ be the remaining graph after deleting $x$ and all incident edges of $x$. Then we only need to find a minimum vertex cover $C-x$ for $G_{x}$. Note that $\left|V\left[G_{x}\right]\right|=|V[G]|-1$, i.e., the measure decreases by 1 . If $x$ is not in $C$, then $y$ must be in $C$. Let $G_{y}$ be the graph after deleting $y$. Then we only need to find a minimum vertex cover $C-y$ for $G_{y}$. The measure of $G_{y}$ also decreases by 1. The characteristic polynomial for these measures is $x^{-1}+x^{-1}=1$. Solve it and find its positive root $\alpha=2$. So $f(n) \leq 2^{n}$. That is, the minimum VERTEX COVER problem can be solved in time $O^{*}\left(2^{n}\right)$.

Next we show how to design fpt-algorithms with branch-and-search tree. The parameter for the VERTEX COVER problem is the size of a vertex cover to find. Let $k$ be the size of a vertex cover to find. Formally, the $k$-VERTEX COVER problem is to: either find a vertex cover of $k$ vertices if such a vertex cover exists, or report 'NO' if none exists. The same algorithm above can be applied to solve the $k$-VERTEX COVER problem. However, we should redefine $C$ to be the vertex cover of $k$ vertices to search, and calculate the measure changes accordingly: (1) for $G_{x}$, the measure decreases by 1 since we only need to find $C-x$ in $G_{x}$, and (2) for $G_{y}$, the measure decreases by 1 since we only need to find $C-y$ in $G_{y}$. The characteristic polynomial for these measures is still $x^{-1}+x^{-1}=1$, whose positive root $\alpha=2$. Therefore, the $k$-vertex cover problem can be solved in time $O^{*}\left(2^{k}\right)$.

It is difficult to design better algorithm of time $O^{*}\left(c^{n}\right)$ or $O^{*}\left(c^{k}\right)$ with small $c$. The best algorithm for the MINIMUM VERTEX COVER problem has running time of $O^{*}\left(1.1893^{n}\right)$ [91]. The best algorithm for the $k$-VERTEX COVER problem has running time of $O^{*}\left(1.2738^{k}\right)$ [21]. But these algorithms are much more complicated than the simple algorithm above.

Before our works [78, 23, 18, 24, 22], researchers had taken only one parameter when using the approach of branch-and-search tree. This common usage has made
it difficult to design fast algorithms. For example, let us consider the $k$ feedback vertex set problem. This problem either finds a set $S$ of $k$ vertices such that every cycle contains at least one vertex in $S$, or reports no such set exists. To apply the traditional approach of branch search tree, we naturally pick a cycle $C$ and have $p$ choices to put a vertex in $S$, where $p$ is the length of $C$. Then we have a characteristic polynomial $f(k)=p f(k-1)$, whose positive root is $p$. Thus any algorithm by this way has time complexity of $O^{*}\left(p^{k}\right)$. It is hard to bound the length of cycles by some constant in graphs. Thus it is difficult to design algorithms of time $O^{*}\left(c^{k}\right)$, where $c$ is a constant not related to $k$, by this traditional approach. In this dissertation, we show that using two measures, instead of single measure, is much more powerful than the traditional approach of branch-and-search tree because more properties can be applied to design algorithms. With this new approach, we present an algorithm of time $O^{*}\left(5^{k}\right)$ for the $k$ FEEDBACK VERTEX SET problem on undirected graphs in Chapter IV.
D. Outline of This Dissertation

In Chapters II to VI, we apply our new approach to design better fpt-algorithms for several problems: 3D-Matching and 3D-Packing problems, multiway-Cut problem, feedback vertex set problems on undirected graphs and directed graphs, and max leaf problem. Detailed analyses are given. The results in these chapters illustrate the power of our new approach for fpt-algorithm design and analysis.

Then in Chapter VII, we illustrate how to design an improved algorithm for the well-known SATISFIABILITY problem using two measures. The result in this chapter demonstrates that our new approach also works well for exact algorithm design and analysis.

In Chapter VIII, we give a summary of our work and directions for future research.

## CHAPTER II

## 3D-MATCHING AND 3-SET PACKING

In this chapter, we give improved randomized and deterministic fpt-algorithms for the 3-D MATCHING and 3-SET PACKING problems. Our randomized algorithm for the 3-D MATCHING problem has running time of $O^{*}\left(2.32^{3 k}\right)$, which improves the previous best randomized algorithm of running time $O^{*}\left(2.52^{3 k}\right)$. Our deterministic algorithm for the 3 -D matching is of running time $O^{*}\left(2.77^{3 k}\right)$, which improves the previous best deterministic algorithms of running time $O^{*}\left(12.8^{3 k}\right)$. Our deterministic algorithm for the 3 -SET PACKING problems are of running time $O^{*}\left(4.61^{3 k}\right)$, which improves the previous best deterministic algorithms of running time $O^{*}\left(12.8^{3 k}\right)$.

The previous algorithms for the 3-D matching and 3-SET PACKING problems focus on the parameter $k$ as a measure, which makes it difficult to design faster algorithms. We still study these problems with the greedy localization approach [19]: start from a matching (packing) of size $p$ and to find a matching (packing) of size $k$ where $k>p$. However, instead of focusing on the parameter $k$ only, we consider two measures: the number of colors used for coloring and the number of elements used in dynamic programming. With this new approach, we discover that only $p+3$ colors are needed, and the number of elements used in dynamic programming is $4 p+3$, if we try to find a matching (packing) of size $p+1$. This reduces the number of colors from $2 k$ to $p+3$, and the number of elements for dynamic programming from $4 k+2$ to $4 p+3$, thus improving the time complexity significantly to $O^{*}\left(4.61^{3 k}\right)$ for the 3 -D matching and 3 -set packing problems. Moreover, we can do better for the 3-D matching PROBLEM with this approach. It is observed that if we keep only two columns of the known matching, then the number of elements for dynamic programming is only $8 p / 3+2$, while the number of colors is only $2 p / 3+2$. This results in a randomized
algorithm of running time $O^{*}\left(2.32^{3 k}\right)$ and a deterministic algorithm of running time $O^{*}\left(2.77^{3 k}\right)$ for the 3-D matching problem.

## A. Introduction

Matching and packing problems have formed an important class of NP-hard problems. In particular, the 3-D matching problem is one of the six "basic" NP-complete problems in terms of Garey and Johnson [53], and the 3-SET PACKING problem is a natural extension of the 3-D matching problem. There has been a remarkable line of research in the study of parameterized algorithms for 3-D mATCHING and 3-SET PACKING problems.

Downey and Fellows [37] proved that the 3-D matching problem is fixedparameter tractable and gave an algorithm of time $O^{*}\left((3 k)!(3 k)^{9 k+1}\right)$ Chen et al. [19] improved the time complexity for 3-D matching to $O^{*}\left((5.7 k)^{k}\right)$, and Jia, Zhang, and Chen [19] improved the time complexity for 3 -SET PACKING to $O^{*}\left((5.7 k)^{k}\right)$.

More progress has been made recently. For the 3-SET Packing problem, Koutis [71] developed a randomized algorithm of time $O^{*}\left(10.88^{3 k}\right)$ and a deterministic algorithm of time $O^{*}\left(2^{O(k)}\right)$. Koutis [71] did not give the exact constant factor in the exponent of the time complexity $O^{*}\left(2^{O(k)}\right)$ for his deterministic algorithm. He used the perfect hashing families proposed by Schmidt and Siegel [92], in which the number of hashing functions to hash $n$ elements into $3 k$ colors is larger than $2^{\log \log n+12 k}$. It can be derived that his deterministic algorithm has time complexity of at least $O^{*}\left(32000^{3 k}\right)$. These algorithms can be applied to the 3-D matching problem without any changes. Fellows et al. [44] studied the complexity of matching and packing problems. They first showed that the 3-D matching problem has a kernel of size $O\left(k^{3}\right)$, and then presented an algorithm of time $O^{*}\left(2^{O(k)}\right)$ for the problem, where
the term $O(k)$ was also not specified in detail. It can be deduced that the running time of the algorithm given in [44] for 3-D matching is at least $O^{*}\left(12.67^{3 k} T(k)\right)$, where $T(k)$ is the running time of a dynamic programming algorithm that, on a set of triples whose symbols are colored with $13 k$ colors, searches for a matching of $k$ triples in which all symbols are colored with distinct colors $\left(T(k)\right.$ is at least $O^{*}\left(10.4^{3 k}\right)$ using currently known techniques). The chapter also discussed how these techniques are applied to solve 3-SET PACKING and various graph packing problems.

There are at least two very recent works that give further improved algorithms for 3-D matching and 3-SEt packing problems. Chen et al. [25] proposed a new technique based on divide-and-conquer that leads to randomized algorithms of time $O^{*}\left(2.52^{3 k}\right)$ for 3-D matching and 3-SET Packing problems. Moreover, they proposed a color-coding scheme of $O^{*}\left(6.1^{k}\right) k$-colorings which, when combined with standard dynamic programming techniques, gives deterministic algorithms of running time $O^{*}\left(12.8^{3 k}\right)$ for 3-D matching and 3-SET PaCKing problems. We point out that using this new color-coding scheme, the time complexity of the algorithms by Coitus [71] for 3-D matching and 3-SET PaCKing can be improved to $O^{*}\left(25.6^{3 k}\right)$, and the time complexity of the algorithms by Fellows et al. [44] can be improved to $O^{*}\left(13.78^{3 k}\right)$. In a work performed independently of that in [25], Kneis et al. [70] developed a divide-and-conquer method that leads to randomized algorithms for 3-D MATCHING and 3 -SET PACKING problems with time complexity similar to that in [25]. Moreover, a different de-randomization method was proposed in [70] based on the work of [84], which leads to deterministic algorithms of running time $O^{*}\left(16^{3 k}\right)$ for 3-D MATCHING and 3-SET PACKING problems.

The known parameterized algorithms for 3-D MATCHING and 3-SET PACKING have used either the technique of greedy localization [19, 32, 64], the technique of colorcoding [3] plus dynamic programming [25, 44, 71], or the divide-and-conquer method
$[25,70]$. In this chapter, we show how a combination of these techniques and new techniques will yield further improved algorithms for these problems. We start with the 3 -SEt Packing problem. In difference from the approach used in $[19,64]$ that constructs a packing of $k 3$-sets directly from a maximal packing, we concentrate on the construction of a packing of $k+13$-sets based on a packing of $k 3$-sets. This slight modification enables us to derive a property for packing that is much stronger than the one given in [19]. Moreover, instead of coloring all elements in an instance of 3-SET PACKING, we color only part of the elements and use either ordering or pre-selected elements to reduce the complexity of the coloring stage in the algorithms. Using these new techniques, we are able to develop a parameterized algorithm of running time $O^{*}\left(4.61^{3 k}\right)$ for the 3-SET PACKING problem, significantly improving the previous best algorithm of running time $O^{*}\left(12.8^{3 k}\right)$ for the problem [25]. For the 3-D matching problem, we further show that the complexity of the dynamic programming stage in the algorithms, which seems to have been largely neglected in the previous research, can also be improved using a pre-ordering technique. Combining this new technique and those developed for 3-SET PACKING, we achieve further improved algorithms for the 3-D matching problem. More specifically, our new randomized algorithm for 3-D MATCHING runs in time $O^{*}\left(2.32^{3 k}\right)$, and our new deterministic algorithm for 3-D matching runs in time $O^{*}\left(2.77^{3 k}\right)$, both significantly improving the previous best algorithms for the problem.

We would like to point out that all previous parameterized algorithms for 3D matching and 3-SET PACKing have the same time complexity for both problems, although it is obvious that 3-SET PACKING is a nontrivial generalization of 3-D matching. The results in the current chapter seem to give faster algorithms for 3-D matching than for 3 -SET Packing. We also mention that the difference in complexity between our deterministic algorithm (i.e., $O^{*}\left(2.77^{3 k}\right)$ ) and our random-

Table I. Comparison of algorithms for 3-D matching

| References | Randomized algorithm | Deterministic algorithm |
| :--- | :--- | :--- |
| Downey and Fellows [37] |  | $O^{*}\left((3 k)!(3 k)^{9 k+1}\right)$ |
| Chen et al. [19] | $O^{*}\left((5.7 k)^{k}\right)$ |  |
| Koutis [71] |  | $>O^{*}\left(32000^{3 k}\right)$ |
| Fellows et al. [44] ${ }^{*}$ | $O^{*}\left(2.52^{3 k}\right)$ | $>O^{*}\left(12.67^{3 k} T(k)\right)$ |
| Kneis et al. [70] | $O^{*}\left(2.52^{3 k}\right)$ | $O^{*}\left(16^{3 k}\right)$ |
| Chen et al. [25] | $O^{*}\left(2.32^{3 k}\right)$ | $O^{*}\left(12.8^{3 k}\right)$ |
| Our new result |  | $O^{*}\left(2.77^{3 k}\right)$ |

${ }^{*} T(k)$ is the running time of a dynamic programming process that, on a set of triples whose symbols are colored with $13 k$ colors, searches for a matching of $k$ triples in which all symbols are colored with distinct colors. Based on currently known techniques, $T(k)$ is at least $O^{*}\left(10.4^{3 k}\right)$.
ized algorithm (i.e., $O^{*}\left(2.32^{3 k}\right)$, which is also currently the best upper bound) for 3-D matching has been significantly narrowed down, which is remarkable considering the fact that in the previous research on the problem, the difference between these two kinds of algorithms is in general very significant. Table I gives a specific comparison of our new algorithms and the previous algorithms for the 3-D matching problem.

## B. Preliminaries and Reformulations

Let $X, Y$, and $Z$ be three pairwise disjoint symbol sets, and let $U=X \times Y \times Z$ be the product set of $X, Y$, and $Z$. Each element $t=(x, y, z)$ in $U$, where $x \in X, y \in Y$, and $z \in Z$, is called a triple. For a triple $t=(x, y, z)$ in $U$, denote by $\operatorname{Val}(t)$ the set
$\{x, y, z\}$, and let $\operatorname{Val}^{1}(t)=\{x\}, \operatorname{Val}^{2}(t)=\{y\}, \operatorname{Val}^{3}(t)=\{z\}$. We say that a triple $t_{1}$ conflicts with another triple $t_{2}$ if $t_{1} \neq t_{2}$ and $\operatorname{Val}\left(t_{1}\right) \cap \operatorname{Val}\left(t_{2}\right) \neq \emptyset$. Let $S$ be a set of triples in $U$. Denote $\operatorname{Val}(S)=\bigcup_{t \in S} \operatorname{Val}(t)$, and $\operatorname{Val}^{i}(S)=\bigcup_{t \in S} \operatorname{Val}^{i}(t)$ for $i=1,2,3$. A matching in $S$ is a subset $M$ of triples in $S$ such that no two triples in $M$ conflict with each other. A matching $M$ in $S$ is a $k$-matching if $M$ contains exactly $k$ triples.

Packing problems are a generalization of matching problems. We say that a set $\rho_{1}$ conflicts with another set $\rho_{2}$ if $\rho_{1} \neq \rho_{2}$ and $\rho_{1} \cap \rho_{2} \neq \emptyset$. Let $S$ be a collection of sets. Denote $\operatorname{Val}(S)=\cup_{\rho \in S} \rho$. A packing in $S$ is a sub-collection $P$ of $S$ such that no two sets in $P$ conflict with each other. A packing $P$ in $S$ is a $k$-packing if $P$ contains exactly $k$ sets.

The main problems we study in this chapter are formally defined as follows.
(PARAMETERIZED) 3-D MATCHING:
Given a pair $(S, k)$, where $S$ is a set of $n$ triples, and $k$ is an integer, either construct a $k$-matching in $S$ or report that no such matching exists.
(PARAMETERIZED) 3-SET PACKING:
Given a pair $(S, k)$, where $S$ is a collection of $n$ sets, each containing at most three elements, and $k$ is an integer, either construct a $k$-packing in $S$ or report that no such packing exists.

A set is a 3-set if it contains exactly three elements. For an instance $(S, k)$ of 3 -SET PACKING, we can assume, without loss of generality, that all sets in $S$ are 3 -sets (otherwise, we can add new elements, i.e., elements not in $S$, to convert each set with fewer than three elements to a 3 -set). Instead of working on the above problems, we will concentrate on the following related problems.

3-D MATCHING AUGMENTATION:
Given a pair $\left(S, M_{k}\right)$, where $S$ is a set of triples, and $M_{k}$ is a $k$-matching in $S$, either construct a $(k+1)$-matching $M_{k+1}$ in $S$, or report that no such matching exists.

## 3-SET PACKING AUGMENTATION:

Given a pair $\left(S, P_{k}\right)$, where $S$ is a collection of $n 3$-sets, and $P_{k}$ is a $k$ packing in $S$, either construct a $(k+1)$-packing $P_{k+1}$ in $S$, or report that no such packing exists.

Lemma B. 1 For any constant $c>1$, the 3-D matching augmentation problem can be solved in time $O^{*}\left(c^{k}\right)$ if and only if the 3-D MATChing problem can be solved in time $O^{*}\left(c^{k}\right)$. Similarly, the 3 -SET PACKING AUGMENTATION problem can be solved in time $O^{*}\left(c^{k}\right)$ if and only if the 3 -SET PACKING problem can be solved in time $O^{*}\left(c^{k}\right)$.

According to Lemma C.1, we only need to concentrate on the 3-D matching aUGMENTATION and 3-SET PACKING AUGMENTATION problems.

## C. Improved Packing Algorithms

The method of greedy localization has been heavily used in early algorithms for matching and packing problems $[19,64]$. The method takes advantage of the fact that information for a larger matching/packing can be obtained from a given smaller matching/packing, which narrows down the size of the search space during the construction of the larger matching/packing. We show that this property can be significantly enhanced and more effectively used to develop algorithms for the 3-SET PACKING AUGMENTATION problem.

Lemma C. 1 Let $\left(S, P_{k}\right)$ be an instance of 3-SEt Packing augmentation, where $P_{k}$ is a $k$-packing in $S$. If $S$ also has $(k+1)$-packings, then there exists a $(k+1)$-packing $P_{k+1}$ in $S$ such that every set in $P_{k}$ contains at least two elements in $\operatorname{Val}\left(P_{k+1}\right)$.

Proof. We prove the lemma by contradiction. Suppose that the lemma does not hold. Then there is a $k$-packing $P_{k}$ such that for every $(k+1)$-packing $P$ in $S$, there is a set in $P_{k}$ that contains at most one element in $\operatorname{Val}(P)$. Let $P_{k+1}$ be a $(k+1)$ packing in $S$ such that the number of common sets in $P_{k}$ and $P_{k+1}$ is maximized over all $(k+1)$-packings in $S$. By our assumption, there is a set $\rho$ in $P_{k}$ that contains at most one element in $\operatorname{Val}\left(P_{k+1}\right)$.

Case 1. Exactly one element $a$ in the set $\rho$ is in $\operatorname{Val}\left(P_{k+1}\right)$. Then let $\rho^{\prime}$ be the set in $P_{k+1}$ that contains the element $a$. Since no other element in $\rho$ is in $\operatorname{Val}\left(P_{k+1}\right)$, if we replace $\rho^{\prime}$ in $P_{k+1}$ by $\rho$, we get a new $(k+1)$-packing that has one more common set (i.e., $\rho$ ) with the $k$-packing $P_{k}$ (note that $\rho^{\prime}$ cannot be in $P_{k}$ because $\rho^{\prime}$ and $\rho$ share a common element $a$ while $\rho$ contains another two elements not in $\left.\operatorname{Val}\left(P_{k+1}\right)\right)$. This contradicts our assumption that the $(k+1)$-packing $P_{k+1}$ maximizes the number of common sets with $P_{k}$.

Case 2. No element in $\rho$ is in $\operatorname{Val}\left(P_{k+1}\right)$. Since $P_{k}$ contains $k$ sets while $P_{k+1}$ contains $k+1$ sets, there must be a set $\rho^{\prime \prime}$ in $P_{k+1}$ that is not in $P_{k}$. Since $\rho$ contains no element in $\operatorname{Val}\left(P_{k+1}\right)$, replacing $\rho^{\prime \prime}$ in $P_{k+1}$ by $\rho$ gives a new $(k+1)$-packing that has one more common set (i.e., $\rho$ ) with $P_{k}$, again contradicting the assumption that the $(k+1)$-packing $P_{k+1}$ maximizes the number of common sets with $P_{k}$.

This contradiction shows that the set $\rho$ in $P_{k}$ that contains at most one element in $\operatorname{Val}\left(P_{k+1}\right)$ cannot exist.

According to Lemma C.1, to construct a $(k+1)$-packing from a given instance ( $S, P_{k}$ ) of 3 -set packing augmentation, we can aim at the $(k+1)$-packing $P_{k+1}$ with the property described in the lemma. The advantage of this $(k+1)$-packing $P_{k+1}$ is that at least $2 k$ elements in $P_{k+1}$ are already present in the $k$-packing $P_{k}$, and we only need to identify at most $k+3$ other elements in $P_{k+1}$. We use the technique of color-coding, first introduced by Alon, Yuster, and Zwick [3], to search for these elements that are in $P_{k+1}$ but not in $P_{k}$.

Let $B$ be a set of elements. A coloring of $B$ is a function mapping $B$ to the natural numbers $\{1,2, \ldots\}$, and an $h$-coloring of $B$ is a function mapping $B$ to $\{1,2, \ldots, h\}$. A subset $B^{\prime}$ of $B$ is colored properly by a coloring $f$ if no two elements in $B^{\prime}$ are colored with the same color under $f$. A collection $\mathcal{C}$ of $h$-colorings of a set $B$ is an $h$-color coding scheme if for every subset $B^{\prime}$ of $h$ elements in $B$, there is an $h$-coloring in $\mathcal{C}$ that colors $B^{\prime}$ properly. The following proposition has been proved in [25].

Proposition C. 2 [25] For any finite set $B$ and any integer $h$, there is an $h$-color coding scheme $\mathcal{C}$ of $O^{*}\left(6.1^{h}\right) h$-colorings of the set $B$. Moreover, the $h$-colorings in $\mathcal{C}$ can be constructed and enumerated in time $O^{*}\left(6.1^{h}\right)$.

Let $S$ be a collection of 3 -sets and let $f$ be a coloring of the set $\operatorname{Val}(S)$. We say that a packing $P$ in $S$ is colored properly if the set $\operatorname{Val}(P)$ is colored properly under the coloring $f$. Let $\left(S, P_{k}\right)$ be an instance of 3 -Set Packing augmentation. Since the set of elements that are in $\operatorname{Val}\left(P_{k+1}\right)$ but not in $\operatorname{Val}\left(P_{k}\right)$ contains at most $k+3$ elements, by introducing $3 k$ new colors to properly color the $3 k$ elements in $P_{k}, P_{k+1}$ can be colored properly with at most $4 k+3$ colors.

Lemma C. 3 Let $\left(S, P_{k}\right)$ be an instance of 3-set packing augmentation, and let $P_{k+1}$ be a $(k+1)$-packing in $S$ such that each 3-set in $P_{k}$ contains at least two elements in $\operatorname{Val}\left(P_{k+1}\right)$. Then there is a collection $\mathcal{C}_{0}$ of $O^{*}\left(6.1^{k}\right)(4 k+3)$-colorings of the set
$\operatorname{Val}(S)$ in which at least one properly colors $P_{k+1}$. Moreover, the collection $\mathcal{C}_{0}$ can be constructed in time $O^{*}\left(6.1^{k}\right)$.

Now we turn to the problem of constructing a properly colored $(k+1)$-packing $P_{k+1}$. Alon, Yuster, and Zwick [3] in their seminal work on color-coding suggested a general principle in which a $(3 k+3)$-coloring that properly colors the $3 k+3$ elements is first constructed in $\operatorname{Val}\left(P_{k+1}\right)$, then a dynamic programming process is applied to find the properly colored $(k+1)$-packing $P_{k+1}$. Koutis [71] proposed an algebraic formulation to find the properly colored $(k+1)$-packing $P_{k+1}$. Fellows et al. [44] considered a more general approach that first uses $g$ colors to properly color the $(k+1)$-packing $P_{k+1}$, where $g \geq 3 k+3$, then perform a dynamic programming algorithm. For completeness, we present such a generalized dynamic programming algorithm in detail, as given in Fig. 1, verify its correctness, and analyze its precise complexity.

Lemma C. 4 The algorithm $3 \operatorname{SetPack}(S, k, f, g)$ runs in time $O^{*}\left(\sum_{j=0}^{k}\binom{g}{3 j}\right)$, and constructs a properly colored $k$-packing in $S$ if such $k$-packings exist.

Proof. From steps 4.1-4.2 of the algorithm, it can be seen that every collection $P$ of 3 -sets added to the super-collection $\mathcal{Q}$ in step 4.4 is a properly colored packing. Therefore, if the algorithm returns a packing in step 5 , the packing must be a properly colored $k$-packing.

For each $i$, let $S_{i}=\left\{\rho_{1}, \ldots, \rho_{i}\right\}$. We prove by induction on $i$ that for all $j \leq k$, if $S_{i}$ has a properly colored $j$-packing $P_{j}$, then after the $i$-th execution of the for-loop in step 4 of the algorithm, the super-collection $\mathcal{Q}$ contains a properly colored $j$-packing $P_{j}^{\prime}$ such that $P_{j}$ and $P_{j}^{\prime}$ use exactly the same $3 j$ colors.

Algorithm 3SetPack $(S, k, f, g)$
input: A collection $S$ of 3 -sets, an integer $k$, a $g$-coloring $f$ of $\operatorname{Val}(S)$ output: A properly colored $k$-packing if such a packing exists

1. remove all 3 -sets in $S$ in which any two elements have the same color;
2. Let the remaining 3 -sets in $S$ be $\rho_{1}, \rho_{2}, \ldots, \rho_{n}$;
3. $\mathcal{Q}=\{\emptyset\}$;
4. for $i=1$ to $n$ do
4.1. for each packing $P$ in $\mathcal{Q}$ such that no element in $P$ is colored with the same color as an element in $\rho_{i}$ do
4.2. $\quad P^{\prime}=P \cup\left\{\rho_{i}\right\} ;$
4.3. $\quad$ if $P^{\prime}$ is a $j$-packing with $j \leq k$ and $\mathcal{Q}$ contains no packing that uses exactly the same colors as that used by $P^{\prime}$ then add $P^{\prime}$ to $\mathcal{Q}$;
4.4.
5. return a $k$-packing in $\mathcal{Q}$ if such a packing exists.

Fig. 1. Dynamic programming for 3-SET PACKING

The initial case $i=0$ is trivial since $\mathcal{Q}=\{\emptyset\}$. Consider $i \geq 1$. Suppose that the collection $S_{i}$ has a properly colored $j$-packing $P_{j}=\left\{\rho_{i_{1}}, \rho_{i_{2}}, \ldots, \rho_{i_{j}}\right\}$, where $1 \leq i_{1}<i_{2}<\cdots<i_{j} \leq i$. Then the collection $S_{i_{j}-1}$ contains the properly colored $(j-1)$-packing $P_{j-1}=\left\{\rho_{i_{1}}, \rho_{i_{2}}, \ldots, \rho_{i_{j-1}}\right\}$. By the inductive hypothesis, after the $\left(i_{j}-1\right)$-st execution of the for-loop in step 4 , the super-collection $\mathcal{Q}$ contains a properly colored $(j-1)$-packing $P_{j-1}^{\prime}$ such that the $(j-1)$-packings $P_{j-1}$ and $P_{j-1}^{\prime}$ use exactly the same $3(j-1)$ colors. Since $P_{j}=\left\{\rho_{i_{1}}, \rho_{i_{2}}, \ldots, \rho_{i_{j}}\right\}$ is a properly colored $j$-packing, and $P_{j-1}=\left\{\rho_{i_{1}}, \rho_{i_{2}}, \ldots, \rho_{i_{j-1}}\right\}$ and $P_{j-1}^{\prime}$ use exactly the same $3(j-1)$ colors, no element in $\operatorname{Val}\left(P_{j-1}^{\prime}\right)$ is colored with the same color as an element in the set $\rho_{i_{j}}$. Therefore, in the $i_{j}$-th execution of the for-loop in step 4 , a properly colored $j$-packing $P_{j-1}^{\prime} \cup\left\{\rho_{i_{j}}\right\}$ will be added to the super-collection $\mathcal{Q}$ if no properly
colored $j$-packing that uses exactly the same $3 j$ colors exists in $\mathcal{Q}$ yet. Note that the $j$-packing $P_{j-1}^{\prime} \cup\left\{\rho_{i_{j}}\right\}$ and the $j$-packing $P_{j}$ use exactly the same $3 j$ colors. Therefore, after the $i_{j}$-th execution of the for-loop in step 4 , a $j$-packing that uses exactly the same $3 j$ colors as $P_{j}$ will exist in the super-collection $\mathcal{Q}$. Finally, since packings in $\mathcal{Q}$ are never removed from $\mathcal{Q}$ and $i_{j} \leq i$, we conclude that after the $i$-th execution of the for-loop in step 4 , a $j$-packing that uses exactly the same $3 j$ colors as $P_{j}$ will exist in the super-collection $\mathcal{Q}$. This completes the inductive proof.

Now if we let $i=n$, for any $j \leq k$, if the original collection $S$ contains a properly colored $j$-packing $P_{j}$, then the super-collection $\mathcal{Q}$ contains a $j$-packing that uses exactly the same $3 j$ colors as $P_{j}$. In particular, if the collection $S$ contains properly colored $k$-packings, then the algorithm $3 \operatorname{SetPack}(S, k, f, g)$ must return a properly colored $k$-packing.

Finally, we analyze the complexity of the algorithm. For each $0 \leq j \leq k$ and for each set of $3 j$ colors, the super-collection $\mathcal{Q}$ keeps at most one properly colored $j$-packing that uses exactly these $3 j$ colors. Since there are $\binom{g}{3 j}$ different subsets of $3 j$ colors over a total of $g$ colors, the total number of packings recorded in $\mathcal{Q}$ is bounded by $\sum_{j=0}^{k}\binom{g}{3 j}$. For each $i, 1 \leq i \leq k$, we examine each packing $P$ in $\mathcal{Q}$ in step 4.1 and check if we can construct a larger packing by adding the set $\rho_{i}$ to the packing $P$. This can be done for each packing $P$ in time $O(k)$. In consequence, the algorithm $\operatorname{3SetPack}(S, k, f, g)$ runs in time $O\left(n k \sum_{j=0}^{k}\binom{g}{3 j}\right)=O^{*}\left(\sum_{j=0}^{k}\binom{g}{3 j}\right)$.

Combining Lemmas C.1, C.3, and C.4, the 3-Set packing augmentation problem can be solved in time $O^{*}\left(6.1^{k} 2^{4 k}\right)=O^{*}\left(4.61^{3 k}\right)$.

Theorem C. 5 The 3-set packing augmentation problem can be solved in time $O^{*}\left(4.61^{3 k}\right)$.

Corollary C. 6 The 3 -SET PaCking problem can be solved in time $O^{*}\left(4.61^{3 k}\right)$.

## D. Matching Algorithms Further Improved

All the previous results are applicable to the 3-D matching problem. In fact, if we regard each triple as a 3 -set, then each instance $S_{M}$ of 3-D matching is also an instance $S_{P}$ of 3 -SET PACKING, and a triple set is a matching in $S_{M}$ if and only if it is a packing in $S_{P}$.

As shown in the previous section, a dynamic programming algorithm is used as a second stage in parameterized algorithms for 3-SET PACKING/3-D MATCHING. In this section, we develop a new technique for the dynamic programming stage for the 3-D matching problem so that fewer colors will be needed. This technique has two advantages. First, the use of fewer colors will significantly reduce the time complexity of the coloring stage. Second, since fewer colors are used, the number of different color sets is reduced, which will reduce the time complexity of the dynamic programming stage remarkably.

Let the universal triple set be $U=X \times Y \times Z$, where $X, Y$, and $Z$ are three pairwise disjoint symbol sets. The symbols in the sets $X, Y$, and $Z$ will be called the symbols in column-1, column-2, and column-3, respectively.

Definition Let $p$ and $q$ be any two indices in the index set $\{1,2,3\}$, and let $S$ be a set of triples in $U$. A matching $M$ in the set $S$ is $(p, q)$-properly colored by a coloring $f$ of $\operatorname{Val}^{p}(M) \cup \operatorname{Val}^{q}(M)$ if no two symbols in $\operatorname{Val}^{p}(M) \cup \operatorname{Val}^{q}(M)$ are colored with the same color under $f$.

Theorem D. 1 Let $p$ and $q$ be any two indices in the index set $\{1,2,3\}$. There is an algorithm of time $O^{*}\left(\sum_{i=0}^{k}\binom{g}{2 i}\right)$ that, on an integer $k$ and a set $S$ of triples in which the symbols in $\operatorname{Val}^{p}(S) \cup \operatorname{Val}^{q}(S)$ are colored by ag-coloring $f$, constructs a

Algorithm 3DMatch( $S, k, f, g ; p, q)$
input: A set $S$ of triples, an integer $k$, a $g$-coloring $f$ of the symbols in $\operatorname{Val}^{p}(S) \cup \operatorname{Val}^{q}(S)$
output: A $(p, q)$-properly colored $k$-matching in $S$ if such a matching exists

1. remove any triples in $S$ in which any two symbols have the same color under $f$;
2. let the set of remaining triples be $S^{\prime}$;
3. $r=\{1,2,3\}-\{p, q\}$;
4. let the symbols in $\operatorname{Val}^{r}\left(S^{\prime}\right)$ be $x_{1}, x_{2}, \ldots, x_{m}$;
5. $\mathcal{Q}_{\text {old }}=\{\emptyset\} ; \mathcal{Q}_{\text {new }}=\{\emptyset\}$;
6. for $i=1$ to $m$ do
6.1. for each set $C$ of symbol pairs in $\mathcal{Q}_{\text {old }}$ do
6.2. for each $t \in S^{\prime}$ with $\operatorname{Val}^{r}(t)=x_{i}$ do
6.3. if no symbol in $C$ is of the same color as a symbol in $\operatorname{Val}^{p}(t) \cup \operatorname{Val}^{q}(t)$
6.4. $\quad$ then $C^{\prime}=C \cup\left\{\left(\operatorname{Val}^{p}(t), \operatorname{Val}^{q}(t)\right)\right\}$;
6.5. if $C^{\prime}$ contains no more than $k$ symbol pairs and $\mathcal{Q}_{\text {new }}$ contains no set of symbol pairs that uses exactly the same colors as that used by $C^{\prime}$
6.6. then add $C^{\prime}$ to $\mathcal{Q}_{\text {new }}$;
6.7. $\mathcal{Q}_{\text {old }}=\mathcal{Q}_{\text {new }}$;
7. return a set $C$ of $k$ symbol pairs in $\mathcal{Q}_{\text {old }}$ if such a set exists.

Fig. 2. Dynamic programming for 3-D matching
$(p, q)$-properly colored $k$-matching in $S$ when such matchings exist in $S$.

Proof. Consider the algorithm in Fig. 11. By steps 6.3-6.6, for every set $C$ in the collection $\mathcal{Q}_{\text {old }}$, all symbols in $C$ are from $\operatorname{Val}^{p}(S) \cup \operatorname{Val}^{q}(S)$, and no two symbols in $C$ are of the same color. The algorithm $\mathbf{3 D M a t c h}(S, k, f, g ; p, q)$ either outputs a set of $k$ symbol pairs in the collection $\mathcal{Q}_{\text {old }}$ or reports that no $(p, q)$-properly colored $k$-matchings exist in $S$. We say that a set $C=\left\{w_{1}, \ldots, w_{i}\right\}$ of $i$ symbol pairs is
extendable to an $i$-matching in $S$ if there is an $i$-matching $M=\left\{t_{1}, \ldots, t_{i}\right\}$ in $S$ such that for all $j$, the pair $\left(\operatorname{Val}^{p}\left(t_{j}\right), \operatorname{Val}^{q}\left(t_{j}\right)\right)$ is identical to the symbol pair $w_{j}$. For each $i$, let $S_{i}^{\prime}$ be the set of triples in $S^{\prime}$ whose symbols in column- $r$ are among $\left\{x_{1}, x_{2}, \ldots, x_{i}\right\}$. For a matching $M$, we will denote by $c l(M)=\left\{f(y) \mid y \in \operatorname{Val}^{p}(M) \cup \operatorname{Val}^{q}(M)\right\}$ the set of colors used by the symbols in $\operatorname{Val}^{p}(M) \cup \operatorname{Val}^{q}(M)$.

We prove the following claim by induction on $i$ :

Claim. For each $i, 0 \leq i \leq m$, and for all $h \leq k$, there is a $(p, q)$-properly colored $h$-matching $M_{h}$ in $S_{i}^{\prime}$ if and only if after the $i$-th execution of the loop 6.1-6.7 of algorithm $\mathbf{3 D M a t c h}(S, k, f, g ; p, q)$, the collection $\mathcal{Q}_{\text {old }}$ contains a set $C_{h}$ of $h$ symbol pairs such that the set of colors used for the symbols in $C_{h}$ is exactly $\operatorname{cl}\left(M_{h}\right)$. Moreover, each set $C_{h}$ of $h$ symbol pairs in the collection $\mathcal{Q}_{\text {old }}$ after the $i$-th execution of the loop 6.1-6.7 is extendable to an $h$-matching in $S_{i}^{\prime}$.

The case $i=0$ is obvious because we initially set $\mathcal{Q}_{\text {old }}$ to $\{\emptyset\}$. Consider $i \geq 1$. First note that the claim is always true for $h=0$ because the collection $\mathcal{Q}_{\text {old }}$ always contains the empty set $\emptyset$ while the set $S_{i}^{\prime}$ always contains a 0 -matching (which by the definition is ( $p, q$ )-properly colored).

Suppose that after the $i$-th execution of the loop 6.1-6.7, the collection $\mathcal{Q}_{\text {old }}$ contains a set $C_{h}$ of $h$ symbol pairs, where $h \geq 1$. Suppose that the set $C_{h}$ was created during the $j$-th execution of the loop 6.1-6.7, where $j \leq i$, by adding a symbol pair $\left(\operatorname{Val}^{p}(t), \operatorname{Val}^{q}(t)\right)$ to a set $C_{h-1}$ of $h-1$ symbol pairs, where $t$ is a triple with $\operatorname{Val}^{r}(t)=x_{j}$ and the set $C_{h-1}$ is contained in $\mathcal{Q}_{\text {old }}$ after the $(j-1)$-st execution of the loop 6.1-6.7. By the inductive hypothesis, the set $C_{h-1}$ is extendable to an $(h-1)-$ matching $M_{h-1}$ in $S_{j-1}^{\prime}$, which is obviously $(p, q)$-properly colored. Since no symbol in $C_{h-1}$ uses the same color as a symbol in $\operatorname{Val}^{p}(t) \cup \operatorname{Val}^{q}(t)$, and the matching $M_{h-1}$
does not contain the symbol $x_{j}$, the set $M_{h}=M_{h-1} \cup\{t\}$ makes a $(p, q)$-properly colored $h$-matching in $S_{j}^{\prime}$. Since $j \leq i$ and $S_{j}^{\prime} \subseteq S_{i}^{\prime}$, we conclude that the set $S_{i}^{\prime}$ contains a $(p, q)$-properly colored $h$-matching $M_{h}$ such that the symbols in the set $C_{h}$ use exactly the color set $\operatorname{cl}\left(M_{h}\right)$. Moreover, it is obvious that the symbol set $C_{h}$ is extendable to the $h$-matching $M_{h}$.

To prove the other direction, suppose that the set $S_{i}^{\prime}$ contains a $(p, q)$-properly colored $h$-matching $M_{h}$.

Case 1. There is a $(p, q)$-properly colored $h$-matching $M_{h}^{\prime}$ in $S_{j}^{\prime}$ for some $j<i$ such that $c l\left(M_{h}^{\prime}\right)=c l\left(M_{h}\right)$. By the inductive hypothesis, after the $j$-th execution of the loop 6.1-6.7, the collection $\mathcal{Q}_{\text {old }}$ contains a set $C_{h}$ of $h$ symbol pairs such that (1) the set of colors used for the symbols of $C_{h}$ is exactly $c l\left(M_{h}^{\prime}\right)$; and (2) $C_{h}$ is extendable to an $h$-matching in $S_{j}^{\prime}$. Since $j<i, S_{j}^{\prime} \subseteq S_{i}^{\prime}$, and we never remove symbol pairs from $\mathcal{Q}_{\text {old }}$, we conclude that in this case, after the $i$-th execution of the loop 6.1-6.7, the set $C_{h}$ is still contained in the collection $\mathcal{Q}_{\text {old }}$ such that (1) the set of colors used for the symbols of $C_{h}$ is exactly $\operatorname{cl}\left(M_{h}^{\prime}\right)=\operatorname{cl}\left(M_{h}\right)$; and (2) $C_{h}$ is extendable to an $h$-matching in $S_{j}^{\prime} \subseteq S_{i}^{\prime}$.

Case 2. There is no $(p, q)$-properly colored $h$-matching $M_{h}^{\prime}$ in $S_{j}^{\prime}$ for any $j<i$ such that $\operatorname{cl}\left(M_{h}^{\prime}\right)=\operatorname{cl}\left(M_{h}\right)$. Then by the inductive hypothesis, after the $j$-th execution of the loop 6.1-6.7 for any $j<i$, the collection $\mathcal{Q}_{\text {old }}$ contains no set $C$ of symbol pairs such that the symbols in $C$ use exactly the color set $c l\left(M_{h}\right)$. Let the $(p, q)$-properly colored $h$-matching $M_{h}$ be $M_{h}=\left\{t_{1}, \cdots, t_{h}\right\}$, where for each $j$, $\operatorname{Val}^{r}\left(t_{j}\right)=x_{d_{j}}$, with $d_{1}<\cdots<d_{h-1}<d_{h}$. In this case, we must have $d_{h}=i$ and $\operatorname{Val}^{r}\left(t_{h}\right)=x_{i}$. Let $y=\operatorname{Val}^{p}\left(t_{h}\right)$ and $z=\operatorname{Val}^{q}\left(t_{h}\right)$. Since $d_{h-1}<d_{h}=i$, the triple set $M_{h-1}=M_{h}-\left\{t_{h}\right\}$ is a $(p, q)$-properly colored $(h-1)$-matching in $S_{i-1}^{\prime}$. By the inductive hypothesis, after the $(i-1)$-st execution of the loop 6.1-6.7 in the algorithm, the collection $\mathcal{Q}_{\text {old }}$ contains a set $C_{h-1}$ of $h-1$ symbol pairs such that the set of colors used for the
symbols in $C_{h-1}$ is exactly $c l\left(M_{h-1}\right)$. Now in the $i$-th execution of the loop 6.1-6.7 when the set $C_{h-1}$ and the triple $t_{h}$ are examined in step 6.3 , a set $C$ of symbol pairs using the color set $\operatorname{cl}\left(M_{h-1}\right) \cup\{f(y), f(z)\}=\operatorname{cl}\left(M_{h}\right)$ will be created. Therefore, after the $i$-th execution of the loop 6.1-6.7, a set $C_{h}$ of $h$ symbol pairs using the color set $c l\left(M_{h}\right)$ must be contained in the collection $\mathcal{Q}_{\text {old }}$. Suppose that the set $C_{h}$ was created during the $i$-th execution by adding a symbol pair $\left(\operatorname{Val}^{p}\left(t_{h}^{\prime}\right), \operatorname{Val}^{q}\left(t_{h}^{\prime}\right)\right)$ to a set $C_{h-1}^{\prime}$ of $h-1$ symbol pairs, where $t_{h}^{\prime}$ satisfies $\operatorname{Val}^{r}\left(t_{h}^{\prime}\right)=x_{i}$ and $C_{h-1}^{\prime}$ is contained in the collection $\mathcal{Q}_{\text {old }}$ after the $(i-1)$-st execution of the loop 6.1-6.7 (note that $t_{h}^{\prime}$ and $C_{h-1}^{\prime}$ are not necessarily $t_{h}$ and $C_{h-1}$, respectively). By the inductive hypothesis, the set $C_{h-1}^{\prime}$ is extendable to a $(p, q)$-properly colored $(h-1)$-matching $M_{h-1}^{\prime}$ in $S_{i-1}^{\prime}$. In consequence, the set $C_{h}$ is extendable to the $(p, q)$-properly colored $h$-matching $M_{h-1}^{\prime} \cup\left\{t_{h}^{\prime}\right\}$ in $S_{i}^{\prime}$. This completes the proof of the claim.

By the claim and let $i=m$, the algorithm $\operatorname{3DMatch}(S, k, f, g ; p, q)$ returns a set $C_{k}$ of $k$ symbol pairs if and only if the triple set $S$ contains a $(p, q)$-properly colored $k$-matching, and the set $C_{k}$ is extendable to a $k$-matching in $S$. To construct such a $k$-matching from $C_{k}$, we can use the graph matching technique suggested in [19]. Formally, from the set $C_{k}$ of symbol pairs, we construct a bipartite graph $B_{k}=\left(V_{L} \cup V_{R}, E\right)$, where $V_{L}$ contains $k$ vertices, corresponding to the $k$ symbol pairs in $C_{k}$, and $V_{R}$ is the set of all symbols in $\operatorname{Val}^{r}(S)$. There is an edge in $B_{k}$ from a vertex $(y, z)$ in $V_{L}$ to a vertex $x$ in $V_{R}$ if and only if the symbols $y, z$, and $x$ form a triple in $S$. It is easy to see that a $(p, q)$-properly colored $k$-matching $M_{k}$ in $S$ can be obtained by constructing a graph matching of $k$ edges in the bipartite graph $B_{k}$, which takes polynomial time [28].

In terms of the time complexity of the above algorithm, note that since for each set of $2 i$ colors, we record at most one set of symbol pairs that uses exactly these $2 i$ colors, the collection $\mathcal{Q}_{\text {old }}$ contains at most $\sum_{i=0}^{k}\binom{g}{2 i}$ sets of symbol pairs.

For each set $C$ of symbol pairs in $\mathcal{Q}_{\text {old }}$, steps 6.2-6.6 of the algorithm take time polynomial in $n$ and $k$. Therefore, each execution of the loop 6.1-6.7 of the algorithm runs in time $O^{*}\left(\sum_{i=0}^{k}\binom{g}{2 i}\right)$. In consequence, the running time of the algorithm $\operatorname{3DMatch}(S, k, f, g ; p, q)$ is bounded by $O^{*}\left(\sum_{i=0}^{k}\binom{g}{2 i}\right)$.

To solve the 3-D matching augmentation problem $\left(S, M_{k}\right)$, we only color two columns of a $(k+1)$-matching properly if it exists. In this case, by Lemma C.1, there is a $(k+1)$-matching $M_{k+1}$ such that $M_{k+1}$ has two columns that contain at least $4 k / 3$ symbols in $M_{k}$. Thus at most $2 k / 3+2$ symbols in these two columns in $M_{k+1}$ are missing in $M_{k}$. By introducing $2 k$ new colors for each symbol in these two columns in $M_{k}$ and by Proposition C.2, in time $O^{*}\left(6.1^{2 k / 3}\right)$, the two columns of $M_{k+1}$ can be colored properly into $8 k / 3+2$ colors. By Theorem D.1, the 3-D matching AUGMENTATION problem $\left(S, M_{k}\right)$ can be solved in time $O^{*}\left(6 \cdot 1^{2 k / 3} 2^{8 k / 3}\right)=O^{*}\left(2.77^{3 k}\right)$.

Theorem D. 2 The 3-D matching problem can be solved in time $O^{*}\left(2.77^{3 k}\right)$.

If we use a randomized color-coding scheme that properly colors a subset of size $k$ into $k$ colors with high probability in time $O^{*}\left(e^{k}\right)$ [3], then the time complexity to solve 3-D matching problem can be improved to $O^{*}\left(e^{2 k / 3} 2^{8 k / 3}\right)=O^{*}\left(2.32^{3 k}\right)$.

Theorem D. 3 The 3-D matching problem can be solved by a randomized algorithm of time $O^{*}\left(2.32^{3 k}\right)$.

## E. Final Remarks

Recently there has been much interest in parameterized algorithms for graph packing problems, i.e., algorithms for constructing $k$ disjoint isomorphic subgraphs in a given graph [45, 70, 83, 86]. In particular, Fellows et al. [45] presented a parameterized algorithm of time $O^{*}\left(2^{2 k \log k+1.869 k}\right)$ for packing $k$ vertex-disjoint triangles
in a given graph, and Mathieson, Prieto, and Shaw [83] proposed a parameterized algorithm of time $O^{*}\left(2^{4.5 k \log k+4.5 k}\right)$ for packing $k$ edge-disjoint triangles in a given graph. Since these problems can be trivially reduced to the 3 -SET PACKING problem, by Corollary C. 6 , they can be solved in time $O^{*}\left(4.61^{3 k}\right)$. This again gives significant improvements over the previous algorithms.

There are more new results by the time this dissertation is completed: the best randomized algorithm of $O^{*}\left(2^{k}\right)$ [72] for the 3-D MATCHING and 3-SET PACKING problems, the best deterministic algorithm of $O^{*}\left(4^{(r-1) k}\right)$ for the weighted R-D matching problem [46], and the best deterministic algorithm of $O^{\left(2^{(2 r-1) k}\right)}$ for the weighted R-SET PACKING problem.

## CHAPTER III

## MULTIWAY CUT*

In this chapter, we present an fpt-algorithm of running time $O^{*}\left(4^{k}\right)$ for the MULTIWAY CUT problem, which significantly improves the previous best algorithm of running time $O^{*}\left(4^{k^{3}}\right)$.

The previous algorithm uses only the parameter $k$, the size of a multiway cut to search, as the measure. Our algorithm considers both the parameter $k$ and the minimum cut $m$ between a vertex and the other vertices as measures. We find an operation which either decreases $k$ by 1 or increases $m$ by 1 . Moreover, we uncover that either such operation is executed or the vertices for further consideration decreases. With these discoveries, we design an algorithm of running time $O^{*}\left(4^{k}\right)$.

## A. Introduction

The multiterminal cut problem is a well-known problem, and has been extensively studied ([16, 66, 85]). Applications of this problem are found in distributed computing [96], VLSI [27], computer vision [15], and many other fields. The problem is defined as follows: given an undirected graph $G=(V, E)$ and a set of $l$ vertices $\left\{t_{1}, \ldots, t_{l}\right\}$ in $G$ (the vertices $t_{i}$ are called terminals), find an edge set $E^{\prime}$ of minimum size in $G$ such that after the deletion of $E^{\prime}$, no two terminals are in the same connected component. This problem is NP-hard for general graphs for any fixed integer $l \geq 3$, and is also NP-hard for planar graphs when $l$ is not fixed [30].

A generalization of the multiterminal cut problem is the minimum node

[^1]mULTIWAY CUT problem, which, for a given graph and a given set of terminals, is to find a vertex set $S$ of minimum size such that after the deletion of $S$, no two terminals are in the same connected component. The minimum node multiway cut problem is at least as hard as the MULTITERMINAL CUT problem, as the latter can be reduced to the former in time $O(|V|+|E|)$, if we require that no terminal be in the separator $S$ [29]. Therefore, the minimum node multiway cut problem is also NP-hard if the number $l$ of terminals is at least 3 .

When there are only two terminals $s$ and $t$, the multiterminal cut problem and the Minimum node multiway cut problem become the edge version and the vertex version of the minimum $s$ - $t$ CUT problem, respectively. According to the maxflow min-cut theorem [74], the MINIMUM $s$ - $t$ CUT problem, for both the edge version and the vertex version, can be solved via algorithms for the MAXIMUM $s$ - $t$ FLOW problem. For an undirected graph $G$ of $n$ vertices and $m$ edges, the maximum $s-t$ FLOW problem can be solved in time $O\left(n^{7 / 6} m^{2 / 3}\right)$ [67]. In consequence, the MULTITERMINAL CUT problem and the MINIMUM NODE MULTIWAY CUT problem can also be solved in time $O\left(n^{7 / 6} m^{2 / 3}\right)$.

A natural extension of the minimum node multiway cut problem is to have a collection of terminal sets, instead of a collection of individual terminals. Formally, let $G=(V, E)$ be an undirected graph, and let $\left\{T_{1}, \ldots, T_{l}\right\}$ be a collection of terminal sets where each $T_{i}$ is a subset of vertices in $G$. A separator $S$ for $\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}$ is a subset of vertices in $G$ such that no vertex in $S$ is in any terminal set, and after deleting $S$ from the graph $G$, no connected component in the resulting graph contains vertices from more than one terminal set.

In certain real world applications, one may expect that the size of the separator be small. For example, suppose that we are given a network (i.e., a graph) $G=(V, E)$ and a collection of network node groups $\left\{T_{1}, \ldots, T_{l}\right\}$ in $G$, and we want to monitor the
message communication among the node groups. A separator for $\left\{T_{1}, \ldots, T_{l}\right\}$ in the network $G$ will well serve for this purpose: any communication path between any two node groups must pass through at least one node in the separator. Therefore, if we set up a monitor process in each of the nodes in the separator, then we can monitor all communications among the node groups. Naturally, we may want to limit the cost of this monitoring system by using only a small number of "monitor nodes" in the network $G$.

This motivates a parameterized version of the MINIMUM NODE MULTIWAY CUT problem, which will be called the Parameterized node multiway cut problem and is defined as follows: given an undirected graph $G=(V, E)$, a collection of pairwise disjoint terminal sets $\left\{T_{1}, \ldots, T_{l}\right\}$ (where each $T_{i}$ is a subset of vertices in $G$ ), and a parameter $k$, either construct a separator of at most $k$ vertices in $G$, or report that no such a separator exists. Our goal is, for the Parameterized node MULTIWAY CUT problem, to develop a fixed-parameter tractable algorithm [37], i.e., an algorithm whose running time is of the form $f(k) n^{c}$ with a function $f$ independent of the input size $n$ and a constant $c$. In particular, when the parameter value $k$ is small, such a fixed-parameter tractable algorithm will be practically effective. In fact, the study of fixed-parameter tractable algorithms for a variety of parameterized problems has drawn considerable attention recently and has direct impact on real word applications where the selected parameter varies in a small range [37].

It can be derived from the graph minor theory of Robertson and Seymour [37] that there is a fixed-parameter tractable algorithm for the PARAMETERIZED NODE multiway cut problem. However, the proof is not constructive. An explicit constructive algorithm for the problem was given by Marx [82], who developed an algorithm of running time $O\left(n^{5} 4^{k^{3}}\right)$ for the PARAMETERIZED NODE MULTIWAY CUT problem for its original version (i.e., in which each terminal set is restricted to con-
tain a single terminal). To our knowledge, this is the only known constructive fixedparameter tractable algorithm for the problem.

In this chapter, we present an algorithm of running time $O\left(n^{3} k 4^{k}\right)$ for the PARAMETERIZED NODE MULTIWAY CUT problem, which significantly improves the algorithm given in [82]. In the real world of computing, this improvement makes it become possible to practically solve the problem for some reasonable values of the parameter $k$. For example, for the case of $k=10$, our algorithm has running time $O\left(n^{3} 4^{10}\right)$, which is practically feasible using the currently available computation power. On the other hand, the algorithm in [82] in this case has running time $O\left(n^{5} 4^{1000}\right)$, which is totally infeasible from the practical point of view. Theoretically, our result gives the first polynomial time algorithm for the MINIMUM NODE MULTIWAY CUT problem when the size of the optimal separator is of order $O(\log n)$.

Finally, we remark that the techniques we developed in this chapter seem to be very powerful for solving various kinds of multiway cut problems. In particular, very recently the techniques have been extended to directed graphs, and led to a fixed parameter tractable algorithm for the FEEDBACK VERTEX SET problem on directed graphs [24], thus resolving an outstanding open problem in the area of parameterized computation and complexity [37, 34].

## B. Minimum V-cuts between Two Terminal Sets

We start with some terminologies. All graphs in our discussion are supposed to be directed.

Let $G=(V, E)$ be a graph and let $u$ and $v$ be two vertices in $G$. A path between $u$ and $v$ is a simple path in $G$ whose two ends are $u$ and $v$, respectively. We say that there is a path between a vertex $u$ and a vertex subset $V^{\prime}$ if there is a path between
the vertex $u$ and a vertex $v$ in the subset $V^{\prime}$. For two vertex subsets $V_{1}$ and $V_{2}$, we say that there is a path between $V_{1}$ and $V_{2}$ if there exist a vertex $u$ in $V_{1}$ and a vertex $v$ in $V_{2}$ such that there is a path between $u$ and $v$. Two paths are internally disjoint if there is no vertex that is an internal vertex for both the paths.

Let $G$ be a graph, and let $\left\{T_{1}, \ldots, T_{l}\right\}$ be a collection of pairwise disjoint terminal sets (each terminal set is a subset of vertices in $G$ ). A subset $S$ of vertices in $G$ is a separator for $\left\{T_{1}, \ldots, T_{l}\right\}$ if $S$ contains no vertex in any of the sets $T_{1}, \ldots, T_{l}$, and if after deleting all vertices in $S$ from $G$, there is no path between any two different subsets $T_{i}$ and $T_{j}$ in the resulting graph. In particular, a separator $S$ for two terminal sets $T_{1}$ and $T_{2}$ is also called a $V$-cut between the two sets $T_{1}$ and $T_{2}$.

Let $T$ be a subset of vertices in the graph $G=(V, E)$. By merging $T$ (into a single vertex), we mean the operation that first deletes all vertices in $T$ then creates a new vertex $w$ adjacent to each $v$ of the vertices in $V-T$ where $v$ is a neighbor of a vertex in $T$ in the original graph $G$.

Finally, for a subset $V^{\prime}$ of vertices in the graph $G$, we will denote by $G\left(V^{\prime}\right)$ the subgraph of $G$ that is induced by the vertex subset $V^{\prime}$.

Proposition B. 1 [17] (Menger's Theorem-Vertex Version) Let $u$ and $v$ be two distinct and nonadjacent vertices in a graph $G$. Then the maximum number of internally disjoint paths between $u$ and $v$ in $G$ is equal to the size of a minimum $V$-cut between $u$ and $v$ in $G$.

Proposition B. 1 can be generalized from the case for two vertices to the case of two vertex subsets, as given in the following lemma.

Lemma B. 2 Let $T_{1}$ and $T_{2}$ be two disjoint vertex subsets in a graph $G$ such that no vertex in $T_{1}$ is adjacent to a vertex in $T_{2}$. Then the maximum number $h$ of internally disjoint paths between $T_{1}$ and $T_{2}$ in $G$ is equal to the size of a minimum $V$-cut between
$T_{1}$ and $T_{2}$ in $G$. Moreover, for any set $\pi$ of $h$ internally disjoint paths between $T_{1}$ and $T_{2}$ in $G$, every minimum $V$-cut between $T_{1}$ and $T_{2}$ in $G$ contains exact one vertex in each of the paths in $\pi$.

Proof. Let $G^{\prime}$ be the graph obtained from the graph $G$ by merging the two vertex subsets $T_{1}$ and $T_{2}$ into two vertices $t_{1}$ and $t_{2}$, respectively. Note that $t_{1}$ and $t_{2}$ are not adjacent in $G^{\prime}$.

By the definition of the merge operation, it is easy to verify that a vertex subset $S$ is a V-cut between the vertex subsets $T_{1}$ and $T_{2}$ in the graph $G$ if and only if $S$ is a V-cut between the vertices $t_{1}$ and $t_{2}$ in the graph $G^{\prime}$. In particular, the size of a minimum $V$-cut between $T_{1}$ and $T_{2}$ in $G$ is equal to the size of a minimum V -cut between $t_{1}$ and $t_{2}$ in $G^{\prime}$. Moreover, it is also easy to verify that for any integer $h^{\prime}$, from a set of $h^{\prime}$ internally disjoint paths between $T_{1}$ and $T_{2}$ in $G$, we can construct a set of $h^{\prime}$ internally disjoint paths between $t_{1}$ and $t_{2}$ in $G^{\prime}$, and vice versa. Therefore, the maximum number of internally disjoint paths between $T_{1}$ and $T_{2}$ in $G$ is equal to the maximum number of internally disjoint paths between $t_{1}$ and $t_{2}$ in $G^{\prime}$. Now the first part of the lemma follows by applying Proposition B. 1 to the graph $G^{\prime}$.

To prove the second part of the lemma, let $S$ be a minimum V-cut, of size $h$, between $T_{1}$ and $T_{2}$ in $G$, and let $\pi$ be a set of $h$ internally disjoint paths between $T_{1}$ and $T_{2}$. The vertex set $S$ must contain at least one vertex from each of the paths in $\pi$ : otherwise there would be a path between $T_{1}$ and $T_{2}$ in $G-S$, contradicting the assumption that $S$ is a V-cut between $T_{1}$ and $T_{2}$. Moreover, the set $S$ cannot contain more than one vertex in any path in $\pi$ : otherwise $S$ would not be able to contain at least one vertex for each of the paths in $\pi$ (note that the paths in $\pi$ are internally disjoint).

Lemma B. 2 provides an efficient algorithm that constructs the maximum number of internally disjoint paths and a minimum-size V-cut between two given vertex subsets in a graph.

Lemma B. 3 Let $T_{1}$ and $T_{2}$ be two disjoint vertex subsets in a graph $G=(V, E)$ such that no vertex in $T_{1}$ is adjacent to a vertex in $T_{2}$. Then in time $O((|V|+|E|) k)$, we can decide if the size $h$ of a minimum $V$-cut between $T_{1}$ and $T_{2}$ is bounded by $k$, and in case $h \leq k$, construct $h$ internally disjoint paths between $T_{1}$ and $T_{2}$.

Proof. Let $G^{\prime}$ be the graph obtained from the graph $G$ by merging the two vertex subsets $T_{1}$ and $T_{2}$ into two vertices $t_{1}$ and $t_{2}$, respectively. As discussed in the proof of Lemma B.2, it suffices to show how to decide if the size $h$ of a minimum V-cut between $t_{1}$ and $t_{2}$ in $G^{\prime}$ is bounded by $k$, and in case $h \leq k$, how to construct $h$ internally disjoint paths between $t_{1}$ and $t_{2}$.

This can be done based on the standard approach to the maximum $t_{1}-t_{2}$ FLOW problem [17]. For this, we first transform the undirected graph $G^{\prime}$ into a directed graph by replacing each edge by two reverse arcs. Then we modify the new directed graph by replacing each vertex $u$ (except the vertices $t_{1}$ and $t_{2}$ ) by two vertices $u_{1}$ and $u_{2}$ with an arc from $u_{1}$ to $u_{2}$, connecting all $u$ 's incoming arcs to the vertex $u_{1}$ and connecting all $u$ 's outgoing arcs to the vertex $u_{2}$. Finally we set all edges to have capacity 1 . Let the resulting flow graph be $G^{\prime \prime}$.

Applying Ford-Fulkerson's standard approach using augmenting paths, in time $O((|V|+|E|) k)$, we can either construct a $t_{1}-t_{2}$ flow of value larger than $k$ in $G^{\prime \prime}$, or end up with a maximum $t_{1}-t_{2}$ flow of value $h$ bounded by $k$. In the former case, we conclude that the size of a minimum V -cut between $t_{1}$ and $t_{2}$ in $G^{\prime}$ is larger than $k$, which implies that the size of a minimum V-cut between $T_{1}$ and $T_{2}$ in $G$ is larger than
$k$. In the latter case, $h$ internally disjoint paths between $t_{1}$ and $t_{2}$ in $G^{\prime}$ can be easily constructed from the maximum $t_{1}-t_{2}$ flow of value $h$ in $G^{\prime \prime}$, from which $h$ internally disjoint paths between $T_{1}$ and $T_{2}$ in $G$ can be constructed.

## C. The Main Algorithm

Now we return back to the parameterized node multiway cut problem. Formally, an instance $\left(G,\left\{T_{1}, \ldots, T_{l}\right\}, k\right)$ of the Parameterized node multiway cut problem consists of an undirected graph $G$, a collection $\left\{T_{1}, \ldots, T_{l}\right\}$ of pairwise disjoint terminal sets (each terminal set is a vertex subset in $G$ ), and a parameter $k$. The objective is to either construct a separator of at most $k$ vertices for $\left\{T_{1}, \ldots, T_{l}\right\}$, or conclude that no such a separator exists.

Before we formally present our algorithm, we give a less formal but intuitive explanation on the basic idea of the algorithm. Let the size of a minimum V-cut between $T_{1}$ and $\bigcup_{j \neq 1} T_{j}$ be $m$.

Pick a vertex $u$ that is not in any terminal set and has a neighbor in $T_{1}$. If $u$ also has a neighbor in another terminal set $T_{i}, i \neq 1$, then we can directly include $u$ in the separator (this is necessary because the separator must separate $T_{1}$ and $T_{i}$ ), and recursively find a separator of size $k-1$ in the remaining graph. On the other hand, if $u$ has no neighbor in other terminal sets, then we compute the size $m^{\prime}$ of a minimum V-cut between the sets $T_{1}^{\prime}=T_{1} \cup\{u\}$ and $\bigcup_{i \neq 1} T_{i}$. It can be proved that we must have $m \leq m^{\prime}$. Note that by Lemma B.3, the values $m$ and $m^{\prime}$ can be computed in polynomial time.

In the case $m=m^{\prime}$, we will show that the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator of size bounded by $k$ if and only if the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator of size bounded by $k$. Then we recursively work on the new instance
$\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$. Thus, in the case of $m=m^{\prime}$, we can reduce the number of vertices that are not in the separator by 1 .

On the other hand, suppose $m<m^{\prime}$. Then we branch on the vertex $u$ in two cases, one includes $u$ in the separator and the other excludes $u$ from the separator. In the case of including the vertex $u$ in the separator, we recursively work on the instance $\left(G-\{u\},\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k-1\right)$, in which the parameter value is decreased by 1 ; and in the case of excluding the vertex $u$ from the separator, we recursively work on the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$, in which the size of the minimum V-cut between $T_{1}^{\prime}$ and $\bigcup_{i \neq 1} T_{i}$ is increased by at least 1 .

Therefore, for the given instance ( $G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k$ ), we can either (1) apply a polynomial time process that either decreases the parameter value by 1 or reduces the number of vertices not in the separator by 1, or (2) branch into two cases, of which one decreases the parameter value by 1 and the other increases the value $m$ by at least 1 (see the definition of $m$ given in the second paragraph in this section). Note that all these generated new instances will be "simpler" than the original given instance: (i) reducing the number of vertices not in the separator will narrow down our search space for the separator; (ii) an instance of parameter value bounded by 1 can be solved in polynomial time; and (iii) an instance in which the value $m$ is larger than the parameter value $k$ obviously has no separator of size bounded by $k$.

To present our formal discussions, we fix an instance $\left(G,\left\{T_{1}, \ldots, T_{l}\right\}, k\right)$ of the PARAMETERIZED NODE MULTIWAY CUT problem, where $G=(V, E)$ is a graph, and $\left\{T_{1}, \ldots, T_{l}\right\}$ is a collection of terminal sets in $G$. Let the size of a minimum V-cut between $T_{1}$ and $\bigcup_{j \neq 1} T_{j}$ be $m$. Moreover, fix a vertex $u$ that is not in any of the terminal sets but has a neighbor in the terminal set $T_{1}$. Let $T_{1}^{\prime}=T_{1} \cup\{u\}$.

Lemma C. 1 Let $m$ be the size of a minimum $V$-cut between the two sets $T_{1}$ and
$\bigcup_{j \neq 1} T_{j}$, and let $m^{\prime}$ be the size of a minimum $V$-cut between the two sets $T_{1}^{\prime}$ and $\bigcup_{j \neq 1} T_{j}$. Then $m^{\prime} \geq m$.

Proof. The lemma follows from the observation that every V-cut between the sets $T_{1}^{\prime}$ and $\bigcup_{j \neq 1} T_{j}$ is also a V-cut between the sets $T_{1}$ and $\bigcup_{j \neq 1} T_{j}$.

The following theorem is the most crucial observation for our algorithm.

Theorem C. 2 If the minimum $V$-cuts between the sets $T_{1}$ and $\bigcup_{j \neq 1} T_{j}$ and the minimum $V$-cuts between the sets $T_{1}^{\prime}$ and $\bigcup_{j \neq 1} T_{j}$ have the same size, then the instance ( $\left.G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator of size bounded by $k$ if and only if the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator of size bounded by $k$.

Proof. If the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator $S$ of size bounded by $k$, then it is obvious that $S$ is also a separator for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. In consequence, the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ also has a separator of size bounded by $k$.

Now we consider the other direction. Suppose that the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator $S_{k}$ of size bounded by $k$.

To simplify the discussion, denote by $T_{o t h e r}$ the set $\bigcup_{j \neq 1} T_{j}$. Let $S_{m}$ be a minimum V-cut between $T_{1}^{\prime}$ and $T_{\text {other }}$ (note that $S_{m}$ does not contain $u$ ). Then $S_{m}$ is also a V-cut between $T_{1}$ and $T_{\text {other }}$. In fact, by the assumption of the theorem, $S_{m}$ is also a minimum V-cut between $T_{1}$ and $T_{\text {other }}$. Let $C\left(T_{1}\right)$ be the set of vertices $x$ such that either $x \in T_{1}$ or there is a path between $x$ and $T_{1}$ in the subgraph $G-S_{m}$. In particular, since $u$ is not in $S_{m}$ and $u$ is adjacent to $T_{1}$, we have $u \in C\left(T_{1}\right)$. Moreover, let $C\left(T_{\text {other }}\right)=V-C\left(T_{1}\right)-S_{m}$.

By Lemma B.2, there exist $\left|S_{m}\right|$ internally disjoint paths between $T_{1}$ and $T_{\text {other }}$,


Fig. 3. Decomposition of separators
each contains exactly one vertex in the set $S_{m}$. Therefore, each of these $\left|S_{m}\right|$ paths is cut into two subpaths by a vertex in $S_{m}$, such that one subpath is in the induced subgraph $G\left(C\left(T_{1}\right)\right)$ and the another subpath is in the induced subgraph $G\left(C\left(T_{\text {other }}\right)\right)$. From this, we derive that there are $\left|S_{m}\right|$ internally disjoint paths between $T_{1}$ and $S_{m}$ in the induced subgraph $G\left(C\left(T_{1}\right) \cup S_{m}\right)$, each contains a distinct vertex in the set $S_{m}$.

Define $A=S_{k} \cap C\left(T_{1}\right), B=S_{k} \cap S_{m}$, and $C=S_{k} \cap C\left(T_{\text {other }}\right)$. Finally, let $S_{m}^{\prime}$ be the set of vertices $x$ in $S_{m}$ such that there is a path between $x$ and $T_{\text {other }}$ in the induced subgraph $G\left(C\left(T_{\text {other }}\right) \cup S_{m}-S_{k}\right)$ (see Figure 3 for an intuitive illustration of these sets).

We first prove that $|A| \geq\left|S_{m}^{\prime}\right|$.
From the fact that there are $\left|S_{m}\right|$ internally disjoint paths between $T_{1}$ and $S_{m}$ in the induced subgraph $G\left(C\left(T_{1}\right) \cup S_{m}\right)$ in which each path contains a distinct vertex in the set $S_{m}$, we derive that there are $\left|S_{m}^{\prime}\right|$ internally disjoint paths between $T_{1}$ and $S_{m}^{\prime}$ in the induced subgraph $G\left(C\left(T_{1}\right) \cup S_{m}^{\prime}\right)$. If $|A|<\left|S_{m}^{\prime}\right|$, then there must be a path $P_{1}$ between $T_{1}$ and a vertex $v^{\prime}$ in $S_{m}^{\prime}$ in the subgraph $G\left(C\left(T_{1}\right) \cup S_{m}^{\prime}-A\right)=$ $G\left(C\left(T_{1}\right) \cup S_{m}^{\prime}-S_{k}\right)$. Moreover, by the definition of the set $S_{m}^{\prime}$, there is also a path $P_{2}$ between $v^{\prime}$ and $T_{\text {other }}$ in the induced subgraph $G\left(C\left(T_{\text {other }}\right) \cup S_{m}-S_{k}\right)$. The concatenation of the paths $P_{1}$ and $P_{2}$ would give a path between $T_{1}$ and $T_{\text {other }}$ in the
induced subgraph $G\left(V-S_{k}\right)$, which contradicts the assumption that $S_{k}$ is a separator of the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. Therefore, we must have $|A| \geq\left|S_{m}^{\prime}\right|$.

Define a set $S_{k}^{\prime}=S_{m}^{\prime} \cup B \cup C$. We now prove that the set $S_{k}^{\prime}$ is a separator of the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$. Suppose that the set $S_{k}^{\prime}$ is not a separator of the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$, then there are two vertices $v_{1}$ and $v_{2}$ that are in two different terminal sets in $\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}$ and there exists a path $P$ between $v_{1}$ and $v_{2}$ in the induced subgraph $G\left(V-S_{k}^{\prime}\right)$. We discuss this in two possible cases.

Case 1: There is a vertex $w$ in the path $P$ such that $w \in C\left(T_{1}\right)$. Because (1) at least one of the vertices $v_{1}$ and $v_{2}$ is in the set $T_{\text {other }},(2)$ there is a path between $T_{1}$ and $w$ in the induced subgraph $G\left(C\left(T_{1}\right)\right)$, and (3) $S_{m}$ is a V-cut between $T_{1}$ and $T_{o t h e r}$, we conclude that there must be a vertex $s \in S_{m}$ that is also on the path $P$. Without loss of generality, we can suppose that the vertex $v_{1}$ is in the set $T_{\text {other }}$, and that the subpath $P^{\prime}$ of $P$ that begins from $v_{1}$ and ends at $s$ has no vertices from $C\left(T_{1}\right)$ - for this we only have to pick the first vertex $s$ in $S_{m}$ when we traverse on the path $P$ from $v_{1}$ to $v_{2}$. Then the path $P^{\prime}$ is in the induced subgraph $G\left(C\left(T_{o t h e r}\right) \cup S_{m}-S_{k}^{\prime}\right)$, which is a subgraph of the induced subgraph $G\left(C\left(T_{\text {other }}\right) \cup S_{m}-S_{k}\right)$. Now by the definition of the set $S_{m}^{\prime}$, the vertex $s$ is in the set $S_{m}^{\prime}$, thus in the set $S_{k}^{\prime}$. But this is impossible because we assumed that the path $P$ is in the induced subgraph $G\left(V-S_{k}^{\prime}\right)$.

Case 2: All vertices of the path $P$ are in $V-S_{k}^{\prime}-C\left(T_{1}\right)$. Then neither of the vertices $v_{1}$ and $v_{2}$ can be from the set $T_{1}$. Moreover, since the induced subgraph $G\left(V-S_{k}^{\prime}-C\left(T_{1}\right)\right)$ is a subgraph of the induced subgraph $G\left(V-S_{k}\right)$, the path $P$, which is between two different terminal sets in $\left\{T_{2}, \ldots, T_{l}\right\}$, would contain no vertex in $S_{k}$. But this again contradicts the assumption that $S_{k}$ is a separator of the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$.

Combining the discussions in Case 1 and Case 2, we conclude that the set $S_{k}^{\prime}$ is a separator for the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$.

Since $|A| \geq\left|S_{m}^{\prime}\right|, S_{k}=A \cup B \cup C$, and $S_{k}^{\prime}=S_{m}^{\prime} \cup B \cup C$, and $A$ does not intersect $B \cup C$, we conclude that $\left|S_{k}\right| \geq\left|S_{k}^{\prime}\right|$. In particular, if the instance ( $\left.G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has the separator $S_{k}$ of size bounded by $k$, then the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has the separator $S_{k}^{\prime}$ of size also bounded by $k$.

This completes the proof of the theorem.

The proof of Theorem C. 4 becomes complicated partially because the vertex $u$ may be included in a separator for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. If we restrict that the vertex $u$ is not in the separators for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$, then a result similar to Theorem C. 4 can be obtained much more easily, even without the need of the condition that the minimum V-cuts between $T_{1}$ and $\bigcup_{j \neq 1} T_{j}$ and the minimum V-cuts between $T_{1}^{\prime}$ and $\bigcup_{j \neq 1} T_{j}$ have the same size. This is given in the following lemma. This result will also be needed in our algorithm.

Lemma C. 3 Let $S$ be a vertex subset in the graph $G$ such that $S$ does not include the vertex $u$. Then $S$ is a separator for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ if and only if $S$ is a separator for the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$.

Proof. If $S$ is a separator for the instance ( $G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k$ ), then as explained in Theorem C.4, $S$ is also a separator for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$.

We are done once we show the other direction. Suppose that $S$ is a separator for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. We show that $S$ is also a separator for the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$. Suppose that $S$ is not a separator for the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$. Then there is a path $P$ in $G-S$ between two different terminal sets in $\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}$. Let one of these two terminal sets in $\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}$ be $T_{i}$, where $i \neq 1$. The path $P$ must contain the vertex $u$ (recall that $S$ does not contain $u$ ) - otherwise the path $P$ in $G-S$ would be between two different termi-
nal sets in $\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}$, contradicting the assumption that $S$ is a separator for $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. However, this would imply that the path from $T_{1}$ to $u$ (recall that $u$ has a neighbor in $T_{1}$ ) then following the path $P$ to the terminal set $T_{i}$ would give a path in $G-S$ between $T_{1}$ and $T_{i}$, again contradicting the assumption that $S$ is a separator for $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. This contradiction shows that the set $S$ must be also a separator for the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$.

Now, we are ready to present our algorithm. For an instance $\left(G,\left\{T_{1}, \ldots, T_{l}\right\}, k\right)$ of the parameterized node multiway cut problem, a vertex in the graph $G$ that does not belong to any terminal sets will be called a "non-terminal".

The algorithm is given in Figure 10.

Theorem C. 4 The algorithm $\operatorname{NMC}\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ in Figure 10 solves the PARAMETERIZED NODE MULTIWAY CUT problem in time $O\left(n^{3} k 4^{k}\right)$.

Proof. We first prove the correctness of the algorithm. Let $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ be an input to the algorithm, which is an instance of the PARAMETERIZED NODE MULTIWAY CUT problem, where $G=(V, E)$ is a graph, $\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}$ is a collection of terminal sets, and $k$ is the upper bound of the size of the separator we are looking for.

If there is an edge whose two ends are in two different terminal sets, then we have no way to separate these two terminal sets since all vertices in a separator are supposed to be non-terminals. Step 1 handles this case correctly.

If a non-terminal $w$ has two neighbors that are in two different terminal sets, then $w$ must be in the separator because otherwise the two terminal sets will not be separated. Thus, we can simply include the vertex $w$ in the separator, and recursively find a separator of size bounded by $k-1$ for the same collection of terminal sets

Algorithm NMC $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$
input: an instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ of the Parameterized node mULTIWAY CUT problem $(l \geq 2)$
output: a separator of size bounded by $k$ for $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$, or report "No" (i.e., no such a separator)

1. if an edge has its two ends in two different terminal sets then return "No";
2. if a non-terminal $w$ has two neighbors in two different terminal sets then return $w+\mathbf{N M C}\left(G-w,\left\{T_{1}, \ldots, T_{l}\right\}, k-1\right)$;
3. find the size $m_{1}$ of a minimum V-cut between $T_{1}$ and $\bigcup_{j=2}^{l} T_{j}$;
4. if $m_{1}>k$ then return "No";
5. if ( $m_{1}=0$ and $l=2$ ) then return $\emptyset$;
5.1 if ( $m_{1}=0$ and $l>2$ ) then return $\operatorname{NMC}\left(G,\left\{T_{2}, \ldots, T_{l}\right\}, k\right)$;
6. else pick a non-terminal $u$ that has a neighbor in $T_{1}$; let $T_{1}^{\prime}=T_{1}+u$;
6.1 if the size of a minimum V-cut between $T_{1}^{\prime}$ and $\bigcup_{j=2}^{l} T_{j}$ is equal to $m_{1}$ then return $\operatorname{NMC}\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$;
6.2 else $S=u+\operatorname{NMC}\left(G-u,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k-1\right)$;
if $S$ is not "No" then return $S$;
6.3 else return $\operatorname{NMC}\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$.

Fig. 4. An algorithm for the Parameterized node multiway cut problem
$\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}$ in the remaining graph $G-w$. This case is correctly handled by step $2 .{ }^{2}$

Step 3 computes the size $m_{1}$ of a minimum V-cut between the sets $T_{1}$ and $\bigcup_{j=2}^{l} T_{j}$. By Lemma B.3, the value $m_{1}$ can be computed in time $O((|V|+|E|) k)$.

If $m_{1}>k$, then the size of a minimum V -cut between $T_{1}$ and $\bigcup_{j=2}^{l} T_{j}$ is larger than $k$, which means that even separating the set $T_{1}$ from the other sets $\bigcup_{j=2}^{l} T_{j}$

[^2]requires more than $k$ vertices. Thus, no separator of size bounded by $k$ can exist for the terminal sets $T_{1}, T_{2}, \ldots, T_{l}$. This is handled by step 4 .

In step 5 we handle the case $m_{1}=0$ and $l=2$, which we do not need to remove any vertex to separate $T_{1}$ and $T_{2}$, i.e. the problem is solved. So we return an empty set $\emptyset$ as a separator of size 0 (note that because of step 4 , here we must have $k \geq 0$ ). In step $5.1, m_{1}=0$ and $l>2$, which means that $T_{1}$ is already separated from $T_{2}, \ldots, T_{l}$. Hence we only need to find a separator to separate $T_{2}, \ldots, T_{l}$. Therefore, step 5.1 handles this case correctly.

When the algorithm reaches step 6 , the following conditions hold true: (1) no edge has its two ends in two different terminal sets (because of step 1); (2) no nonterminal has two neighbors in two different terminal sets (because of step 2); (3) $0<m_{1} \leq k$ (because of steps 4-5). In particular, by condition (3), there must be a non-terminal $u$ that has a neighbor in $T_{1}$.

Let $m^{\prime}$ be the size of a minimum V-cut between the sets $T_{1}^{\prime}$ and $\bigcup_{j=2}^{l} T_{j}$. If $m^{\prime}=m_{1}$, then by Theorem C.4, the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ has a separator of size bounded by $k$ if and only if the instance ( $G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k$ ) has a separator of size bounded by $k$. In particular, as shown in the proof of Theorem C.4, a separator of size bounded by $k$ for the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$ is actually also a separator for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. Therefore, in this case, we can recursively work on the instance $\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$, as given in step 6.1. On the other hand, if $m^{\prime} \neq m_{1}$, which means $m^{\prime}>m_{1}$, then we simply branch on the vertex $u$ in two cases: (1) including $u$ in the separator and recursively working on the remaining graph for a separator of size bounded by $k-1$, as given by step 6.2 ; and (2) excluding $u$ from the separator thus looking for a separator that does not include $u$ and is of size bounded by $k$ for the instance $\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$. By Lemma C.1, the second case is equivalent to finding a separator of size bounded by $k$ for the instance
$\left(G,\left\{T_{1}^{\prime}, T_{2}, \ldots, T_{l}\right\}, k\right)$. This case is thus handled by step 6.3.
This completes the proof of the correctness of the algorithm. Now we analyze the complexity of the algorithm.

The recursive execution of the algorithm can be described as a search tree $\mathcal{T}$. We first count the number of leaves in the search tree $\mathcal{T}$. Note that only steps 6.2-6.3 of the algorithm correspond to branches in the search tree $\mathcal{T}$. Let $D\left(k, m_{1}\right)$ be the total number of leaves in the search tree $\mathcal{T}$ for the algorithm $\operatorname{NMC}\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$, where $m_{1}$ is the size of a minimum V-cut between the sets $T_{1}$ and $\bigcup_{j=2}^{l} T_{j}$. Then steps 6.2-6.3 induce the following recurrence relation:

$$
\begin{equation*}
D\left(k, m_{1}\right) \leq D\left(k-1, m^{\prime \prime}\right)+D\left(k, m^{\prime \prime \prime}\right) \tag{3.1}
\end{equation*}
$$

where $m^{\prime \prime}$ is the size of a minimum V-cut between $T_{1}$ and $\bigcup_{j=2}^{l} T_{j}$ in the graph $G-u$ as given in step 6.2, and $m^{\prime \prime \prime}$ is the size of a minimum V-cut between $T_{1}^{\prime}$ and $\bigcup_{j=2}^{l} T_{j}$ in the graph $G$ as given in step 6.3. Note that $m_{1}-1 \leq m^{\prime \prime} \leq m_{1}$ because removing the vertex $u$ from $G$ cannot increase the size of a minimum V-cut between two sets, and can decrease the size of a minimum V-cut between the two sets by at most 1 . Moreover, by Lemma C. 1 and because of step 6.1, the size $m^{\prime \prime \prime}$ of a minimum V-cut between $T_{1}^{\prime}$ and $\bigcup_{j=2}^{l} T_{j}$ in step 6.3 is at least $m_{1}+1$. Summarizing these, we have

$$
\begin{equation*}
m_{1}-1 \leq m^{\prime \prime} \leq m_{1} \quad \text { and } \quad m^{\prime \prime \prime} \geq m_{1}+1 \tag{3.2}
\end{equation*}
$$

Introduce a new function $D^{\prime}$ such that $D^{\prime}\left(2 k-m_{1}\right)=D\left(k, m_{1}\right)$, and let $t=$ $2 k-m_{1}$. Then by Inequalities (5.1) and (5.2), the branch in step 6.2-6.3 in the algorithm becomes

$$
D^{\prime}(t) \leq D^{\prime}\left(t_{1}\right)+D^{\prime}\left(t_{2}\right)
$$

where when $t=2 k-m_{1}$ then $t_{1}=2(k-1)-m^{\prime \prime} \leq t-1$, and $t_{2}=2 k-m^{\prime \prime \prime} \leq t-1$.

We also point out that certain non-branching steps (i.e., steps 2, 5.1, and 6.1) may also change the values of $k$ and $m_{1}$, thus changing the value $t=2 k-m_{1}$. However, none of these steps increases the value $t=2 k-m_{1}$ : (1) step 2 decreases the value $k$ by 1 and the value $m_{1}$ by at most 1 , which as a total will decrease the value $t=2 k-m_{1}$ by at least 1 ; (2) step 5.1 keeps the value $k$ unchanged and, since we have $m_{1}=0$ before the execution of this step, the new value $m_{1}$ is at least as large as the old value $m_{1}$. As a consequence, the value $t=2 k-m_{1}$ is not increased; (3) finally, step 6.1 does not change the values $k$ and $m_{1}$, thus neither changes the value $t=2 k-m_{1}$. In summary, the value $t=2 k-m_{1}$ after a branching step to the next branching step can never be increased.

Our initial instance starts with $t=2 k-m_{1} \leq 2 k$. In the case $t=2 k-m_{1}=0$, because we also have the conditions $k \geq m_{1} \geq 0$, we must have $m_{1}=0$ and $k=0$, in this case the algorithm can solve the instance without further branching. Therefore, we have $D^{\prime}(0)=1$. Combining all these, we derive

$$
D\left(k, m_{1}\right)=D^{\prime}\left(2 k-m_{1}\right) \leq 2^{2 k}
$$

and the search tree $\mathcal{T}$ has at most $2^{2 k}$ leaves.
Finally, it is easy to verify that along each root-leaf path in the search tree $\mathcal{T}$, the running time of the algorithm is bounded by $O\left(n^{3} k\right)$, where $n$ is the number of vertices in the graph. In conclusion, the running time of the algorithm $\operatorname{NMC}\left(G,\left\{T_{1}, T_{2}, \ldots, T_{l}\right\}, k\right)$ is bounded by $O\left(n^{3} k 4^{k}\right)$.

This completes the proof of the theorem.

## D. Final Remarks

We developed new and powerful techniques that lead to an algorithm of running time $O\left(n^{3} k 4^{k}\right)$ for a generalized version of the PARAMETERIZED NODE MULTIWAY CUT problem. The algorithm significantly improves previous algorithms for the problem. More recently, our techniques have been extended to directed graphs that lead to a fixed parameter tractable algorithm for the FEEDBACK VERTEX SET problem [24], thus resolving an outstanding open problem in parameterized computation and complexity.

Our algorithm finds a separator that has no vertices in any terminal set. We call such a separator a restricted separator. If a separator is allowed to include vertices from terminal sets, the separator is called an unrestricted separator. It can be verified easily that the instance $\left(G,\left\{T_{1}, \ldots, T_{l}\right\}, k\right)$ has an unrestricted separator of size $k$ if and only if the instance $\left(G^{\prime},\left\{\left\{x_{1}\right\}, \ldots,\left\{x_{l}\right\}\right\}, k\right)$ has a restricted separator of size $k$, where the graph $G^{\prime}$ is obtained from the graph $G$ by adding $l$ new vertices $x_{1}, \ldots, x_{l}$ and connecting $x_{i}$ to each vertex in $T_{i}$ for all $1 \leq i \leq l$. Therefore, our algorithm can also be used to construct unrestricted separators for undirected graphs.

One related problem is the Parameterized node multicut problem [82], where we look for a separator of size $k$ to separate each of the $l$ given pairs of terminals. When both $k$ and $l$ are used as parameters, based on the techniques developed in the current chapter, the fixed parameter tractable algorithm presented in [82] for the PARAMETERIZED NODE MULTICUT problem can be improved. On the other hand, if only $k$ is used as the parameter, or if the graph $G$ is a directed graph (or even just a directed acyclic graph), it is currently unknown whether the PARAMETERIZED NODE MULTICUT problem has fixed parameter tractable algorithms, which seem very interesting topics for further research.

## CHAPTER IV

## UNDIRECTED FEEDBACK VERTEX SET*

In this chapter, we give an fpt-algorithm of running time $O^{*}\left(5^{k}\right)$ for the FEEDBACK VERTEX SET problem in weighted graphs. We first present an fpt-algorithm for the FEEDBACK VERTEX SET problem in unweighted graphs, then extend that algorithm for the FEEDBACK VERTEX SET problem in weighted graphs.

The previous algorithms have been focused on the parameter $k$, the size of a feedback vertex set to search, as the only measure. The previous algorithm takes the approach based on the iterative compression method [90]. Our algorithm still takes the same approach. But we consider both the parameter $k$ and the number of connected components when we apply the iterative compression method. We discover an operation which either decreases $k$ by 1 or decreases the number of connected components by 1 . When either $k$ become 0 or the number of connected component is 1 , the FEEDBACK VERTEX SET problem can be solved in polynomial time. With these interesting properties, we design an algorithm of running time $O^{*}\left(5^{k}\right)$ for the FEEDBACK VERTEX SET problem in weighted graphs.

## A. Introduction

Let $G$ be an undirected graph. A feedback vertex set (FVS) $F$ in $G$ is a set of vertices in $G$ whose removal results in an acyclic graph (or equivalently, every cycle in $G$ contains at least one vertex in $F$ ). The problem of finding a minimum feedback vertex set in a graph is one of the classical NP-complete problems [68] and has many

[^3]applications. The history of the problem can be traced back to early '60s. For several decades, many different algorithmic approaches were tried on this problem, including approximation algorithms, linear programming, local search, polyhedral combinatorics, and probabilistic algorithms (see the survey of Festa et al. [47]). There are also exact algorithms finding a minimum FVS in a graph of $n$ vertices in time $\mathcal{O}\left(1.9053^{n}\right)$ [89] and in time $\mathcal{O}\left(1.7548^{n}\right)$ [48].

An important application of the FVS problem is deadlock recovery in operating systems [94], in which a deadlock is presented by a cycle in a system resource-allocation graph $G$. Therefore, in order to recover from deadlocks, we need to abort a set of processes in the system, i.e., to remove a set of vertices in the graph $G$, so that all cycles in $G$ are broken. Equivalently, we need to find an FVS in $G$. The problem also has a version on weighted graphs, where the weight of a vertex can be interpreted as the cost of aborting the corresponding process. In this case, we are looking for an FVS in $G$ whose weight is minimized.

In a practical system resource-allocation graph $G$, it can be expected that the size $k$ of the minimum FVS in $G$, i.e., the number of vertices in the FVS, is fairly small. This motivated the study of parameterized algorithms for the FVS problem that find an FVS of $k$ vertices in a graph of $n$ vertices (where the weight of the FVS is minimized, in the case of weighted graphs), and run in time $f(k) n^{\mathcal{O}(1)}$ for a fixed function $f$ (thus, the algorithms become practically efficient when the value $k$ is small).

This line of research has received considerable attention, mostly on the unweighted version of the problem. The first group of parameterized algorithms of running time $f(k) n^{\mathcal{O}(1)}$ for the FVS problem on unweighted graphs was given by Bodlaender [10] and by Downey and Fellows [35]. Since then a chain of dramatic improvements was obtained by different researchers (see table II for references).

Table II. History of parameterized algorithms for the UNWEIGHTED FEEDBACK VERTEX SET problem

| Bodlaender, Fellows [10, 35] | $\mathcal{O}\left(17\left(k^{4}\right)!n^{\mathcal{O}(1)}\right)$ |
| :--- | :---: |
| Downey and Fellows [37] | $\mathcal{O}\left((2 k+1)^{k} n^{2}\right)$ |
| Raman et al.[87] | $\mathcal{O}\left(\max \left\{12^{k},(4 \log k)^{k}\right\} n^{2.376}\right)$ |
| Kanj et al.[65] | $\mathcal{O}\left((2 \log k+2 \log \log k+18)^{k} n^{2}\right)$ |
| Raman et al.[88] | $\mathcal{O}\left((12 \log k / \log \log k+6)^{k} n^{2.376}\right)$ |
| Guo et al.[57] | $\mathcal{O}\left((37.7)^{k} n^{2}\right)$ |
| Dehne et al.[33] | $\mathcal{O}\left((10.6)^{k} n^{3}\right)$ |

Randomized parameterized algorithms have also been studied in the literature for the FVS problem, for both unweighted and weighted graphs. The best known randomized parameterized algorithms for the FVS problems are due to Becker et al. [6], who developed a randomized algorithm of running time $\mathcal{O}\left(4^{k} k n^{2}\right)$ for the FVS problem on unweighted graphs, and a randomized algorithm of running time $\mathcal{O}\left(6^{k} k n^{2}\right)$ for the FVS problem on weighted graphs. To our knowledge, no deterministic algorithm of running time $f(k) n^{O(1)}$ for any function $f$ was known prior to our results for the weighted FVS problem.

The main result of this chapter is an algorithm that for a given integer $k$ and a weighted graph $G$, either finds a minimum weight FVS in $G$ of at most $k$ vertices, or correctly reports that $G$ contains no FVS of at most $k$ vertices. The running time of our algorithm is $\mathcal{O}\left(5^{k} k n^{2}\right)$. This improves and generalizes a long chain of results in parameterized algorithms. Let us remark that the running time of our (deterministic) algorithm comes close to that of the best randomized algorithm for the FVS problem on unweighted graphs and is better than the running time of the
previous best randomized algorithm for the FVS problem on weighted graphs.
The general approach of our algorithm is based on the iterative compression method [90], which has been successfully used recently for improved algorithms for the FVS and other problems [33, 57, 90]. The method starts with an FVS of $k$ vertices for a small subgraph of the given graph, and iteratively grows the small subgraph while keeping an FVS of $k$ vertices in the grown subgraph until the subgraph becomes the original input graph. This method makes it possible to reduce the original FVS problem on general graphs to the FVS problem on graphs with a special decomposition structure. The main contribution of the current chapter is the development of a general algorithmic technique that identifies a dual parameter in problem instances that limits the number of times where the original parameter $k$ cannot be effectively reduced during a branch and search process. In particular, for the FVS problem on graphs of the above special decomposition structure, a measure is introduced that nicely combines the original parameter and the dual parameter and bounds effectively the running time of a branch and search algorithm for the FVS problem. This technique leads to a simpler but significantly more efficient parameterized algorithm for the FVS problem on unweighted graphs. Moreover, the introduction of the measure greatly simplifies the process of degree-2 vertices in a weighted graph, and enables us to solve the FVS problem on weighted graphs in the same complexity as that for the problem on unweighted graphs. Note that this is significant because no previous algorithms for the FVS problem on unweighted graphs can be extended to weighted graphs mainly because of the lack of effective method for handling degree-2 vertices. Finally, the technique of dual parameters seems to be of general usefulness for the development of parameterized algorithms, and has been used more recently in solving other parameterized problems $[23,24]$.

The remaining part of this chapter is organized as follows. In Section B, we provide in full details a simpler algorithm and its analysis for unweighted graphs. This is done for clearer demonstration of our approach. We also indicate why this simpler algorithm does not work for weighted graphs. In Section C, we obtain the main result of the chapter, the algorithm for the weighted FVS problem. This generalization of the results from Section B is not straightforward and requires a number of new structures and techniques.

## B. Feedback Vertex Set in Unweighted Graphs

In this section, we consider the FVS problem on unweighted graphs. We start with some terminologies. A forest is a graph that contains no cycles. A tree is a forest that is connected (therefore, a forest can be equivalently defined as a collection of disjoint trees). Let $W$ be a subset of vertices in a graph $G=(V, E)$. We will denote by $G[W]$ the subgraph of $G$ that is induced by the vertex set $W$. For simplicity we will use the notation $G-w$ and $G-W$ for respectively $G[V \backslash\{w\}]$ and $G[V \backslash W]$ where $w \in V$ and $W \subseteq V$. A pair $\left(V_{1}, V_{2}\right)$ of vertex subsets in a graph $G=(V, E)$ is a forest bipartition of $G$ if $V_{1} \cup V_{2}=V, V_{1} \cap V_{2}=\emptyset$, and both induced subgraphs $G\left[V_{1}\right]$ and $G\left[V_{2}\right]$ are forests. For a vertex $u \in V$ the degree of $u$ will be the number of edges incident to $u$.

Let $G$ be a graph and let $F$ be a subset of vertices in $G$. The set $F$ is a feedback vertex set (shortly, FVS) of $G$ if $G-F$ is a forest. The size of an FVS $F$ is the number of vertices in $F$.

Our main problem is formally defined as follows.

FEEDBACK VERTEX SET: given a graph $G$ and an integer $k$, either find an FVS of size at most $k$ in $G$, or report that no such an FVS exists.

Before we present our algorithm for the FEEDBACK VERTEX SET problem, we first consider a special version of the problem, defined as follows:

F-BIPARTITION FVS: given a graph $G$, a forest bipartition $\left(V_{1}, V_{2}\right)$ of $G$, and an integer $k$, either find an FVS of size at most $k$ for the graph $G$ in the subset $V_{1}$, or report that no such an FVS exists.

Note that the main difference between the F-bipartition fVs problem and the original FEEDBACK VERTEX SET problem is that we require that the FVS in the F-BIPARTITION FVS is contained in the given subset $V_{1}$.

Observe that certain structures in the input graph $G$ can be easily processed and then removed from $G$. For example, if a vertex $v$ has a self-loop (i.e., an edge whose both ends are incident to $v$ ), then the vertex $v$ is necessarily contained in every FVS in $G$. Thus, we can directly include $v$ in the objective FVS. If two vertices $v$ and $w$ are connected by multiple edges (i.e., there are more than one edge whose one end is $v$ and the other end is $w$ ), then one of $v$ and $w$ must be contained in the objective FVS. Thus, we can branch into two recursive calls, one includes $v$, and the other includes $w$, in the objective FVS. All these operations are more efficient than the algorithm of running time $\mathcal{O}\left(5^{k} k n^{2}\right)$ developed in the current chapter. Therefore, for a given input graph $G$, we always first apply a preprocessing that applies the above operations and remove all self-loops and multiple edges in the graph $G$. In consequence, we can assume, without loss of generality, that the input graph $G$ contains neither self-loops nor multiple edges.

The algorithm, Feedback $\left(G, V_{1}, V_{2}, k\right)$, for the F-bipartition FVs problem is given in Figure 5. We first discuss the correctness of the algorithm. The correctness of step 1 and step 2 of the algorithm is obvious. Now consider step 3. Let $w$ be a vertex in $V_{1}$ that has at least two neighbors in $V_{2}$.

Algorithm-1 Feedback $\left(G, V_{1}, V_{2}, k\right)$
Input: $G=(V, E)$ is a graph with a forest bipartition $\left(V_{1}, V_{2}\right), k$ is an integer.
Output: An FVS $F$ of $G$ such that $|F| \leq k$ and $F \subseteq V_{1}$; or report "No" (i.e., no such an FVS exists).

1. if $(k<0)$ or ( $k=0$ and $G$ is not a forest) then return "No";
2. if $(k \geq 0)$ and $G$ is a forest then return $\emptyset$;
3. if a vertex $w$ in $V_{1}$ has at least two neighbors in $V_{2}$ then
3.1. if two neighbors of $w$ in $V_{2}$ is in the same tree in $G\left[V_{2}\right]$ then
$F_{1}=\operatorname{Feedback}\left(G-w, V_{1} \backslash\{w\}, V_{2}, k-1\right)$;
if $F_{1}=$ "No" then return "No"
else return $F_{1} \cup\{w\} ;$
3.2. else
$F_{1}=$ Feedback $\left(G-w, V_{1} \backslash\{w\}, V_{2}, k-1\right)$;
$F_{2}=$ Feedback $\left(G, V_{1} \backslash\{w\}, V_{2} \cup\{w\}, k\right)$;
if $F_{1} \neq$ "No" then return $F_{1} \cup\{w\}$
else return $F_{2}$;
4. else pick any vertex $w$ that has degree $\leq 1$ in $G\left[V_{1}\right]$;
4.1. if $w$ has degree $\leq 1$ in the original graph $G$ then return Feedback $\left(G-w, V_{1} \backslash\{w\}, V_{2}, k\right)$;
4.2. else return Feedback $\left(G, V_{1} \backslash\{w\}, V_{2} \cup\{w\}, k\right)$.

Fig. 5. Algorithm for the UNWEIGHTED FEEDBACK VERTEX SET problem.

If the vertex $w$ has two neighbors in $V_{2}$ that belong to the same tree $T$ in the induced subgraph $G\left[V_{2}\right]$, then the tree $T$ plus the vertex $w$ contains at least one cycle. Since our search for an FVS is restricted to $V_{1}$, the only way to break the cycles in $T \cup\{w\}$ is to include the vertex $w$ in the objective FVS. Moreover, the objective FVS of size at most $k$ exists in $G$ if and only if the remaining graph $G-w$ has an FVS of size at most $k-1$ in the subset $V_{1} \backslash\{w\}$ (note that $\left(V_{1} \backslash\{w\}, V_{2}\right)$ is a valid forest bipartition of the graph $G-w)$. Therefore, step 3.1 correctly handles this case.

If no two neighbors of the vertex $w$ belong to the same tree in the induced
subgraph $G\left[V_{2}\right]$, then the vertex $w$ is either in the objective FVS or not in the objective FVS. If $w$ is in the objective FVS, then we should be able to find an FVS $F_{1}$ in the graph $G-w$ such that $\left|F_{1}\right| \leq k-1$ and $F_{1} \subseteq V_{1} \backslash\{w\}$ (again note that ( $V_{1} \backslash\{w\}, V_{2}$ ) is a valid forest bipartition of the graph $G-w$ ). On the other hand, if $w$ is not in the objective FVS, then the objective FVS for $G$ must be contained in the subset $V_{1} \backslash\{w\}$. Also note that in this case, the subgraph $G\left[V_{2} \cup\{w\}\right]$ induced by the subset $V_{2} \cup\{w\}$ is still a forest since no two neighbors of $w$ in $V_{2}$ belong to the same tree in $G\left[V_{2}\right]$. In consequence, $\left(V_{1} \backslash\{w\}, V_{2} \cup\{w\}\right)$ still makes a valid forest bipartition for the graph $G$. Therefore, step 3.2 handles this case correctly.

Now we consider step 4. At this point, every vertex in $V_{1}$ has at most one neighbor in $V_{2}$. Moreover, since the induced subgraph $G\left[V_{1}\right]$ is a forest, there must be a vertex $w$ in $V_{1}$ that has degree at most 1 in $G\left[V_{1}\right]$ (note that $V_{1}$ cannot be empty at this point since otherwise the algorithm would have stopped at step 2). If the vertex $w$ also has degree at most 1 in the original graph $G$, then removing $w$ does not help breaking any cycles in $G$. Therefore, the vertex $w$ can be discarded. This case is correctly handled by step 4.1 . Otherwise, the vertex $w$ has degree at most 1 in the induced subgraph $G\left[V_{1}\right]$ but has degree larger than 1 in the original graph $G$. Observing that $w$ has at most one neighbor in $V_{2}$, we can derive that the degree of $w$ in the original graph $G$ must be exactly 2 . Moreover, $w$ has exactly two neighbors $u$ and $v$ such that $v$ is in $V_{1}$ and $u$ is in $V_{2}$.

Since the vertex $w$ has degree 2 in the original graph $G$, and the vertex $v$ is adjacent to $w$, we have that every cycle in $G$ that contains $w$ has to contain $v$. In consequence, if $w$ is contained in the objective FVS, then we can simply replace it by $v$. Therefore, in this case, we can safely assume that the vertex $w$ is not in the objective FVS. This can be easily implemented by moving the vertex $w$ from the set $V_{1}$ to the set $V_{2}$, and recursively working on the modified instance, as given in step
4.2 of the algorithm (note that $\left(V_{1} \backslash\{w\}, V_{2} \cup\{w\}\right)$ is a valid forest bipartition of the graph $G$, because by our assumption, the vertex $w$ will be a degree- 1 vertex in the induced subgraph $\left.G\left[V_{2} \cup\{w\}\right]\right)$.

Now we are ready to present the following lemma.
Lemma B. 1 The algorithm Feedback $\left(G, V_{1}, V_{2}, k\right)$ correctly solves the F-BIPARTITION FVS problem. The running time of the algorithm is $\mathcal{O}\left(2^{k+l} n^{2}\right)$, where $n$ is the number of vertices in $G$, and $l$ is the number of connected components in the induced subgraph $G\left[V_{2}\right]$.

Proof. The correctness of the algorithm has been verified by the above discussion. Now we consider the complexity of the algorithm.

The recursive execution of the algorithm can be described as a search tree $\mathcal{T}$. We first count the number of leaves in the search tree $\mathcal{T}$. Note that only step 3.2 of the algorithm corresponds to branches in the search tree $\mathcal{T}$. Let $T(k, l)$ be the total number of leaves in the search tree $\mathcal{T}$ for the algorithm Feedback $\left(G, V_{1}, V_{2}, k\right)$, where $l$ is the number of connected components (i.e., trees) in the forest $G\left[V_{2}\right]$. Inductively, the number of leaves in the search tree $\mathcal{T}_{1}$ corresponding to the recursive call Feedback $\left(G-w, V_{1} \backslash\{w\}, V_{2}, k-1\right)$ is at most $T(k-1, l)$. Moreover, we assumed at step 3.2 that $w$ has at least two neighbors in $V_{2}$ and that no two neighbors of $w$ in $V_{2}$ belong to the same tree in $G\left[V_{2}\right]$. Therefore, the vertex $w$ "merges" at least two trees in $G\left[V_{2}\right]$ into a single tree in $G\left[V_{2} \cup\{w\}\right]$. Hence, the number of trees in $G\left[V_{2} \cup\{w\}\right]$ is at most $l-1$. In consequence, the number of leaves in the search tree $\mathcal{I}_{2}$ corresponding to the recursive call Feedback $\left(G, V_{1} \backslash\{w\}, V_{2} \cup\{w\}, k\right)$ is at most $T(k, l-1)$. This gives the following recurrence relation:

$$
T(k, l) \leq T(k-1, l)+T(k, l-1)
$$

Also note that none of the (non-branching) recursive calls in the algorithm (steps 3.1, 4.1, and 4.2) would increase the values $k$ and $l$, and that $T(0, l)=1$ for all $l$ and $T(k, 0)=1$ for all $k$ (by steps 1-2). From all these facts, we can easily derive that $T(k, l)=\mathcal{O}\left(2^{k+l}\right)$.

Finally, observe that along each root-leaf path in the search tree $\mathcal{T}$, the total number of executions of steps $1,2,3,3.1,4.1$, and 4.2 of the algorithm is $\mathcal{O}(n)$ because each of these steps either stops immediately, or reduces the size of the set $V_{1}$ by at least 1 (and the size of $V_{1}$ is never increased during the execution of the algorithm). It remains to explain how each of the steps can be executed in $\mathcal{O}(n)$ time.

Before the first call to the Feedback algorithm, we use $\mathcal{O}\left(n^{2}\right)$ time, because this will happen only once. The three graphs, $G_{1}=G\left[V_{1}\right], G_{2}=G\left[V_{2}\right]$, and $G_{12}=\left(V, E \backslash\left(E\left(G_{1}\right) \cup E\left(G_{2}\right)\right)\right)$ can be trivially constructed in $\mathcal{O}\left(n^{2}\right)$ time. $G_{1}$ and $G_{2}$ are forests, and $G_{12}$ is a bipartite graph with the two vertex sets $V_{1}$ and $V_{2}$ as independent sets.

Steps 1, 2, 4.1, and 4.2 can be easily performed in $\mathcal{O}(n)$ time. For step 3, we simply search for a vertex of $V_{1}$ that has degree at least 2 in $G_{12}$, and for step 4 we search in $G_{1}$ for a leaf (vertex of degree at most 1). The condition for step 3.1 is that no two neighbors of $w$ belong to the same tree in $G\left[V_{2}\right]$, which can be checked by simply marking each neighbor of $w$, and doing a search in the forest $G\left[V_{2}\right]$.

Each of the steps $3.1,3.2$, 4.1, and 4.2 changes one or more of the graphs $G_{1}, G_{2}, G_{12}$, and we have to argue that these manipulations can also be done in $\mathcal{O}(n)$ time. Looking closely at these steps, we can observe that only two operations are required. The first is to delete a vertex in $V_{1}$, which corresponds to deleting the vertex and all incident edges in $G_{1}$ and $G_{12}$. The second operation is to move a vertex $w$ from $V_{1}$ to $V_{2}$, which corresponds to deleting $w$ from $G_{1}$ and updating $G_{2}$ and $G_{12}$
as follows: add $w$ to $V\left(G_{2}\right)$ and to $V_{2}$ of $G_{12}$, and read the set of edges incident to $w$ in $G$, and add edges between $w$ and vertices in $V_{2}$ to $G_{2}$ and between $w$ and $V_{1}$ to $G_{12}$. Using double linked lists and pointers it is possible to delete a vertex and all incident edges in $\mathcal{O}(n)$ time, and to insert edges in $\mathcal{O}(1)$ time.

Therefore, the computation time along each root-leaf path in the search tree $\mathcal{T}$ is $\mathcal{O}\left(n^{2}\right)$. In conclusion, the running time of the algorithm Feedback $\left(G, V_{1}, V_{2}, k\right)$ is $\mathcal{O}\left(2^{k+l} n^{2}\right)$. This completes the proof of the lemma.

Following the idea of iterative compression proposed by Reed et al. [90], we formulate the following problem:

FVS REDUCTION: given a graph $G$ and an FVS $F$ of size $k+1$ for $G$, either construct an FVS of size at most $k$ for $G$, or report that no such an FVS exists.

Lemma B. 2 The FVS Reduction problem on an n-vertex graph $G$ can be solved in time $\mathcal{O}\left(5^{k} n^{2}\right)$.

Proof. We use the algorithm Feedback to solve the fVS Reduction problem. Let $F$ be the FVS of size $k+1$ in the graph $G=(V, E)$. Every FVS $F^{\prime}$ of size at most $k$ for $G$ is a union of a subset $F_{1}$ of at most $k-j$ vertices in $V \backslash F$ and a subset $F_{2}$ of $j$ vertices in $F$, for some integer $j, 0 \leq j \leq k$. Note that since we assume that no vertex in $F \backslash F_{2}$ is in the FVS $F^{\prime}$, the induced subgraph $G\left[F \backslash F_{2}\right]$ must be a forest. For each $j, 0 \leq j \leq k$, we enumerate all subsets of $j$ vertices in $F$. For each such subset $F_{2}$ in $F$ such that $G\left[F \backslash F_{2}\right]$ is a forest, we seek a subset $F_{1}$ of at most $k-j$ vertices in $V \backslash F$ such that $F_{1} \cup F_{2}$ is an FVS in $G$.

Fix a subset $F_{2}$ in $F$, where $\left|F_{2}\right|=j$. Note that the graph $G$ has an FVS $F_{1} \cup F_{2}$ of size at most $k$, where $F_{1} \subseteq V \backslash F$, if and only if the subset $F_{1}$ of $V \backslash F$ is an FVS
for the graph $G-F_{2}$ and $\left|F_{1}\right| \leq k-j$. Therefore, to solve the original problem, we construct an FVS $F_{1}$ for the graph $G-F_{2}$ such that $\left|F_{1}\right| \leq k-j$ and $F_{1} \subseteq V \backslash F$.

Since $F$ is an FVS for $G$, we have that the induced subgraph $G[V \backslash F]=G-F$ is a forest. Moreover, by our assumption, the induced subgraph $G\left[F \backslash F_{2}\right]$ is also a forest. Note that $(V \backslash F) \cup\left(F \backslash F_{2}\right)=V \backslash F_{2}$, which is the vertex set for the graph $G^{\prime}=G-F_{2}$. Therefore, $\left(V \backslash F, F \backslash F_{2}\right)$ is a forest bipartition of the graph $G^{\prime}$. Thus, an FVS $F_{1}$ for the graph $G^{\prime}$ such that $\left|F_{1}\right| \leq k-j$ and $F_{1} \subseteq V \backslash F$ can be constructed by the algorithm Feedback $\left(G^{\prime}, V \backslash F, F \backslash F_{2}, k-j\right)$.

Since $|F|=k+1$ and $\left|F_{2}\right|=j$, we have that $\left|F \backslash F_{2}\right|=k+1-j$. Therefore, the forest $G\left[F \backslash F_{2}\right]$ contains at most $k+1-j$ connected components. By Lemma B.1, the running time of the algorithm Feedback $\left(G^{\prime}, V \backslash F, F \backslash F_{2}, k-j\right)$ is $\mathcal{O}\left(2^{(k-j)+(k+1-j)} n^{2}\right)=\mathcal{O}\left(4^{k-j} n^{2}\right)$. Now for all integers $j, 0 \leq j \leq k$, we enumerate all subsets $F_{2}$ of $j$ vertices in $F$ and apply the algorithm Feedback $\left(G^{\prime}, V \backslash F, F \backslash F_{2}, k-j\right)$ for those $F_{2}$ such that $G\left[F \backslash F_{2}\right]$ is a forest. As we discussed above, the graph $G$ has an FVS of size at most $k$ if and only if for some $F_{2} \subseteq F$, the algorithm Feedback $\left(G^{\prime}, V \backslash F, F \backslash F_{2}, k-j\right)$ produces an FVS $F_{1}$ for the graph $G^{\prime}$. The running time of this procedure is

$$
\sum_{j=0}^{k}\binom{k+1}{j} \cdot \mathcal{O}\left(4^{k-j} n^{2}\right)=\sum_{j=0}^{k}\binom{k+1}{k-j+1} \mathcal{O}\left(4^{k-j+1} n^{2}\right)=\mathcal{O}\left(5^{k} n^{2}\right)
$$

This completes the proof of the lemma.
Finally, by combining Lemma E. 1 with iterative compression, we obtain the main result of this section.

Theorem B. 3 The FEEDBACK VERTEX SET problem on an n-vertex graph is solvable in time $\mathcal{O}\left(5^{k} k n^{2}\right)$.

Proof. To solve the FEEDBACK vertex set problem, for a given graph $G=$ $(V, E)$, we start by applying Bafna et al.'s 2-approximation algorithm for the MINIMUM feedback vertex set problem [5]. This algorithm runs in $\mathcal{O}\left(n^{2}\right)$ time, and either returns an FVS $F^{\prime}$ of size at most $2 k$, or verifies that no FVS of size at most $k$ exists. If no FVS is returned, the algorithm is terminated with the conclusion that no FVS of size at most $k$ exists. In the case of the opposite result, we use any subset $V^{\prime}$ of $k$ vertices in $F^{\prime}$, and put $V_{0}=V^{\prime} \cup\left(V \backslash F^{\prime}\right)$. Of course, the induced subgraph $G\left[V_{0}\right]$ has an FVS of size $k$, namely the set $V^{\prime}\left(G\left[V_{0} \backslash V^{\prime}\right]\right.$ is a forest). Let $F^{\prime} \backslash V^{\prime}=\left\{v_{1}, v_{2}, \ldots, v_{\left|F^{\prime}\right|-k}\right\}$, and let $V_{i}=V_{0} \cup\left\{v_{1}, \ldots, v_{i}\right\}$ for $i \in\left\{0,1, \ldots,\left|F^{\prime}\right|-k\right\}$. Inductively, suppose that we have constructed an FVS $F_{i}$ for the graph $G\left[V_{i}\right]$, where $\left|F_{i}\right|=k$. Then the set $F_{i+1}^{\prime}=F_{i} \cup\left\{v_{i+1}\right\}$ is obviously an FVS for the graph $G\left[V_{i+1}\right]$ and $\left|F_{i+1}^{\prime}\right|=k+1$.

Now the pair $\left(G\left[V_{i+1}\right], F_{i+1}^{\prime}\right)$ is an instance for the FVS REDUCTION problem. Therefore, in time $\mathcal{O}\left(5^{k} n^{2}\right)$, we can either construct an FVS $F_{i+1}$ of size $k$ for the graph $G\left[V_{i+1}\right]$, or report that no such an FVS exists. Note that if the graph $G\left[V_{i+1}\right]$ does not have an FVS of size $k$, then the original graph $G$ cannot have an FVS of size $k$. In this case, we simply stop and claim the non-existence of an FVS of size $k$ for the original graph $G$. On the other hand, with an FVS $F_{i+1}$ of size $k$ for the graph $G\left[V_{i+1}\right]$, our induction proceeds to the next graph $G\left[V_{i+1}\right]$, until we reach the graph $G=G\left[V_{\left|F^{\prime}\right|-k}\right]$. This process runs in time $\mathcal{O}\left(5^{k} k n^{2}\right)$ since $\left|F^{\prime}\right|-k \leq k$, and solves the feedback vertex set problem.

## C. Feedback Vertex Set in Weighted Graphs

In this section, we discuss the feedback vertex set problem on weighted graphs. A weighted graph $G=(V, E)$ is an undirected graph, where each vertex $u \in V$ is
assigned a weight that is a positive real number. The weight of a vertex set $A \subseteq V$ is the sum of the vertex weights of all vertices in $A$. We denote by $|A|$ the cardinality of $A$. The (parameterized) FEEDBACK VERTEX SET problem on weighted graphs is formally defined as follows:

WEIGHTED-FVS: given a weighted graph $G$ and an integer $k$, either find an FVS $F$ of minimum weight for $G$ such that $|F| \leq k$, or report that no FVS of size at most $k$ exists in $G$.

The algorithm for the weighted case has several similarities with the unweighted case, but has also a significant difference. The difference is that step 4.2 of

Algorithm-1 becomes invalid for weighted graphs. A degree- 2 vertex $w$ in the set $V_{1}$ cannot simply be excluded from the objective FVS. If the weight of $w$ is smaller than that of its parent $v$ in $G\left[V_{1}\right]$, it may become necessary to include $w$ instead of $v$ in the objective FVS.

To overcome this difficulty, we introduce a new partition structure of the vertices in a weighted graph.

Definition A triple ( $V_{0}, V_{1}, V_{2}$ ) is an independent-forest partition (IF-partition) of a graph $G=(V, E)$ if $\left(V_{0}, V_{1}, V_{2}\right)$ is a partitioning of $V$ (i.e., $V_{0} \cup V_{1} \cup V_{2}=V$, and $V_{0}$, $V_{1}$, and $V_{2}$ are pairwise disjoint), such that
(1) $G\left[V_{1}\right]$ and $G\left[V_{2}\right]$ are forests;
(2) $G\left[V_{0}\right]$ is an independent set;
(3) Every vertex $u$ in $V_{0}$ is of degree 2 in $G$, with both neighbors in $V_{2}$.

The following problem is the analogue of the F-BIPARTITION FVS problem on weighted graphs.

WEIGHTED IF-PARTITION FVS: given a weighted graph $G$, an IF-partition $\left(V_{0}, V_{1}, V_{2}\right)$ of $G$, and an integer $k$, either find an FVS $F$ of minimum weight for $G$ that satisfies the conditions $|F| \leq k$ and $F \subseteq V_{0} \cup V_{1}$, or report that no such an FVS exists.

To develop and analyze our algorithm for the WEIGHTED IF-PARTITION FVS problem, we need the following concept of measure for the problem instances. For a vertex subset $V^{\prime}$ in the graph $G$, we will denote by $\# c\left(V^{\prime}\right)$ the number of connected components in the induced subgraph $G\left[V^{\prime}\right]$.

Definition Let $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ be an instance of the WEIGHTED IF-PARTITION FVS problem with an IF-partition $\left(V_{0}, V_{1}, V_{2}\right)$. The deficiency of the instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ is defined as

$$
\tau\left(k, V_{0}, V_{1}, V_{2}\right)=k-\left(\left|V_{0}\right|-\# c\left(V_{2}\right)+1\right)
$$

Intuitively, $\tau\left(k, V_{0}, V_{1}, V_{2}\right)$ of the instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ is an upper bound on the number of vertices in the objective FVS that are in the set $V_{1}$ (this will become clearer during our discussion below). Our algorithm for the WEIGHTED IF-PARTITION FVS problem is based on the following observation: once we have correctly determined all vertices in the objective FVS that are in the set $V_{1}$, the problem will become solvable in polynomial time, as shown in the following lemma.

Lemma C. 1 Let $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ be an instance of the Weighted if-Partition FVS problem with an IF-partition $\left(V_{0}, V_{1}, V_{2}\right)$ of an n-vertex graph $G$. If $V_{1}=\emptyset$, or $V_{2}=\emptyset$, or $\tau\left(k, V_{0}, V_{1}, V_{2}\right) \leq 0$, then a solution to the instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ can be constructed in time $\mathcal{O}\left(n^{2}\right)$.

Proof. First of all, note that if $k<0$, then the solution to the instance is "No": we cannot remove a negative number of vertices from $G$. Thus, in the following discussion, we assume that $k \geq 0$.

If $V_{2}=\emptyset$, then by the definition, $V_{0}$ should also be an empty set. Thus, the graph $G=G\left[V_{1}\right]$ is a forest, and the solution to the instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ is the empty set $\emptyset$.

Now consider the case $V_{1}=\emptyset$. Then we need to find a minimum-weight subset of at most $k$ vertices in the set $V_{0}$ whose removal from the graph $G=G\left[V_{0} \cup V_{2}\right]$ results in a forest.

Construct a new graph $\mathcal{H}=(\mathcal{V}, \mathcal{E})$, where each vertex $\mu$ in $\mathcal{V}$ corresponds to a connected component in the induced subgraph $G\left[V_{2}\right]$, and each edge $[\mu, \nu]$ in $\mathcal{E}$ corresponds to a vertex $v$ in the set $V_{0}$ such that the two neighbors of $v$ are in the connected components in $G\left[V_{2}\right]$ that correspond to the two vertices $\mu$ and $\nu$, respectively, in $\mathcal{H}$. Intuitively, the graph $\mathcal{H}$ can be obtained from the graph $G=G\left[V_{0} \cup V_{2}\right]$ by "shrinking" each connected component in $G\left[V_{2}\right]$ into a single vertex and "smoothing" each degree- 2 vertex in $V_{0}$ (note that the graph $\mathcal{H}$ may contain multiple edges and self-loops). Moreover, we give each edge in $\mathcal{H}$ a weight that is equal to the weight of the corresponding vertex in $V_{0}$. Thus, the graph $\mathcal{H}$ is a graph with edge weights. Observe that there is a one-to-one correspondence between the connected components in the graph $\mathcal{H}$ and the connected components in the graph
$G$. Moreover, since each connected component in the induced subgraph $G\left[V_{2}\right]$ is a tree, a connected component in the graph $\mathcal{H}$ is a tree if and only if the corresponding connected component in the graph $G$ is a tree. Most importantly, removing a vertex in $V_{0}$ in the graph $G$ corresponds to removing the corresponding edge in the graph $\mathcal{H}$. Therefore, the problem of constructing a minimum-weight vertex set in $V_{0}$ whose removal from $G$ results in a forest is equivalent to the problem of constructing a minimum-weight edge set in the graph $\mathcal{H}$ whose removal from $\mathcal{H}$ results in a forest.

Let $\mathcal{H}_{1}, \ldots, \mathcal{H}_{s}$ be the connected components of the graph $\mathcal{H}$, where for each $i$, the component $\mathcal{H}_{i}$ has $n_{i}$ vertices and $m_{i}$ edges. An edge set $\mathcal{E}_{i}$ in $\mathcal{H}_{i}$ whose removal from $\mathcal{H}_{i}$ results in a forest is of the minimum weight if and only if the complement graph $\mathcal{H}_{i}-\mathcal{E}_{i}$ is a spanning tree of the maximum weight in $\mathcal{H}_{i}$. Thus, the union $\mathcal{E}^{\prime}=\bigcup_{i=1}^{s}\left(\mathcal{H}_{i}-\mathcal{T}_{i}\right)$ is a minimum-weight edge set whose removal from $\mathcal{H}$ results in a forest, where for each $i, \mathcal{T}_{i}$ is a maximum-weight spanning tree in the graph $\mathcal{H}_{i}$. Since the maximum-weight spanning tree $\mathcal{T}_{i}$ in $\mathcal{H}_{i}$ can be constructed in time
$\mathcal{O}\left(n_{i}^{2}\right)$ by modifying the well-known minimum spanning tree algorithms (the algorithms work even for graphs with self-loops and multiple edges) [28], we conclude that the minimum-weight edge set $\mathcal{E}^{\prime}$ in $\mathcal{H}$ can be constructed in time $\sum_{i=1}^{s} \mathcal{O}\left(n_{i}^{2}\right)=\mathcal{O}\left(n^{2}\right)$. Also note that the number of edges in the set $\mathcal{E}^{\prime}$ is equal to $\sum_{i=1}^{s}\left(m_{i}-n_{i}+1\right)=|\mathcal{E}|-|\mathcal{V}|+s$.

Correspondingly, in case $V_{1}=\emptyset$, each minimum-weight FVS in $V_{0}$ for the graph $G$ contains exactly $|\mathcal{E}|-|\mathcal{V}|+s$ vertices, and such an FVS can be constructed in time $\mathcal{O}\left(n^{2}\right)$. Note that $|\mathcal{E}|=\left|V_{0}\right|,|\mathcal{V}|$ is equal to the number $\# c\left(V_{2}\right)$ of connected components in the induced subgraph $G\left[V_{2}\right]$, and $s$ (which is the number of connected components in $\mathcal{H})$ is equal to the number $\# c(G)$ of connected components in the graph $G=G\left[V_{0} \cup V_{2}\right]$. Thus, each minimum-weight FVS in $V_{0}$ for the graph $G$ contains exactly $\left|V_{0}\right|-\# c\left(V_{2}\right)+\# c(G)$ vertices, and such a minimum-weight FVS
can be constructed in time $\mathcal{O}\left(n^{2}\right)$. Therefore, for the given instance ( $G, V_{0}, V_{1}, V_{2}, k$ ) of the WEIGHTED IF-PARTITION FVS problem with $V_{1}=\emptyset$, the solution is "No" if $k<\left|V_{0}\right|-\# c\left(V_{2}\right)+\# c(G)$; and the solution is an $\mathcal{O}\left(n^{2}\right)$ time constructible FVS of $\left|V_{0}\right|-\# c\left(V_{2}\right)+\# c(G)$ vertices in $V_{0}$ if $k \geq\left|V_{0}\right|-\# c\left(V_{2}\right)+\# c(G)$.

This completes the proof that when $V_{1}=\emptyset$, a solution to the instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ can be constructed in time $\mathcal{O}\left(n^{2}\right)$.

Now consider the case $\tau\left(k, V_{0}, V_{1}, V_{2}\right) \leq 0$. If $V_{2}=\emptyset$, then by the first part of the proof, the lemma holds. Thus, we assume that $V_{2} \neq \emptyset$. As analyzed above, to break every cycle in the induced subgraph $G\left[V_{0} \cup V_{2}\right]$ we have to remove at least $\left|V_{0}\right|-\# c\left(V_{2}\right)+\# c\left(V_{0} \cup V_{2}\right)$ vertices in the set $V_{0}$. Therefore, if $\tau\left(k, V_{0}, V_{1}, V_{2}\right) \leq 0$, then $k \leq\left|V_{0}\right|-\# c\left(V_{2}\right)+1 \leq\left|V_{0}\right|-\# c\left(V_{2}\right)+\# c\left(V_{0} \cup V_{2}\right)$ (note that $V_{2} \neq \emptyset$ so $\left.\# c\left(V_{0} \cup V_{2}\right) \geq 1\right)$. Thus, in this case, all $k$ vertices in the objective FVS must be in the set $V_{0}$ in order to break all cycles in the induced subgraph $G\left[V_{0} \cup V_{2}\right]$, and no vertex in the objective FVS can be in the set $V_{1}$. Hence, if the induced subgraph $G\left[V_{1} \cup V_{2}\right]$ contains a cycle, then the solution to the instance is "No". On the other hand, suppose that $G\left[V_{1} \cup V_{2}\right]$ is a forest, then the graph $G$ has another IF-partition $\left(V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right)$, where $V_{0}^{\prime}=V_{0}, V_{1}^{\prime}=\emptyset$, and $V_{2}^{\prime}=V_{1} \cup V_{2}$. It is easy to verify that in this case the instance $\left(G, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}, k\right)$ with the IF-partition $\left(V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right)$ has the same solution set as the instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ with the IF-partition $\left(V_{0}, V_{1}, V_{2}\right)$. Since $V_{1}^{\prime}=\emptyset$, by the second part of the proof, a solution to the instance $\left(G, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}, k\right)$ with the IF-partition $\left(V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right)$ can be constructed in time $\mathcal{O}\left(n^{2}\right)$. This completes the proof of the lemma.

We are now in a position to introduce our main algorithm, which is given in Figure 6 and solves the WEIGHTED IF-PARTITION FVS problem. The subroutine min-w $\left(S_{1}, S_{2}\right)$ on two vertex subsets $S_{1}$ and $S_{2}$ in the algorithm returns among $S_{1}$ and $S_{2}$ the one with a smaller weight (or any one of them if the weights are tied). To
simplify our descriptions, we take the conventions that "No" is a special vertex set of an infinitely large weight and that any set plus "No" gives a "No". Therefore, the value of min-w $\left(S_{1}, S_{2}\right)$ will be (1) "No" if both $S_{1}$ and $S_{2}$ are "No"; (2) $S_{1}$ if $S_{2}$ is "No"; (3) $S_{2}$ if $S_{1}$ is "No"; and (4) the one of smaller weight among $S_{1}$ and $S_{2}$ if both $S_{1}$ and $S_{2}$ are not "No".

For each tree in the forest $G\left[V_{1}\right]$, we fix a root so that we can talk about the "lowest leaf" in a tree in $G\left[V_{1}\right]$.

Lemma C. 2 The algorithm W-Feedback $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ correctly solves the WEIGHTED IF-BIPARTITION FVS problem, and its running time is $\mathcal{O}\left(2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)} n^{2}\right)$, where $n$ is the number of vertices in the graph $G$.

Proof. We first verify the correctness of the algorithm. Step 1 of the algorithm is justified by Lemma C.1. Justifications for steps 2, 3, 4, 4.1, and 4.2 are exactly the same as that for steps 1, 4.1, 3, 3.1, and 3.2 in Algorithm-1 for unweighted graphs. Now consider step 5. When the algorithm reaches step 5, the following conditions hold:
(1) the sets $V_{1}$ and $V_{2}$ are not empty;
(2) every vertex in the set $V_{1}$ has degree at least 2 in the graph $G$; and
(3) every vertex in the set $V_{1}$ has at most one neighbor in the set $V_{2}$.

Condition (1) holds because of step 1; condition (2) holds because of step 3; and condition (3) holds because of step 4.

By condition (1) and because the induced subgraph $G\left[V_{1}\right]$ is a forest, step 5 can always pick the vertex $w_{1}$. By conditions (2) and (3), the vertex $w_{1}$ has a unique
neighbor in $V_{2}$. Also by conditions (2) and (3), the vertex $w_{1}$ must have a parent $w$ in the tree $T$ in $G\left[V_{1}\right]$. In consequence, the vertex $w_{1}$ has degree exactly 2 in $G$. Finally, since $w_{1}$ is the lowest leaf in the tree $T$, all children $w_{1}, \ldots, w_{t}$ of $w$ in the tree $T$ are also leaves in $T$. By conditions (2) and (3) again, each child $w_{i}$ of $w$ has a unique neighbor in the set $V_{2}$, and every child $w_{i}$ of $w$ has degree exactly 2 in the graph $G$.

Step 5.1 simply branches on the vertex $w$. To include the vertex $w$ in the objective FVS, we simply remove $w$ from the graph $G$ (and from the set $V_{1}$ ), and recursively look for an FVS in $V_{0} \cup\left(V_{1} \backslash\{w\}\right)$ of size at most $k-1$. Note that in this case, the sets $V_{0}$ and $V_{2}$ are unchanged, and the triple $\left(V_{0}, V_{1} \backslash\{w\}, V_{2}\right)$ obviously makes a valid IF-partition for the graph $G-w$. On the other hand, to exclude the vertex $w$ from the objective FVS, we move $w$ from $V_{1}$ to $V_{2}$. First note that since the vertex $w$ has at most one neighbor in $V_{2}$, the induced subgraph $G\left[V_{2} \cup\{w\}\right]$ is still a forest. Moreover, since all children $w_{1}, \ldots, w_{t}$ of $w$ have degree 2 in the graph $G$ and each $w_{i}$ has a unique neighbor in the set $V_{2}$, after moving $w$ from $V_{1}$ to $V_{2}$, all these degree- 2 vertices $w_{1}, \ldots, w_{t}$ have their both neighbors in the set $V_{2} \cup\{w\}$. Therefore, these vertices $w_{1}, \ldots, w_{t}$ now can be moved to the set $V_{0}$. In particular, the triple $\left(V_{0} \cup\left\{w_{1}, \ldots, w_{t}\right\}, V_{1} \backslash\left\{w, w_{1}, \ldots, w_{t}\right\}, V_{2} \cup\{w\}\right)$ is a valid IF-partition of the vertex set of the graph $G$. This recursive branching is implemented by the two recursive calls in step 5.1.

If we reach step 5.2 , then the two conditions in step 5.1 do not hold. Therefore, in addition to conditions (1)-(3), the following two conditions also hold:
(4) the vertex $w$ has no neighbor in $V_{2}$; and
(5) the vertex $w$ has a unique child $w_{1}$ in the tree $T$.

By conditions (2), (4), and (5), the vertex $w$ has degree exactly 2 in the graph $G$ (and $w$ is not the root of the tree $T$ ). Therefore, the vertices $w_{1}$ and $w$ are two adjacent degree-2 vertices in the graph $G$. Observe that in this case, a cycle in the graph $G$ contains the vertex $w_{1}$ if and only if it also contains the vertex $w$. Therefore, we can safely assume that the one of larger weight among $w_{1}$ and $w$ is not in the objective FVS. If the larger weight vertex is $w_{1}$, then the first recursive call in step 5.2 is executed, which moves $w_{1}$ from set $V_{1}$ to set $V_{2}$ (note that the triple $\left(V_{0}, V_{1} \backslash\left\{w_{1}\right\}, V_{2} \cup\left\{w_{1}\right\}\right)$ is a valid IF-partition of $G$ because $w_{1}$ has a unique neighbor in $V_{2}$ ). If the larger weight vertex is $w$, then the second recursive call in step 5.2 is executed, which moves $w$ from $V_{1}$ to $V_{2}$. Note that since both neighbors of the degree- 2 vertex $w_{1}$ are in the set $V_{2} \cup\{w\}$, after adding $w$ to the set $V_{2}$, we can also move the vertex $w_{1}$ from $V_{1}$ to $V_{0}$. Thus, the triple $\left(V_{0} \cup\left\{w_{1}\right\}, V_{1} \backslash\left\{w, w_{1}\right\}, V_{2} \cup\{w\}\right)$ is a valid IF-partition of the graph $G$.

We also remark that by our assumption, the input graph $G$ contains neither multiple edges nor self-loops. Moreover, the graph in each of the recursive calls in the algorithm is either the original $G$, or $G$ with a vertex deleted. Therefore, the graph in each of the recursive calls in the algorithm also contains neither multiple edges nor self-loops.

Since all possible cases are covered in the algorithm, we conclude that when the algorithm W-Feedback stops, it must output a correct solution to the given instance $\left(G, V_{0}, V_{1}, V_{2}, k\right)$.

To analyze the running time, as in the unweighted case, we first count the number of leaves in the search tree corresponding to the execution of the algorithm. Let $T\left(k, V_{0}, V_{1}, V_{2}\right)$ be the number of leaves in the search tree for algorithm $\mathbf{W}$ Feedback $\left(G, V_{0}, V_{1}, V_{2}, k\right)$. We prove by induction on the value $\tau\left(k, V_{0}, V_{1}, V_{2}\right)$ that $T\left(k, V_{0}, V_{1}, V_{2}\right) \leq \max \left(1,2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)}\right)$. First of all, if $\tau\left(k, V_{0}, V_{1}, V_{2}\right) \leq 0$, then by step

1 of the algorithm, we have $T\left(k, V_{0}, V_{1}, V_{2}\right)=1$.
First consider the branching steps, i.e., step 4.2 and step 5.1. In case step 4.2 of the algorithm is executed, we have recursively

$$
\begin{align*}
& T\left(k, V_{0}, V_{1}, V_{2}\right) \\
& \leq T\left(k-1, V_{0}, V_{1} \backslash\{w\}, V_{2}\right)+T\left(k, V_{0}, V_{1} \backslash\{w\}, V_{2} \cup\{w\}\right) . \tag{4.1}
\end{align*}
$$

Since

$$
\begin{aligned}
\tau_{1} & =\tau\left(k-1, V_{0}, V_{1} \backslash\{w\}, V_{2}\right) \\
& =(k-1)-\left(\left|V_{0}\right|-\# c\left(V_{2}\right)+1\right) \\
& =\tau\left(k, V_{0}, V_{1}, V_{2}\right)-1 \\
& <\tau\left(k, V_{0}, V_{1}, V_{2}\right)
\end{aligned}
$$

and

$$
\begin{aligned}
\tau_{2} & =\tau\left(k, V_{0}, V_{1} \backslash\{w\}, V_{2} \cup\{w\}\right) \\
& =k-\left(\left|V_{0}\right|-\# c\left(V_{2} \cup\{w\}\right)+1\right) \\
& \leq k-\left(\left|V_{0}\right|-\left(\# c\left(V_{2}\right)-1\right)+1\right) \\
& =\tau\left(k, V_{0}, V_{1}, V_{2}\right)-1 \\
& <\tau\left(k, V_{0}, V_{1}, V_{2}\right),
\end{aligned}
$$

where we have used the fact $\# c\left(V_{2} \cup\{w\}\right) \leq \# c\left(V_{2}\right)-1$ because in this case, we assume that the vertex $w$ has two neighbors in two different trees in $G\left[V_{2}\right]$, therefore, adding $w$ to $V_{2}$ merges at least two connected components in $G\left[V_{2}\right]$ and reduces the number of connected components by at least 1 .

Therefore, by the inductive hypothesis, $T\left(k-1, V_{0}, V_{1} \backslash\{w\}, V_{2}\right) \leq 2^{\tau_{1}}$, and $T\left(k, V_{0}, V_{1} \backslash\{w\}, V_{2} \cup\{w\}\right) \leq 2^{\tau_{2}}$. Combining these with Inequality (4.1), we get

$$
\begin{aligned}
T\left(k, V_{0}, V_{1}, V_{2}\right) & \leq T\left(k-1, V_{0}, V_{1} \backslash\{w\}, V_{2}\right)+T\left(k, V_{0}, V_{1} \backslash\{w\}, V_{2} \cup\{w\}\right) \\
& \leq 2^{\tau_{1}}+2^{\tau_{2}} \\
& \leq 2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)-1}+2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)-1} \\
& =2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)}
\end{aligned}
$$

In conclusion, the induction goes through for step 4.2 of the algorithm.
Now we consider step 5.1, which is the least trivial case, and makes the major difference from the unweighted cases. Let $V_{0}^{\prime}=V_{0} \cup\left\{w_{1}, \ldots, w_{t}\right\}, V_{1}^{\prime}=V_{1} \backslash$ $\left\{w, w_{1}, \ldots, w_{t}\right\}$, and $V_{2}^{\prime}=V_{2} \cup\{w\}$. The execution of step 5.1 gives the following inequality:

$$
T\left(k, V_{0}, V_{1}, V_{2}\right) \leq T\left(k-1, V_{0}, V_{1} \backslash\{w\}, V_{2}\right)+T\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right)
$$

As we have shown above, by the inductive hypothesis, we have

$$
\begin{equation*}
T\left(k-1, V_{0}, V_{1} \backslash\{w\}, V_{2}\right) \leq 2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)-1} \tag{4.2}
\end{equation*}
$$

To estimate the value $T\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right)$, first note that $\left|V_{0}^{\prime}\right|=\left|V_{0}\right|+t$. Moreover, at this point, we must have either that the vertex $w$ has a neighbor in $V_{2}$ or that the vertex $w$ has more than one child in the tree $T$ in $G\left[V_{1}\right]$.

If $w$ has a neighbor in $V_{2}$, then adding $w$ to $V_{2}$ will "attach" the vertex $w$ to a connected component in $G\left[V_{2}\right]$. In consequence, the number of connected components in $G\left[V_{2}\right]$ will be equal to that in $G\left[V_{2} \cup\{w\}\right]$ (recall that $w$ has only one neighbor in
$V_{2}$ ). In this case (note that $t \geq 1$ ), we have

$$
\begin{aligned}
\tau\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right) & =k-\left(\left|V_{0}^{\prime}\right|-\# c\left(V_{2}^{\prime}\right)+1\right) \\
& =k-\left(\left(\left|V_{0}\right|+t\right)-\# c\left(V_{2}\right)+1\right) \\
& \leq \tau\left(k, V_{0}, V_{1}, V_{2}\right)-1
\end{aligned}
$$

Now let us assume that the vertex $w$ has no neighbor in $V_{2}$ but has more than one child in the tree $T$ in $G\left[V_{1}\right]$ (i.e., $t \geq 2$ ). Then $\left|V_{0}^{\prime}\right|=\left|V_{0}\right|+t \geq\left|V_{0}\right|+2$. In this case, adding the vertex $w$ to the set $V_{2}$ increases the number of connected components in $G\left[V_{2}\right]$ by 1 (since $w$ has no neighbor in $V_{2}$, the vertex $w$ will become a single-vertex connected component in the induced subgraph $\left.G\left[V_{2} \cup\{w\}\right]\right)$. That is, $\# c\left(V_{2} \cup\{w\}\right)=\# c\left(V_{2}\right)+1$. Therefore,

$$
\begin{aligned}
\tau\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right) & =k-\left(\left|V_{0}^{\prime}\right|-\# c\left(V_{2}^{\prime}\right)+1\right) \\
& =k-\left(\left(\left|V_{0}\right|+t\right)-\# c\left(V_{2} \cup\{w\}\right)+1\right) \\
& \leq k-\left(\left(\left|V_{0}\right|+2\right)-\left(\# c\left(V_{2}\right)+1\right)+1\right) \\
& \leq \tau\left(k, V_{0}, V_{1}, V_{2}\right)-1
\end{aligned}
$$

In conclusion, in all cases in step 5.1, we will have $\tau\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right) \leq \tau\left(k, V_{0}, V_{1}, V_{2}\right)-$ 1. Therefore, now we can apply the induction and get

$$
T\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right) \leq 2^{\tau\left(k, V_{0}^{\prime}, V_{1}^{\prime}, V_{2}^{\prime}\right)} \leq 2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)-1}
$$

Combining this with the inequalities (4.1) and (4.2), we conclude that

$$
T\left(k, V_{0}, V_{1}, V_{2}\right) \leq 2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)}
$$

holds for the case of step 5.1.

We should also remark that it can be verified that for all non-branching recursive calls in the algorithm, i.e., steps 3, 4.1, and 5.2, the instance deficiency is never increased. In particular, if the first recursive call in step 5.2 is executed, then since the vertex $w_{1}$ has a unique neighbor in $V_{2}, \# c\left(V_{2}\right)=\# c\left(V_{2} \cup\left\{w_{1}\right\}\right)$. Thus,

$$
\begin{aligned}
\tau\left(k, V_{0}, V_{1} \backslash\left\{w_{1}\right\}, V_{2} \cup\left\{w_{1}\right\}\right) & =k-\left(\left|V_{0}\right|-\# c\left(V_{2} \cup\left\{w_{1}\right\}\right)+1\right) \\
& =k-\left(\left|V_{0}\right|-\# c\left(V_{2}\right)+1\right) \\
& =\tau\left(k, V_{0}, V_{1}, V_{2}\right) .
\end{aligned}
$$

If the second recursive call in step 5.2 is executed, then

$$
\# c\left(V_{2} \cup\{w\}\right)=\# c\left(V_{2}\right)+1
$$

because $w$ has no neighbor in $V_{2}$ and $w$ will become a single-vertex connected component in the induced subgraph $G\left[V_{2} \cup\{w\}\right]$. Therefore,

$$
\begin{aligned}
& \tau\left(k, V_{0} \cup\left\{w_{1}\right\}, V_{1} \backslash\left\{w, w_{1}\right\}, V_{2} \cup\{w\}\right) \\
= & k-\left(\left|V_{0} \cup\left\{w_{1}\right\}\right|-\# c\left(V_{2} \cup\{w\}\right)+1\right) \\
= & k-\left(\left(\left|V_{0}\right|+1\right)-\left(\# c\left(V_{2}\right)+1\right)+1\right) \\
= & \tau\left(k, V_{0}, V_{1}, V_{2}\right) .
\end{aligned}
$$

Summarizing all the above discussions, we complete the inductive proof that the number of leaves in the search tree for the algorithm $\mathbf{W}$-Feedback $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ is at most $2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)}$.

In the same way as in the proof for the unweighted case, we observe that along each root-leaf path in the search tree, the total number of executions of steps 1,2 , $3,4,4.1,4.2,5,5.1$, and 5.2 of the algorithm is $\mathcal{O}(n)$ because each of these steps either stops immediately, or reduces the size of the set $V_{1}$ by at least 1 . Step 1
is only preformed in leaf nodes of the tree, and thus only adds $\mathcal{O}\left(n^{2}\right)$ time to the total. By similar arguments as the one used for the unweighted case, all steps except step 1 can be preformed in $\mathcal{O}(n)$ time. Therefore, the running time of the algorithm W-Feedback $\left(G, V_{0}, V_{1}, V_{2}, k\right)$ is $\mathcal{O}\left(2^{\tau\left(k, V_{0}, V_{1}, V_{2}\right)} n^{2}\right)$.

With Lemma C.2, we can now proceed in the same way as for the unweighted case to solve the original WEIGHTED-FVS problem. Consider the following weighted version of the FVS REDUCTION problem.

WEIGHTED FVS REDUCTION: given a weighted graph $G$ and an FVS $F$ of size $k+1$ for $G$, either construct an FVS $F^{\prime}$ of minimum weight that satisfies $\left|F^{\prime}\right| \leq k$, or report that no such an FVS exists.

Note that in the definition of WEIGHTED FVS REDUCTION, we do not require that the given FVS $F$ of size $k+1$ have the minimum weight.

Lemma C. 3 The WEIGHTED FVS REDUCTION problem on an n-vertex graph is solvable in time $\mathcal{O}\left(5^{k} n^{2}\right)$.

Proof. The proof proceeds similarly to the proof of Lemma E.1. For the given FVS $F$ of size $k+1$ in the graph $G=(V, E)$, every FVS $F^{\prime}$ of size at most $k$ for $G$ (including the one with the minimum weight) is a union of a subset $F_{1}$ of at most $k-j$ vertices in $V \backslash F$ and a subset $F_{2}$ of $j$ vertices in $F$, for some integer $j, 0 \leq j \leq k$, where $\left(V \backslash F, F \backslash F_{2}\right)$ is a forest bipartition of the graph $G_{0}=G-F_{2}$. Therefore, we can enumerate all subsets $F_{2}$ of $j$ vertices in $F$, for each $j, 0 \leq j \leq k$, such that $\left(V \backslash F, F \backslash F_{2}\right)$ is a forest bipartition of the graph $G_{0}=G-F_{2}$, and construct the minimum-weight FVS $F_{0}$ of $G_{0}$ satisfying $\left|F_{0}\right| \leq k-j$. Note that the forest bipartition $\left(V \backslash F, F \backslash F_{2}\right)$ of $G_{0}$ is in fact a special IF-partition $\left(V_{0}, V_{1}, V_{2}\right)$ of $G_{0}$, where $V_{0}=\emptyset$,
$V_{1}=V \backslash F$, and $V_{2}=F \backslash F_{2}$. Therefore, by Lemma C.2, a minimum-weight FVS $F_{0}$ of $G_{0}$ satisfying $\left|F_{0}\right| \leq k-j$ can be constructed in time

$$
\mathcal{O}\left(2^{\tau\left(k-j, V_{0}, V_{1}, V_{2}\right)} n^{2}\right)=\mathcal{O}\left(2^{(k-j)-\left(0-\# c\left(F \backslash F_{2}\right)+1\right)} n^{2}\right)=\mathcal{O}\left(4^{k-j} n^{2}\right),
$$

where we have used the fact $\# c\left(F \backslash F_{2}\right) \leq\left|F \backslash F_{2}\right|=k+1-j$. Now the proof proceeds exactly the same way as that in Lemma E.1, and concludes that the WEIGHTED FVS REDUCTION problem can be solved in time $\mathcal{O}\left(5^{k} n^{2}\right)$.

Using Theorem B. 3 and Lemma C.3, we obtain the main result of this chapter.

Theorem C. 4 The WEIGHTED-FVs problem on an n-vertex graph is solvable in time $\mathcal{O}\left(5^{k} k n^{2}\right)$.

Proof. Let $(G, k)$ be a given instance of the weighted-FVs problem. As we explained in the proof of Theorem B.3, we can first construct, in time $\mathcal{O}\left(5^{k} k n^{2}\right)$, an FVS $F$ of size $k+1$ for the graph $G$ (the weight of $F$ is not necessarily the minimum). Then we simply apply Lemma C.3.

Algorithm-2 W-Feedback $\left(G, V_{0}, V_{1}, V_{2}, k\right)$
Input: $G=(V, E)$ is a graph with an IF-partition $\left(V_{0}, V_{1}, V_{2}\right), k$ is an integer.
Output: a minimum-weight FVS $F$ of $G$ such that $|F| \leq k$ and $F \subseteq V_{0} \cup V_{1}$; or report "No" (i.e., no such an FVS exists).

1 if $\left(V_{1}=\emptyset\right)$ or $\left(V_{2}=\emptyset\right)$ or $\left(\tau\left(k, V_{0}, V_{1}, V_{2}\right) \leq 0\right)$ then solve the problem in time $\mathcal{O}\left(n^{2}\right)$;
2 if $(k<0)$ or ( $k=0$ and $G$ is not a forest) then return "No";
3 if a vertex $w$ in $V_{1}$ has degree less than 2 in $G$ then return $\mathbf{W}$-Feedback $\left(G-w, V_{0}, V_{1} \backslash\{w\}, V_{2}, k\right)$;
4 else
if a vertex $w$ in $V_{1}$ has at least two neighbors in $V_{2}$ then
4.1 if two neighbors of $w$ are in the same tree of $G\left[V_{2}\right]$ then return $\left(\{w\} \cup \mathbf{W}\right.$-Feedback $\left.\left(G-w, V_{0}, V_{1} \backslash\{w\}, V_{2}, k-1\right)\right)$;
4.2 else
$F_{1}=\mathbf{W}$-Feedback $\left(G-w, V_{0}, V_{1} \backslash\{w\}, V_{2}, k-1\right)$;
$F_{2}=\mathbf{W}$-Feedback $\left(G, V_{0}, V_{1} \backslash\{w\}, V_{2} \cup\{w\}, k\right)$;
return min-w $\left(F_{1} \cup\{w\}, F_{2}\right)$;
5 else pick a lowest leaf $w_{1}$ in any tree $T$ in $G\left[V_{1}\right]$; let $w$ be the parent of $w_{1}$ in $T$, and let $w_{1}, \ldots, w_{t}$ be the children of $w$ in $T$;
5.1 if ( $w$ has a neighbor in $V_{2}$ ) or ( $w$ has more than one child in $T$ ) then $F_{1}=\mathbf{W}$-Feedback $\left(G-w, V_{0}, V_{1} \backslash\{w\}, V_{2}, k-1\right)$;
$F_{2}=\mathbf{W}$-Feedback $\left(G, V_{0} \cup\left\{w_{1}, \ldots, w_{t}\right\}, V_{1} \backslash\left\{w, w_{1}, \ldots, w_{t}\right\}\right.$,
$\left.V_{2} \cup\{w\}, k\right) ;$
return min-w $\left(F_{1} \cup\{w\}, F_{2}\right)$;
5.2 else
if the weight of $w_{1}$ is larger than the weight of $w$ then return W-Feedback $\left(G, V_{0}, V_{1} \backslash\left\{w_{1}\right\}, V_{2} \cup\left\{w_{1}\right\}, k\right)$;
else return W-Feedback $\left(G, V_{0} \cup\left\{w_{1}\right\}, V_{1} \backslash\left\{w, w_{1}\right\}, V_{2} \cup\{w\}, k\right)$.

Fig. 6. Algorithm for the Weighted feedback vertex set problem

## CHAPTER V

## DIRECTED FEEDBACK VERTEX SET*

In this chapter, we present the first fpt-algorithm of running time $O^{*}\left(4^{k} k!\right)$ for the FEEDBACK VERTEX SET problem on directed graphs. It had been an well-known open problem whether the FEEDBACK VERTEX SET problem on directed graphs is fixed-parameter tractable or not for 16 years.

Our algorithm transforms the FEEDBACK VERTEX SET problem into $O(k!)$ SKEW SEPARATOR problems, then solves each SKEW SEPARATOR problem in time of $O^{*}\left(4^{k}\right)$. Our algorithm for the SKEW SEPARATOR problem takes both the size of the skew separator to search and the minimum cut from the last source to all sinks as measures. The last measure is critical, because it avoids the need of bounding the length of cycles in graphs. In the rest of this chapter, we give detailed analysis of the algorithm for the feedback vertex set problem on directed graphs.

## A. Introduction

Let $G$ be a directed graph. A feedback vertex set $F$ (briefly, FVS) for $G$ is a set of vertices in $G$ such that every directed cycle in $G$ contains at least one vertex in $F$, or equivalently, that the removal of $F$ from the graph $G$ leaves a directed acyclic graph (i.e., a DAG). The (parameterized) FEEDBACK VERTEX SET problem on directed graphs (briefly, the DFVs problem) is defined as follows: given a directed graph $G$ and a parameter $k$, either construct an FVS of at most $k$ vertices for $G$ or report that no such set exists.

[^4]The DFVS problem is a classic NP-complete problem that appeared in the first list of NP-complete problems in Karp's seminal paper [68], and has a variety of applications in areas such as operating systems [94], database systems [52], and circuit testing [76]. In particular, the DFVS problem has played an essential role in the study of deadlock recovery in database systems and in operating systems [94, 52]. In such a system, the status of system resource allocations can be represented as a directed graph $G$ (i.e., the system resource-allocation graph), and a directed cycle in $G$ represents a deadlock in the system. Therefore, in order to recover from deadlocks, we need to abort a set of processes in the system, i.e., to remove a set of vertices in the graph $G$, so that all directed cycles in $G$ are broken. Equivalently, we need to find an FVS in the graph $G$. In practice, one may expect and desire that the number of vertices removed from the graph $G$, which is the number of processes to be aborted in the system, be small. This motivates the study of parameterized algorithms for the DFVS problem that find an FVS of $k$ vertices in a directed graph of $n$ vertices and run in time $f(k) n^{O(1)}$ for a fixed function $f$; thus, the algorithms become practically efficient when the value $k$ is small.

This work has been part of a systematic study of the theory of fixed-parameter tractability [37], which has received considerable attention in recent years. A problem $Q$ is a parameterized problem if each instance of $Q$ contains a specific integral parameter $k$. A parameterized problem is fixed-parameter tractable if it can be solved in time $f(k) n^{c}$ for a function $f(k)$ and a constant $c$, where the function $f(k)$ is independent of the instance size $n$. A large number of NP-hard parameterized problems, such as the VERTEX COVER problem [20] and the ML TYPE-CHECKING problem [77], have been shown to be fixed-parameter tractable. On the other hand, strong evidence has been given that another group of well-known parameterized problems, including the INDEpendent set problem and the Dominating set problem, are not fixed-parameter
tractable [37]. The study of fixed-parameter tractability of parameterized problems has become increasingly interesting, for both theoretical research and practical computation.

The fixed-parameter tractability of the DFVS problem was posted as an open problem in the very first papers on the study of fixed-parameter tractability [35, 34]. After numerous significant efforts, however, the problem still remained open. In the past fifteen years, the problem has been constantly and explicitly posted as an open problem in a large number of publications in the literature (see [58] for a recent survey on this study). The problem has become a well-known and outstanding open problem in parameterized computation and complexity.

In this chapter, we develop new algorithmic techniques that lead to the conclusion that the DFVS problem is fixed-parameter tractable, and thus resolve the above open problem in parameterized computation and complexity. We first show that the DFVS problem can be reduced in time $f(k) n^{O(1)}$ for some function $f$ to a special version of the multi-cut problem, which will be called the SKEW SEPARATOR problem. We then develop an algorithm that shows the fixed-parameter tractability of the SKEW SEPARATOR problem. The combination of these two results gives an algorithm with running time $4^{k} k!n^{O(1)}$ for the DFVS problem, which proves its fixed-parameter tractability.

The relationship between the DFVS problem and multi-cut problems has been studied in the research of approximation algorithms for the FEEDBACK VERTEX SET problem [41, 75]. However, our problem formulations and the corresponding techniques are significantly different from those studied in the approximation algorithms. In particular, our formulations and techniques seem especially suitable for developing faster and more effective exact algorithms (of exponential-time) for NP-hard multicut problems. First of all, instead of seeking a multi-cut that separates a given set
of terminal vertices, as formulated in most multi-cut problems, our problem is more general: we wish to construct a multi-cut that separates a collection of terminal vertex-subsets. This more general version of the multi-cut problem enables us to effectively reduce the search space size when we are searching for an optimal solution of a given problem instance. Secondly, unlike most multi-cut problems whose solutions are multi-cuts that are in general symmetric to the given terminal vertices, the multi-cuts for the SKEW SEPARATOR problem are asymmetric to the terminal vertexsubsets. Thirdly, we develop an (exponential-time) reduction that effectively reduces the problem of multi-cuts for multiple terminal vertex-subsets to the problem of minimum cuts from a single source vertex to a single sink vertex. Note that the latter is solvable in polynomial time via algorithms for the maximum flow problem. Such an exponential-time reduction is obviously very different from the polynomial-time processes used in the development of polynomial-time approximation algorithms. Finally, unlike most parameterized algorithms that are focused on effectively decreasing the sole parameter value $k$, our algorithm for the SKEW SEPARATOR problem adds another dimension of bounds in terms of the size of a minimum cut between two properly chosen terminal vertex-subsets. This dimension of bounds has become crucial in our development of the algorithm for the SKEW SEPARATOR problem because it effectively bounds the number of branches in which the parameter value $k$ is not decreased.

Before we move to the technical discussion of our algorithms, we remark that the FEEDBACK VERTEX SET problem on undirected graphs (briefly, the UFVS problem) has also been an interesting and active research topic in parameterized computation and complexity. Since the first fixed-parameter tractable algorithm for the UFVS problem was published fifteen years ago [9], there has been an impressive list of improved algorithms for the problem. Currently the best algorithm for the UFVS
problem runs in time $O\left(5^{k} k n^{2}\right)$ [23]. The FEEDBACK VERTEX SET problem on directed graphs (i.e., the DFVS problem) seems very different from the problem on undirected graphs (i.e., the UFVS problem). This fact has also been reflected in the study of approximation algorithms for the problems. The feedback vertex set problem on undirected graphs is polynomial-time approximable with a ratio 2 . This holds true even for weighted graphs [5]. On the other hand, it still remains open whether the FEEDBACK VERTEX SET problem on directed graphs has a constant-ratio polynomialtime approximation for the problem on directed graphs has a ratio $O(\log \tau \log \log \tau)$, where $\tau$ is the size of a minimum FVS for the input graph [41].

## B. Preliminaries

Let $G=(V, E)$ be a directed graph and let $e=[u, v]$ be a (directed) edge in $G$. We say that the edge $e$ goes out from the vertex $u$ and comes into the vertex $v$. The edge $e$ is called an outgoing edge of the vertex $u$, and an incoming edge of the vertex $v$. These concepts can be extended from single vertices to general vertex sets. Thus, for two vertex sets $S_{1}$ and $S_{2}$, we can say that an edge goes out from $S_{1}$ and comes into $S_{2}$ if the edge goes out from a vertex in $S_{1}$ and comes into a vertex in $S_{2}$. Moreover, we say that an edge goes out from $S_{1}$ if the edge goes out from a vertex in $S_{1}$ and comes into a vertex not in $S_{1}$, and that an edge comes into $S_{2}$ if the edge goes out from a vertex not in $S_{2}$ and comes into a vertex in $S_{2}$.

A path $P$ from a vertex $v_{1}$ to a vertex $v_{h}$ in the graph $G$ is a sequence $\left\{v_{1}, v_{2}, \ldots, v_{h}\right\}$ of vertices in $G$ such that $\left[v_{i}, v_{i+1}\right]$ is an edge in $G$ for all $1 \leq i \leq h-1$. The path $P$ is simple if no vertex is repeated in $P$. The path $P$ is a cycle if $v_{1}=v_{h}$, and the cycle is simple if no other vertices are repeated. We say that a path is from a vertex set $S_{1}$ to a vertex set $S_{2}$ if the path is from a vertex in $S_{1}$ to a vertex in $S_{2}$. The graph $G$
is a $D A G$ (i.e., directed acyclic graph) if it contains no cycles.
For a vertex subset $V^{\prime} \subseteq V$ in the directed graph $G=(V, E)$, we denote by $G\left[V^{\prime}\right]$ the subgraph of $G$ that is induced by the vertex subset $V^{\prime}$. Without any ambiguity, we will denote by $G-V^{\prime}$ the induced subgraph $G\left[V-V^{\prime}\right]$, and by $G-w$, where $w$ is a vertex in $G$, the induced subgraph $G[V-\{w\}]$.

A vertex subset $F$ in the directed graph $G$ is a feedback vertex set (FVS) if the graph $G-F$ is a DAG. Since a vertex $v$ with a self-loop (i.e., an edge that both goes out from and comes into $v$ ) must be included in every FVS for the graph $G$, we will assume, without loss of generality, that the graphs in our discussion have no self-loops.

Definition Let $\left[S_{1}, \ldots, S_{l}\right]$ and $\left[T_{1}, \ldots, T_{l}\right]$ be two collections of $l$ vertex subsets in a directed graph $G=(V, E)$. A skew separator $X$ for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ is a vertex subset in $V-\bigcup_{i=1}^{l}\left(S_{i} \cup T_{i}\right)$ such that for any pair of indices $i$ and $j$ satisfying $l \geq i \geq j \geq 1$, there is no path from $S_{i}$ to $T_{j}$ in the graph $G-X$.

The subsets $S_{1}, \ldots, S_{l}$ will be called the source sets and the subsets $T_{1}, \ldots, T_{l}$ will be called the sink sets. A vertex is a non-terminal vertex if it is not in $\bigcup_{i=1}^{l}\left(S_{i} \cup T_{i}\right)$. Note that by definition, all vertices in a skew separator must be non-terminal vertices. Moreover, a skew separator $X$ is asymmetric to the source sets and the sink sets: a path from $S_{i}$ to $T_{j}$ with $i<j$ may exist in the graph $G-X$.

When there is only one source set $S_{1}$ and one sink set $T_{1}$, a skew separator for the pair $\left(\left[S_{1}\right],\left[T_{1}\right]\right)$ becomes a regular cut for $S_{1}$ and $T_{1}$, i.e., a vertex set whose removal leaves a graph in which there is no path from the set $S_{1}$ to the set $T_{1}$. Therefore, a skew separator for the pair $\left(\left[S_{1}\right],\left[T_{1}\right]\right)$ is also called a cut from $S_{1}$ to $T_{1}$. A cut from $S_{1}$ to $T_{1}$ is a min-cut (i.e., a minimum cut) if it has the smallest cardinality over all
cuts from $S_{1}$ to $T_{1}$.
The following lemma can be easily derived based on standard maximum flow techniques [93]. Thus, we omit its proof.

Lemma B. 1 There is an $O\left(k n^{2}\right)$ time algorithm that for two given vertex subsets $S$ and $T$ in a directed graph $G$ of $n$ vertices, and a parameter $k$, either constructs $a$ min-cut from $S$ to $T$ whose size is bounded by $k$, or reports that the min-cut from $S$ to $T$ has a size larger than $k$.

The algorithm for the DFVS problem is obtained through careful development of algorithms for a series of problems. In the following, we give the formal definitions of these problems.

SKEW SEPARATOR: given $\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$, where $G$ is a directed graph, $\left[S_{1}, \ldots, S_{l}\right]$ is a collection of $l$ source sets and $\left[T_{1}, \ldots, T_{l}\right]$ is a collection of $l \operatorname{sink}$ sets in $G$, and a parameter $k$, such that

1. all sets $S_{1}, \ldots, S_{l}, T_{1}, \ldots, T_{l}$ are pairwise disjoint;
2. for each $i, 1 \leq i \leq l-1$, there is no edge coming into the set $S_{i}$; and
3. for each $j, 1 \leq j \leq l$, there is no edge going out from the set $T_{j}$,
either construct a skew separator of at most $k$ vertices for the pair ( $\left[S_{1}, \ldots, S_{l}\right]$, $\left[T_{1}, \ldots, T_{l}\right]$ ), or report that no such separator exists.

Note that in an instance of the SKEW SEPARATOR problem, condition (2) on source sets and condition (3) on sink sets are not completely symmetric. Although the first $l-1$ source sets are not allowed to have incoming edges, the last source set $S_{l}$ is allowed to have incoming edges. On the other hand, all sink sets are not allowed to have outgoing edges.

We remark that conditions (1)-(3) in the definition of the SKEW SEPARATOR problem (plus the restriction that the skew separator can consist of only non-terminal vertices) may be relaxed, and our techniques for the problem may still be applicable. However, the above formulation of the problem will make our discussion simpler, and will also be sufficient for our solution to the DFVS problem, which is the focus of the current chapter. We leave the investigation of the separator problems of more general forms to later research.

Let $G=(V, E)$ be a directed graph, and let $\left(D_{1}, D_{2}\right)$ be a bi-partition of the vertex set $V$ of $G$, i.e., $D_{1} \cup D_{2}=V$ and $D_{1} \cap D_{2}=\emptyset$. The bi-partition $\left(D_{1}, D_{2}\right)$ is a DAG-bipartition for the graph $G$ if both induced subgraphs $G\left[D_{1}\right]$ and $G\left[D_{2}\right]$ are DAGs. A vertex subset $F$ in the graph $G$ is a $D_{1}$-FVS if $F$ is an FVS for $G$ and $F \subseteq D_{1}$.

DAG-BIPARTITION FVS: given $\left(G, D_{1}, D_{2}, k\right)$, where $G$ is a directed graph, $\left(D_{1}, D_{2}\right)$ is a DAG-bipartition for $G$, and $k$ is the parameter, either construct a $D_{1}$-FVS of size bounded by $k$ for the graph $G$, or report that no such $D_{1}$-FVS exists.

We will be also interested in a special version of the FEEDBACK VERTEX SET problem.

DFVS REDUCTION: given a triple $(G, F, k)$, where $G$ is a directed graph and $F$ is an FVS of size $k+1$ for $G$, either construct an FVS of size bounded by $k$ for $G$, or report that no such FVS exists.

Finally, our central problem in this chapter is as follows.
DFVS: given a pair $(G, k)$, where $G$ is a directed graph and $k$ is the parameter, either construct an FVS of size bounded by $k$ for $G$, or report that no such FVS exists.

## C. Solving the skew separator Problem

In this section, we study the complexity of the SKEW SEPARATOR problem.
Let $\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$ be an instance of the SKEW SEPARATOR problem. Define $T_{\text {all }}=\bigcup_{1 \leq i \leq l} T_{i}$. There are a few cases in which we can directly reduce the instance size:

Rule R1 There is no path from $S_{l}$ to $T_{\text {all }}$, i.e., the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ is 0: then we only need to find a skew separator of size $k$ that separates $S_{i}$ from $T_{j}$ for all indices $i$ and $j$ satisfying $l-1 \geq i \geq j \geq 1$, i.e., we can work instead on the instance $\left(G,\left[S_{1}, \ldots, S_{l-1}\right],\left[T_{1}, \ldots, T_{l-1}\right], k\right)$. Note that in this case, by definition, if $l=1$, then the solution to the instance $\left(G,\left[S_{1}, \ldots, S_{l-1}\right],\left[T_{1}, \ldots, T_{l-1}\right], k\right)$ is simply the empty set $\emptyset$;

Rule R2 There is an edge from $S_{l}$ to $T_{\text {all }}$ : then there is no way to even separate $S_{l}$ from $T_{\text {all }}$ - we can simply stop and claim that the given instance is a "No" instance;

Rule R3 There exists a non-terminal vertex $w$, an edge from $S_{l}$ to $w$, and an edge from $w$ to $T_{\text {all }}$ : then the vertex $w$ must be included in the skew separator in order to separate $S_{l}$ and $T_{\text {all }}$ - we can simply work on the instance ( $G$ $\left.w,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k-1\right)$ and recursively find a skew separator of size $k-1$.

Note in Rules R1 and R3, the reduced instances $\left(G,\left[S_{1}, \ldots, S_{l-1}\right],\left[T_{1}, \ldots, T_{l-1}\right], k\right)$ and $\left(G-w,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k-1\right)$ are still valid instances of the SKEW SEPARATOR problem.

In the following, we assume that for the instance $\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$, none of the rules above is applicable. In particular, since Rule R1 is not applicable, a
min-cut from $S_{l}$ to $T_{\text {all }}$ has size larger than 0 . Because Rules R1-R3 are not applicable, there must be a non-terminal vertex $u_{0}$ such that (1) there is an edge from $S_{l}$ to $u_{0}$; and (2) there is no edge from $u_{0}$ to $T_{\text {all }}$. Such a vertex $u_{0}$ will be called an $S_{l}$-extended vertex. Fix an $S_{l}$-extended vertex $u_{0}$, let $S_{l}^{\prime}=S_{l} \cup\left\{u_{0}\right\}$.

We start with the following simple but important lemma. The proof of this lemma is straightforward. Thus, we leave it to the reader.

Lemma C. 1 Let $X$ be a subset of vertices in the graph $G$ that does not contain the $S_{l}$-extended vertex $u_{0}$. Then $X$ is a skew separator for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ if and only if $X$ is a skew separator for $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$.

Lemma C. 1 also directly implies the following two useful corollaries.

Corollary C. 2 A skew separator for $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$ is also $a$ skew separator for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$.

Corollary C. 3 The size of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$ in the graph $G$ is at least as large as the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ in $G$.

Now we are ready for our main theorem in this section.

Theorem C. 4 If the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ is equal to the size of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$, then the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ has a skew separator of size bounded by $k$ if and only if the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$ has a skew separator of size bounded by $k$.

Proof. $\Leftarrow$ : Suppose that the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$ has a skew separator $X^{\prime}$ of size bounded by $k$. By Corollary C. $2, X^{\prime}$ is also a skew separator
for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. In consequence, $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ has a skew separator of size bounded by $k$.
$\Rightarrow$ : Suppose that the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ has a skew separator $X$ of size bounded by $k$. If the skew separator $X$ does not contain the $S_{l}$-extended vertex $u_{0}$, then by Lemma C.1, $X$ is also a skew separator of size bounded by $k$ for the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$, and the theorem is proved. Therefore, we can assume that the set $X$ contains the $S_{l}$-extended vertex $u_{0}$. We will define another set $X^{\prime}$ that does not contain $u_{0}$. We will show that $\left|X^{\prime}\right| \leq|X|$ and that $X^{\prime}$ is a skew separator for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. Then the theorem will immediately follow.

Let $Y$ be a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$. Then $Y$ does not contain the $S_{l}$-extended vertex $u_{0}$. Moreover, since there is no edge coming into $S_{i}$ from outside of $S_{i}$ for all $i \leq l-1$, the set $Y$ does not contain any vertex in $\bigcup_{i=1}^{l-1} S_{i}$. In consequence, the set $Y$ consists of only non-terminal vertices. By Corollary C. $2, Y$ is also a cut from $S_{l}$ to $T_{\text {all }}$. Moreover, by the assumption of the theorem that the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ is equal to the size of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}, Y$ is actually also a min-cut from $S_{l}$ to $T_{\text {all }}$. Let $R_{Y}\left(S_{l}\right)$ be the set of vertices $v$ such that either $v \in S_{l}$ or there is a path from $S_{l}$ to $v$ in the subgraph $G-Y$. In particular, $u_{0} \in R_{Y}\left(S_{l}\right)$ because $Y$ does not contain $u_{0}$ and there is an edge from $S_{l}$ to $u_{0}$.

We introduce a number of sets as follows.

$$
\begin{aligned}
Z & =X \cap Y \\
X_{i n} & =X \cap R_{Y}\left(S_{l}\right) \\
X_{o u t} & =X-\left(X_{i n} \cup Z\right) .
\end{aligned}
$$

That is, the skew separator $X$ for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ is decomposed into three disjoint subsets $Z, X_{\text {in }}$, and $X_{\text {out }}$ (note that by definitions, $R_{Y}\left(S_{l}\right)$ and $Y$ do not


Fig. 7. Sets in the proof of Theorem C.4.
intersect).
Let $Y_{T}$ be the set of vertices $v$ in the min-cut $Y$ such that there is a path from $v$ to $T_{\text {all }}$ in the subgraph $G-X$. By definition, we have $Y_{T} \cap Z=\emptyset$. Let

$$
Y_{S}=Y-\left(Y_{T} \cup Z\right)
$$

Thus, the min-cut $Y$ from $S_{l}$ to $T_{\text {all }}$ is decomposed into three disjoint subsets $Z, Y_{T}$, and $Y_{S}$. Figure 7 gives an intuitive illustration of the sets $Z, X_{\text {in }}, X_{o u t}, Y_{T}, Y_{S}$, and $R_{Y}\left(S_{l}\right)$.

We first show that the set $Y^{\prime}=Y_{S} \cup Z \cup X_{\text {in }}$ is also a cut from $S_{l}$ to $T_{\text {all }}$. If by contradiction $Y^{\prime}$ is not a cut from $S_{l}$ to $T_{\text {all }}$, then there is a path $P_{1}$ from $S_{l}$ to $T_{\text {all }}$ in the subgraph $G-Y^{\prime}$. The path $P_{1}$ must contain vertices in the set $Y$ since $Y$ is a cut from $S_{l}$ to $T_{\text {all }}$. Let $w$ be the first vertex on the path $P_{1}$ that is in $Y$
when we traverse from $S_{l}$ to $T_{\text {all }}$ along the path $P_{1}$. Then $w$ must be in $Y_{T}$ since $Y^{\prime}$ contains both $Y_{S}$ and $Z$. Now the partial path $P_{1}^{\prime}$ of $P_{1}$ from $S_{l}$ to $w$ (not including $w)$ must be entirely contained in $R_{Y}\left(S_{l}\right)$ (note that the path $P_{1}$ does not intersect $\left.Y_{S} \cup Z\right)$. Moreover, the path $P_{1}^{\prime}$ contains neither vertices in $X_{i n} \cup Z$ (by the definition of the set $Y^{\prime}$ ) nor vertices in $X_{\text {out }}$ (since the sets $X_{\text {out }}$ and $R_{Y}\left(S_{l}\right)$ are disjoint). In summary, the subpath $P_{1}^{\prime}$ from $S_{l}$ to $w$ contains no vertex in the set $X$. Moreover, by the definition of the set $Y_{T}$, and $w \in Y_{T}$, there is a path $P_{1}^{\prime \prime}$ from $w$ to $T_{\text {all }}$ in the subgraph $G-X$. Now the concatenation of the paths $P_{1}^{\prime}$ and $P_{1}^{\prime \prime}$ would result in a path from $S_{l}$ to $T_{\text {all }}$ in the graph $G-X$, contradicting the fact that $X$ is a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. This contradiction shows that the set $Y^{\prime}$ must be a cut from $S_{l}$ to $T_{\text {all }}$.

Since $Y$ is a min-cut from $S_{l}$ to $T_{\text {all }}$, we have $|Y| \leq\left|Y^{\prime}\right|$. By definition, $Y=$ $Y_{S} \cup Z \cup Y_{T}$ and $Y^{\prime}=Y_{S} \cup Z \cup X_{i n}$. Also note that $Y_{S}, Z$, and $Y_{T}$ are pairwise disjoint, and that $Y_{S}, Z$, and $X_{i n}$ are also pairwise disjoint. Therefore, we must have $\left|Y_{T}\right| \leq\left|X_{i n}\right|$.

Consider the set $X^{\prime}=X_{\text {out }} \cup Z \cup Y_{T}$. The set $X^{\prime}$ has the following properties: (1) $X^{\prime}$ consists of only non-terminal vertices (because both $X$ and $Y$ consist of only non-terminal vertices); (2) $\left|X^{\prime}\right| \leq|X|$ (because $\left|Y_{T}\right| \leq\left|X_{\text {in }}\right|$ ), so the size of $X^{\prime}$ is bounded by $k$; and (3) the set $X^{\prime}$ does not contain the $S_{l}$-extended vertex $u_{0}$ (this is because $u_{0}$ is in $X_{\text {in }}$ and $Y$ does not contain $u_{0}$ ). Therefore, if we can prove that $X^{\prime}$ is a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$, then by Lemma C.1, $X^{\prime}$ is also a skew separator of size bounded by $k$ for the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$. This will complete the proof of the theorem.

Therefore, what remains is to prove that the set $X^{\prime}=X_{\text {out }} \cup Z \cup Y_{T}$ is a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. Let $R_{Y}\left(T_{\text {all }}\right)$ be the set of vertices $v$ such that either $v \in T_{\text {all }}$, or there is a path from $v$ to $T_{\text {all }}$ in the subgraph $G-Y$.

Suppose by contradiction that the set $X^{\prime}$ is not a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. Then there is a path $P_{2}$ in the subgraph $G-X^{\prime}$ from $S_{i}$ to $T_{j}$ for some $i \geq j$. The path $P_{2}$ has the following properties:

1. The path $P_{2}$ must contain a vertex in $R_{Y}\left(S_{l}\right)$ : since $X$ is a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$, the path $P_{2}$ from $S_{i}$ to $T_{j}$ with $i \geq j$ must contain at least one vertex $w_{1}$ in $X=X_{\text {in }} \cup Z \cup X_{\text {out }}$. Now since the path $P_{2}$ is in the subgraph $G-X^{\prime}$, where $X^{\prime}=X_{\text {out }} \cup Z \cup Y_{T}$, the vertex $w_{1}$ must be in $X_{i n}$, which is a subset of $R_{Y}\left(S_{l}\right)$;
2. The path $P_{2}$ must contain a vertex in $Y_{S}$ : by Property (1), $P_{2}$ contains a vertex $w_{1}$ in $R_{Y}\left(S_{l}\right)$. From the vertex $w_{1}$ to $T_{\text {all }}$ along the path $P_{2}$, there must be a vertex $w_{2}$ in $Y=Y_{S} \cup Z \cup Y_{T}$ since $Y$ is a cut from $S_{l}$ to $T_{\text {all }}$ while $w_{1}$ is reachable from $S_{l}$ in the subgraph $G-Y$. Now since $X^{\prime}=X_{o u t} \cup Z \cup Y_{T}$, and the path $P_{2}$ is in the subgraph $G-X^{\prime}$, the vertex $w_{2}$ on the path $P_{2}$ must be in the set $Y_{S}$;
3. The path $P_{2}$ must end at a vertex in $R_{Y}\left(T_{\text {all }}\right)$; this is simply because $P_{2}$ is ended in $T_{\text {all }}$. Note that by definition, no vertex in $Y_{S}$ can be in $R_{Y}\left(T_{\text {all }}\right)$.

By Properties (2)-(3), the path $P_{2}$ contains a vertex not in $R_{Y}\left(T_{\text {all }}\right)$ and ends at a vertex in $R_{Y}\left(T_{\text {all }}\right)$. Thus, there must be an internal vertex $w$ in the path such that $w$ is not in $R_{Y}\left(T_{\text {all }}\right)$ but all vertices after $w$ along the path $P_{2}$ (from $S_{i}$ to $T_{j}$ ) are in $R_{Y}\left(T_{\text {all }}\right)$. Note that no vertex $w^{\prime}$ after the vertex $w$ along the path $P_{2}$ can be in the set $X: w^{\prime}$ in $X$ would imply $w^{\prime}$ in $X_{i n}$ (since $P_{2}$ is a path in the subgraph $G-X^{\prime}$ ), which would imply that there is another vertex after $w^{\prime}$ that is in $Y$ thus is not in $R_{Y}\left(T_{\text {all }}\right)$. Moreover, the vertex $w$ must be in the set $Y$ (otherwise, $w$ would be in $\left.R_{Y}\left(T_{\text {all }}\right)\right)$. Since $P_{2}$ is a path in $G-X^{\prime}$ and $X^{\prime}=X_{\text {out }} \cup Z \cup Y_{T}$, the vertex
$w$ must be in the set $Y_{S}$. However, this derives a contradiction: the subpath of $P_{2}$ from $w$ to $T_{\text {all }}$ shows that the vertex $w$ should belong to the set $Y_{T}$ (note that all vertices after $w$ on the path are not in $X$ ), and the sets $Y_{S}$ and $Y_{T}$ are disjoint. This contradiction proves that the set $X^{\prime}$ must be a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. Since the size of the set $X^{\prime}$ is bounded by $k$ and $X^{\prime}$ does not contain the $S_{l}$-extended vertex $u_{0}$, by Lemma C.1, the set $X^{\prime}$ is also a skew separator for the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$, and the size of $X^{\prime}$ is bounded by $k$.

This completes the proof of the theorem.
Theorem C. 4 enables us to develop a parameterized algorithm for the SKEW SEPARATOR problem. The algorithm is presented in Figure 8.

Theorem C. 5 The algorithm $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right), k\right]$ solves the problem SKEW SEPARATOR in time $O\left(4^{k} k n^{3}\right)$, where $n$ is the number of vertices in the input graph $G$.

Proof. For the correctness of the algorithm, let $\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$ be an input to the algorithm, which is an instance of the SKEW SEPARATOR problem, where $G=(V, E)$ is a directed graph, $\left[S_{1}, \ldots, S_{l}\right]$ and $\left[T_{1}, \ldots, T_{l}\right]$ are the source sets and the sink sets, respectively, and $k$ is the upper bound of the size of the skew separator we are looking for.

If $l=1$, then the problem becomes the construction of a min-cut of size bounded by $k$ from $S_{1}$ to $T_{1}$, which can be solved in $O\left(k n^{2}\right)$ time by Lemma B.1. Steps 2-4 were justified in the discussions of Rules 2, 1, 3, respectively, at the beginning of this section (note that we have also consistently defined that an instance is a "No" instance if the parameter $k$ has a negative value). Therefore, if the algorithm reaches step 5, then none of the Rules 1-3 are applicable. In particular, since Rule 1 is not
applicable and the sets $S_{l}$ and $T_{\text {all }}$ are disjoint, there must be an edge $[v, w]$, where $v \in S_{l}$ and $w \notin S_{l}$. Since Rule 2 is not applicable, the vertex $w$ is not in the set $T_{\text {all }}$. The vertex $w$ also cannot be in any source set $S_{i}$ for $i<l$ because there is no edge coming into $S_{i}$ from outside of $S_{i}$. Therefore, the vertex $w$ is a non-terminal vertex. Finally, since Rule 3 is not applicable, there is no edge from $w$ to $T_{\text {all }}$. Thus, $w$ must be an $S_{l}$-extended vertex. This proves that at step 5 , the algorithm can always find an $S_{l}$-extended vertex $u_{0}$.

In the case $m>k$ in step 7 , i.e., the size $m$ of a min-cut from $S_{l}$ to $T_{\text {all }}$ is larger than the parameter $k$, then even separating a single source set $S_{l}$ from the sink sets $T_{\text {all }}=\bigcup_{j=1}^{l} T_{j}$ requires more than $k$ vertices. Thus, no skew separator of size bounded by $k$ can exist to separate $S_{i}$ from $T_{j}$ for all $l \geq i \geq j \geq 1$. Step 7 correctly handles this case by returning "No".

In the case $m=m^{\prime}$ in step 9 , i.e., the size $m$ of a min-cut from $S_{l}$ to $T_{\text {all }}$ is equal to the size $m^{\prime}$ of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$, by Theorem C.4, the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$ has a skew separator of size bounded by $k$ if and only if the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$ has a skew separator of size bounded by $k$. Moreover, by Corollary C.2, a skew separator of size bounded by $k$ for the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$ is also a skew separator for $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$. Therefore, in this case we can recursively call the algorithm $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right]\right.$, $\left.\left[T_{1}, \ldots, T_{l-1}, T_{l}\right], k\right)$, and look instead for a skew separator of size bounded by $k$ for the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$, as handled by step 9.1.

In the case $m \neq m^{\prime}$, then the algorithm branches into two subcases: step 9.2 includes the $S_{l}$-extended vertex $u_{0}$ in the skew separator and recursively looks for a skew separator of size bounded by $k-1$ in the remaining graph $G-u_{0}$ for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$; and step 9.3 excludes the $S_{l}$-extended vertex $u_{0}$ from the skew separator and recursively looks for a skew separator that does not contain $u_{0}$ and
is of size bounded by $k$ in the graph $G$ for the pair $\left(\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]\right)$, which is a skew separator of size bounded by $k$ for the pair $\left(\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right]\right)$ by Lemma C.1. This completes the verification of the correctness of the algorithm. Now we analyze its complexity.

The recursive execution of the algorithm can be described as a search tree $\mathcal{T}$. We first count the number of leaves in the search tree $\mathcal{T}$. Note that only steps 9.2-9.3 of the algorithm correspond to branches in the search tree $\mathcal{T}$. Let $D(k, m)$ be the total number of leaves in the search tree $\mathcal{T}$ for the algorithm $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$, where $m$ is the size of a min-cut from $S_{l}$ to $T_{\text {all }}$. Then steps 9.2-9.3 induce the following recurrence relation:

$$
\begin{equation*}
D(k, m) \leq D\left(k-1, m_{1}\right)+D\left(k, m_{2}\right) \tag{5.1}
\end{equation*}
$$

where $m_{1}$ is the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ in the graph $G-u_{0}$ as given in step 9.2 , and $m_{2}$ is the size of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$ in the graph $G$ as given in step 9.3. Note that $m-1 \leq m_{1} \leq m$ because removing the vertex $u_{0}$ from the graph $G$ cannot increase the size of a min-cut from $S_{l}$ to $T_{\text {all }}$, and can decrease the size of a min-cut for the two sets by at most 1 . Moreover, by Corollary C.3, in step 9.3 we must have $m_{2} \geq m+1$. Summarizing these, we have

$$
\begin{equation*}
m-1 \leq m_{1} \leq m \quad \text { and } \quad m_{2} \geq m+1 \tag{5.2}
\end{equation*}
$$

We prove, by induction on $t=2 k-m$, that $D(k, m) \leq 2^{2 k-m}$. First note that we always have $t=2 k-m \geq 0$ because by the definitions of $k$ and $m$ we always have $k \geq m \geq 0$. In particular, in the initial case when $t=2 k-m=0$, we must have $k=m=0$; in this case the algorithm can solve the instance without further branching. Therefore, we have $D(k, m)=1$ when $t=2 k-m=0$. For the inductive
step, note that by Inequalities (5.2), we have

$$
t_{1}=2(k-1)-m_{1} \leq 2(k-1)-(m-1)=2 k-m-1,
$$

and

$$
t_{2}=2 k-m_{2} \leq 2 k-(m+1)=2 k-m-1
$$

Therefore, we can apply the inductive hypothesis on Inequality (5.1), which gives

$$
\begin{align*}
D(k, m) & \leq D\left(k-1, m_{1}\right)+D\left(k, m_{2}\right) \\
& \leq 2^{2(k-1)-m_{1}}+2^{2 k-m_{2}} \\
& \leq 2^{2 k-m-1}+2^{2 k-m-1} \\
& =2^{2 k-m} . \tag{5.3}
\end{align*}
$$

This completes the inductive proof. Moreover, we also note that certain non-branching steps (i.e., steps 3, 4, and 9.1) may also change the values of $k$ and $m$, thus changing the value $t=2 k-m$. However, none of these steps increases the value $t=2 k-m$ : (i) step 3 keeps the value $k$ unchanged and does not decrease the value $m$ (because in this case the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ is 0 that cannot be larger than the size of a min-cut from $S_{l-1}$ to $\bigcup_{j=1}^{l-1} T_{j}$ ); (ii) step 4 decreases the value $k$ by 1 and the value $m$ by at most 1 (because removing a vertex from $G$ can reduce the size of a min-cut from $S_{l}$ to $T_{\text {all }}$ by at most 1), which as a total will decrease the value $t=2 k-m$ by at least 1 ; (iii) by the condition assumed, step 9.1 keeps both the values $k$ and $m$ unchanged, thus unchanging the value $t=2 k-m$. As a result, the value $t=2 k-m$ after a branching step to the next branching step can never be increased.

Summarizing the above discussion, we conclude that the total number of leaves, $D(k, m)$, in the search tree $\mathcal{T}$ for the algorithm $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$,
where $m$ is the size of a min-cut from $S_{l}$ to $T_{\text {all }}$, satisfies the following inequality

$$
D(k, m) \leq 2^{2 k-m} \leq 4^{k}
$$

The running time of each execution of the algorithm SMC, not counting the time for the recursive calls in the execution, is bounded by $O\left(k n^{2}\right)$, where $n$ is the number of vertices in the input graph. In particular, by Lemma B.1, step 1 that looks for a min-cut of size bounded by $k$ from $S_{1}$ to $T_{1}$, steps 6-7 that determine if the size $m$ of a min-cut from $S_{l}$ to $T_{\text {all }}$ is bounded by $k$, and steps 8-9 that determine if the size of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$ is equal to $m$ ( $m \leq k$ at this point), all have their running time bounded by $O\left(k n^{2}\right)$.

Observe that for each recursive call in an execution of the algorithm SMC, either the number of source-sink pairs in the instance is decreased by 1 (step 3 ), or the number of non-terminal vertices in the instance is decreased by 1 (steps 4, 9.1, 9.2, and 9.3). When the number of source-sink pairs is equal to 1 , the problem is solved in time $O\left(k n^{2}\right)$ by step 1 , and when the number of non-terminal vertices is equal to 0 , either step 2 or step 3 can be applied directly. In conclusion, along each root-leaf path in the search tree $\mathcal{T}$, there are at most $O(n)$ recursive calls to the algorithm SMC. Therefore, the running time along each root-leaf path in the search tree $\mathcal{T}$ is bounded by $O\left(k n^{3}\right)$.

Summarizing the above discussions, we conclude that the running time of the algorithm SMC is bounded by $O\left(4^{k} k n^{3}\right)$. This completes the proof of the theorem.

## D. Solving the DAG-bipartition fVs Problem

In this section, we describe how to use the results in the previous section to solve the DAG-BIPARTITION FVS problem.

Recall that an instance of DAG-BIPARTITION FVS is given as a tuple ( $G, D_{1}, D_{2}, k$ ), where $G$ is a directed graph, $\left(D_{1}, D_{2}\right)$ is a DAG-bipartition of $G$, and $k$ is the parameter, with the objective of finding an FVS $X$ for the graph $G$ such that $X \subseteq D_{1}$ (recall that such an FVS is called a $D_{1}$-FVS) and that the size of $X$ is bounded by $k$.

Let $\pi=\left\{v_{1}, v_{2}, \ldots, v_{h}\right\}$ be a topologically sorted order of the vertices in the induced DAG $G\left[D_{2}\right]$. We construct an instance of the SKEW SEPARATOR problem as follows:

1. Let $G^{\prime}$ be the graph obtained from $G$ by removing all edges in $G\left[D_{2}\right]$.
2. In the graph $G^{\prime}$, replace each vertex $v_{i}$ in $D_{2}$ by a pair $\left(t_{i}, s_{i}\right)$ of vertices such that all incoming edges into $v_{i}$ are now coming into the vertex $t_{i}$, and that all outgoing edges from $v_{i}$ are now going out from the vertex $s_{i}$. Let the resulting graph be $G_{\pi}$.

Note that in the resulting graph $G_{\pi}$, the vertices $s_{i}, 1 \leq i \leq h$, have no incoming edges, and the vertices $t_{j}, 1 \leq j \leq h$, have no outgoing edges. Moreover, since we have removed all edges between the vertices in $G\left[D_{2}\right]$, every edge going out from a vertex $s_{i}$ must come into a vertex in the set $D_{1}$, and every edge coming into a vertex $t_{j}$ must go out from a vertex in the set $D_{1}$. In particular, $\left(G_{\pi},\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right], k\right)$ is a valid instance for the SKEW SEPARATOR problem, which will be called an instance of the SKEW SEPARATOR induced by the instance $\left(G, D_{1}, D_{2}, k\right)$ of DAG-BIPARTITION FVS and the topologically sorted order $\pi$ of the vertices in $G\left[D_{2}\right]$.

Thus, each vertex $v_{i}$ in the set $D_{2}$ in the graph $G$ is now "split" into the two
vertices $s_{i}$ and $t_{i}$ in the graph $G_{\pi}$. Moreover, there is a one-to-one mapping between the vertices in the set $D_{1}$ in the graph $G$ and the non-terminal vertices in the graph $G_{\pi}$. Thus, in case of no ambiguity, we will use the same vertex name to refer to both a non-terminal vertex in the graph $G_{\pi}$ and a vertex in the set $D_{1}$ in the graph G. In particular, a skew separator for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$ corresponds to a subset of $D_{1}$ in the graph $G$. We have the following important theorem.

Theorem D. 1 Let $\left(G, D_{1}, D_{2}, k\right)$ be an instance of the DAG-BIPARTITION FVs problem, and let $X$ be a $D_{1}$-FVS for the graph $G$. Then there is a topologically sorted order $\pi=\left\{v_{1}, \ldots, v_{h}\right\}$ of the vertices in $G\left[D_{2}\right]$ such that in the instance $\left(G_{\pi},\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right]\right.$, $\left.\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right], k\right)$ induced by $\left(G, D_{1}, D_{2}, k\right)$ and $\pi$ : (1) $X$ is a skew separator for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$; and (2) every skew separator for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in $G_{\pi}$ is a $D_{1}-F V S$ for the graph $G$.

Proof. As assumed in the theorem, let $\left(G, D_{1}, D_{2}, k\right)$ be an instance of the DAGbipartition fVs problem, and let $X$ be a $D_{1}$-FVS for the graph $G$. Consider the subgraph $G-X$. Since $X$ is an FVS for $G$, the graph $G-X$ is a DAG. Therefore, the vertices in $G-X$ can be topologically sorted into an ordered list $\pi^{\prime}$ such that there is no edge in $G-X$ that goes out from a later vertex in $\pi^{\prime}$ and comes into an earlier vertex in $\pi^{\prime}$. Let $\pi=\left\{v_{1}, \ldots, v_{h}\right\}$ be the order of the vertices in $D_{2}$ that is induced from the order $\pi^{\prime}$ (i.e., $\pi$ is obtained from $\pi^{\prime}$ by removing the vertices not in $D_{2}$. Note that all vertices in $X$ are in $D_{1}$ ). The order $\pi$ is obviously a topologically sorted order for the DAG $G\left[D_{2}\right]$. We show that this order $\pi$ of the vertices in $D_{2}$ and the corresponding instance $\left(G_{\pi},\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right], k\right)$ induced by $\left(G, D_{1}, D_{2}, k\right)$ and $\pi$ satisfy the conclusions of the theorem.

We first show that the set $X$ is a skew separator for $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$. If this were not the case, then there would be a path $P$ in the graph $G_{\pi}-X$ that starts from a vertex $s_{i}$ and ends at a vertex $t_{j}$ with $i \geq j$. Since no vertex in $\left\{s_{1}, \ldots, s_{h}\right\}$ has incoming edges and no vertex in $\left\{t_{1}, \ldots, t_{h}\right\}$ has outgoing edges, all internal vertices on the path $P$ are non-terminal vertices in $G_{\pi}$. In consequence, all internal vertices on $P$ are vertices in the set $D_{1}$ in the graph $G$. Therefore, the path $P$ in $G_{\pi}-X$ corresponds to a path $P^{\prime}$ in the graph $G-X$ that starts from the vertex $v_{i}$ and ends at the vertex $v_{j}$, where $i \geq j$, with all internal vertices of $P^{\prime}$ in the set $D_{1}$. But it is impossible: (1) if $i=j$ then the path $P^{\prime}$ would be a cycle in the graph $G-X$, contradicting the assumption that $X$ is an FVS for the graph $G$; and (2) if $i>j$, then $P^{\prime}$ would become a path from $v_{i}$ to $v_{j}$ with $i>j$ in the graph $G-X$, contradicting the assumption that $\pi=\left\{v_{1}, \ldots, v_{h}\right\}$ is an order of the vertices in $D_{2}$ that is induced from the topologically sorted order $\pi^{\prime}$ of the vertices in the DAG $G-X$. In conclusion, the path $P$ does not exist, and the set $X$ is a skew separator for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$.

Now we prove that every skew separator $X^{\prime}$ for $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$ is a $D_{1}$-FVS for the graph $G$. First of all, by definition, a skew separator consists of only non-terminals, thus, all vertices in $X^{\prime}$ are in the set $D_{1}$. Suppose for a contradiction that $X^{\prime}$ is not a $D_{1}$-FVS for the graph $G$. Then there is a cycle $C$ in the graph $G-X^{\prime}$. Without loss of generality, we can assume that $C$ is a simple cycle. Since both the induced subgraphs $G\left[D_{1}\right]$ and $G\left[D_{2}\right]$ are DAGs, the cycle $C$ must contain both vertices in $D_{1}$ and vertices in $D_{2}$. We consider two different cases.

Case 1. The cycle $C$ contains a single vertex $v_{i}$ in the set $D_{2}$. Then all other vertices in the cycle $C$ are in the set $D_{1}$. But then the cycle $C$ would correspond to a path $P_{1}$ in the graph $G_{\pi}-X^{\prime}$ that starts with the vertex $s_{i}$ and ends at the vertex
$t_{i}$ (with all internal vertices being non-terminal vertices). But this contradicts the assumption that $X^{\prime}$ is a skew separator for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ that should have cut all paths from $s_{i}$ to $t_{i}$.

Case 2. The cycle $C$ contains more than one vertex in $D_{2}$. Let $\left\{v_{i_{1}}, v_{i_{2}}, \ldots, v_{i_{d}}, v_{i_{1}}\right\}$ be the order of the vertices in $D_{2}$ that we encounter when traversing along the cycle $C$ (starting from an arbitrary vertex $v_{i_{1}}$ in $D_{2}$ ), where $d>1$. Then there must be an index $j$ such that $i_{j}>i_{j+1}$ (where we take $i_{j+1}=i_{1}$ if $j=d$ ). Now consider the subpath $P_{2}$ of $C$ that starts from the vertex $v_{i_{j}}$ and ends at the vertex $v_{i_{j+1}}$. The path $P_{2}$ cannot be a single edge from $v_{i_{j}}$ to $v_{i_{j+1}}$ since $\pi=\left\{v_{1}, v_{2}, \ldots, v_{h}\right\}$ is a topologically sorted order for the vertices in the DAG $G\left[D_{2}\right]$ and $i_{j}>i_{j+1}$. Thus, the path $P_{2}$ contains at least one internal vertex. Since all internal vertices on the path $P_{2}$ are not in $D_{2}$ thus correspond to non-terminal vertices in the graph $G_{\pi}-X^{\prime}$, the path $P_{2}$ would correspond to a path $P_{2}^{\prime}$ in the graph $G_{\pi}-X^{\prime}$ that starts from the vertex $s_{i_{j}}$ and ends at the vertex $t_{i_{j+1}}$, with $i_{j}>i_{j+1}$. Again this contradicts the assumption that $X^{\prime}$ is a skew separator for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$, which should have cut all paths from $s_{i_{j}}$ to $t_{i_{j+1}}$ when $i_{j}>i_{j+1}$.

This proves that the skew separator $X^{\prime}$ for $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$ must be a $D_{1}$-FVS for the graph $G$. This completes the proof of the theorem.

Theorem D. 1 enables us to reduce the DAG-BIPARTITION FVS problem to the SKEW SEPARATOR problem. An algorithm for the DAG-BIPARTITION FVS problem is given in Figure 9.

Theorem D. 2 The algorithm $\operatorname{DBF}\left(G, D_{1}, D_{2}, k\right)$ solves the DAG-BIPARTITION FVS problem in time $O\left(4^{k} k n^{3} h!\right)$, where $h$ is the number of vertices in the set $D_{2}$, and $n$ is the number of vertices in the input graph $G$.

Proof. The running time of the algorithm is obvious: the for-loop in step 1 is executed at most $h$ ! times, and the time for each execution is dominated by the subroutine call to the algorithm $\mathbf{S M C}$ in step 1.2. By Theorem C.5, the running time of each execution of step 1.2 is bounded by $O\left(4^{k} k n^{3}\right)$.

For the correctness of the algorithm, first note that the algorithm always returns "No" unless it actually constructs a $D_{1}$-FVS of size bounded by $k$ for $G$ in step 1.3. In particular, if the input instance ( $G, D_{1}, D_{2}, k$ ) contains no $D_{1}$-FVS of size bounded by $k$ for the graph $G$, then the algorithm always correctly reports "No".

On the other hand, suppose that there is a $D_{1}$-FVS $X_{0}$ of size bounded by $k$ for the graph $G$. Then by Theorem D.1, there is a topologically sorted order $\pi=$ $\left\{v_{1}, \ldots, v_{h}\right\}$ of the vertices in $G\left[D_{2}\right]$ such that in the instance $\left(G_{\pi},\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right]\right.$, $\left.\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right], k\right)$ of the SKEW SEPARATOR problem induced by $\left(G, D_{1}, D_{2}, k\right)$ and $\pi, X_{0}$ is a skew separator for $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right]\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$, and every skew separator for $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in $G_{\pi}$ is a $D_{1}$-FVS for the graph $G$. In particular, $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ has at least one skew separator of size bounded by $k$ (e.g., $X_{0}$ ) in the graph $G_{\pi}$. Therefore, step 1.2 of the algorithm DBF must return a skew separator $X$ of size bounded by $k$ for the pair $\left(\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right]\right)$ in the graph $G_{\pi}$ (the set $X$ may be different from the set $X_{0}$ ), and this set $X$ is a $D_{1}$-FVS for the graph $G$. In conclusion, if there is a $D_{1}$-FVS of size bounded by $k$ for the graph $G$, then the algorithm $\operatorname{DBF}\left(G, D_{1}, D_{2}, k\right)$ will correctly return a $D_{1}$-FVS of size bounded by $k$ in step 1.3 .

## E. Solving the DFvs Problem

We present our algorithm for the DFVS problem. We start with a more restricted version of the problem, the DFVS REDUCTION problem, defined as follows.

DFVS REDUCTION: given a triple $(G, F, k)$, where $G$ is a directed graph and $F$ is an FVS of size $k+1$ for $G$, either construct an FVS of size bounded by $k$ for $G$, or report that no such FVS exists.

Lemma E. 1 The DFvs reduction problem on a triple $(G, F, k)$ is solvable in time $O\left(n^{3} 4^{k} k^{3} k!\right)$, where $n$ is the number of vertices in the input graph $G$.

Proof. Let $G=(V, E)$ be the input directed graph with $n=|V|$ vertices, and let $F$ be the input FVS of size $k+1$ for the graph $G$. Every FVS $F^{\prime}$ of size bounded by $k$ for $G$ can be split into two disjoint subsets $F_{1}$ and $F_{2}$, where $F_{2}$ consists of $j$ vertices in $F$ for some integer $j, 0 \leq j \leq k$, and $F_{1}$ consists of at most $k-j$ vertices in $V-F$. Note that since we assume that no vertex in $F-F_{2}$ is in the FVS $F^{\prime}$, the induced subgraph $G\left[F-F_{2}\right]$ must be a DAG. Therefore, for each $j, 0 \leq j \leq k$, we enumerate all subsets of $j$ vertices in $F$. For each such subset $F_{2}$ of $F$ such that $G\left[F-F_{2}\right]$ is a DAG, we seek a subset $F_{1}$ of at most $k-j$ vertices in $V-F$ such that $F_{1} \cup F_{2}$ makes an FVS for the graph $G$.

Fix a subset $F_{2}$ of $F$, such that $\left|F_{2}\right|=j$ and that the induced subgraph $G\left[F-F_{2}\right]$ is a DAG. Note that the graph $G$ has an FVS $F_{1} \cup F_{2}$ of size bounded by $k$, where $F_{1} \subseteq V-F$, if and only if the subset $F_{1}$ of $V-F$ is an FVS for the graph $G-F_{2}$ and the size of $F_{1}$ is bounded by $k-j$. Therefore, to solve the original problem, we can instead consider how to construct an FVS $F_{1}$ for the graph $G-F_{2}$ such that $\left|F_{1}\right| \leq k-j$ and $F_{1} \subseteq V-F$.

Since $F$ is an FVS for $G$, we have that the induced subgraph $G[V-F]=G-F$ is a DAG. Moreover, by our assumption, the induced subgraph $G\left[F-F_{2}\right]$ is also a DAG. Note that $(V-F) \cup\left(F-F_{2}\right)=V-F_{2}$, which is the vertex set for the graph $G^{\prime}=G-F_{2}$. Therefore, $\left(V-F, F-F_{2}\right)$ is a DAG-bipartition of the graph
$G^{\prime}$. Thus, an FVS $F_{1}$ for the graph $G^{\prime}$ such that $\left|F_{1}\right| \leq k-j$ and $F_{1} \subseteq V-F$, is actually a $(V-F)$-FVS of size bounded by $k-j$ for the graph $G^{\prime}$ with the DAGbipartition $\left(V-F, F-F_{2}\right)$. Therefore, the set $F_{1}$ can be constructed by the algorithm $\operatorname{DBF}\left(G^{\prime}, V-F, F-F_{2}, k-j\right)$.

Since $|F|=k+1$ and $\left|F_{2}\right|=j$, we have $\left|F-F_{2}\right|=k+1-j$. Therefore, the DAG $G\left[F-F_{2}\right]$ contains exactly $k+1-j$ vertices. By Theorem D. 2 , the running time of the algorithm $\operatorname{DBF}\left(G^{\prime}, V-F, F-F_{2}, k-j\right)$ is bounded by $O\left(4^{k-j}(k-j) n^{3}(k+1-j)!\right)$. Now for all integers $j, 0 \leq j \leq k$, we enumerate all subsets $F_{2}$ of $j$ vertices in $F$ and apply the algorithm $\operatorname{DBF}\left(G^{\prime}, V-F, F-F_{2}, k-j\right)$ for those $F_{2}$ such that $G\left[F-F_{2}\right]$ is a DAG. As we discussed above, the graph $G$ has an FVS of size bounded by $k$ if and only if for some $F_{2}$ of $j$ vertices in $F$, where $0 \leq j \leq k$, the algorithm $\operatorname{DBF}\left(G^{\prime}, V-F, F-F_{2}, k-j\right)$ produces an FVS $F_{1}$ of size bounded by $k-j$ for the graph $G^{\prime}$. The running time of this process is bounded by the order of

$$
\sum_{j=0}^{k}\binom{k+1}{j}\left(4^{k-j}(k-j) n^{3}(k+1-j)!\right)=O\left(n^{3} 4^{k} k^{3} k!\right)
$$

This completes the proof of the lemma.
The rest of our process for solving the original DFVS problem is to apply the iterative compression method. The method was proposed by [90] and has been used for solving the FEEDBACK VERTEX SET problem on undirected graphs [33, 57]. Here we extend the method and apply it to solve the DFVS problem.

Theorem E. 2 The DFVs problem is solvable in time $O\left(n^{4} 4^{k} k^{3} k!\right)$.

Proof. Let $(G, k)$ be an instance of the DFVs problem, where $G=(V, E)$ is a directed graph with $n=|V|$ vertices, and $k$ is the parameter. Pick any subset $V_{0}$ of $k+1$ vertices in $G$, and let $F_{0}$ be any subset of $k$ vertices in $V_{0}$. Note that the set $F_{0}$
is an FVS of $k$ vertices for the induced subgraph $G_{0}=G\left[V_{0}\right]$ since the graph $G_{0}-F_{0}$ consists of a single vertex (note that by our assumption, the graph $G$ contains no self-loops).

Let $V-V_{0}=\left\{v_{1}, v_{2}, \ldots, v_{n-k-1}\right\}$. Let $V_{i}=V_{0} \cup\left\{v_{1}, \ldots, v_{i}\right\}$, and let $G_{i}=G\left[V_{i}\right]$ be the subgraph induced by $V_{i}$, for $i=0,1, \ldots, n-k-1$. Inductively, suppose that for an integer $i, 0 \leq i<n-k-1$, we have constructed an FVS $F_{i}$ of size bounded by $k$ for the induced subgraph $G_{i}$ (this has been the case for $i=0$ ). Without loss of generality, we can assume that the set $F_{i}$ consists of exactly $k$ vertices - otherwise we simply pick $k-\left|F_{i}\right|$ vertices (arbitrarily) from $G_{i}-F_{i}$ and add them to the set $F_{i}$. Now consider the set $F_{i+1}^{\prime}=F_{i}+v_{i+1}$. Since $G_{i+1}-F_{i+1}^{\prime}=G_{i}-F_{i}$ and $F_{i}$ is an FVS for $G_{i}$, the set $F_{i+1}^{\prime}$ is an FVS of size $k+1$ for the induced subgraph $G_{i+1}$. In particular, the triple $\left(G_{i+1}, F_{i+1}^{\prime}, k\right)$ is a valid instance for the DFVS REDUCTION problem.

Apply Theorem E. 1 to the instance $\left(G_{i+1}, F_{i+1}^{\prime}, k\right)$, which either returns an FVS $F_{i+1}$ of size bounded by $k$ for the graph $G_{i+1}$, or claims that no such FVS exists. It is easy to see that if the induced subgraph $G_{i+1}=G\left[V_{i+1}\right]$ does not have an FVS of size bounded by $k$, then the original graph $G$ cannot have an FVS of size bounded by $k$. Therefore, in this case, we can simply stop and conclude that there is no FVS of size bounded by $k$ for the original input graph $G$. On the other hand, suppose that an FVS $F_{i+1}$ of size bounded by $k$ is constructed for the graph $G_{i+1}$ in the above process, then the induction successfully proceeds from $i$ to $i+1$ with a new pair $\left(G_{i+1}, F_{i+1}\right)$.

In conclusion, the above process either stops at some point and correctly reports that the input graph $G$ has no FVS of size bounded by $k$, or eventually ends with an FVS $F_{n-k-1}$ of size bounded by $k$ for the graph $G_{n-k-1}=G\left[V_{n-k-1}\right]=G$.

This process is involved in solving at most $n-k-1$ instances $\left(G_{i}, F_{i}, k\right)$ of the DFVS REDUCTION problem, for $0 \leq i \leq n-k-2$. By Theorem E.1, the running time
of the process is bounded by $O\left(n^{3} 4^{k} k^{3} k!(n-k-1)\right)=O\left(n^{4} 4^{k} k^{3} k!\right)$, and the process correctly solves the DFVS problem.

## F. Final Remarks

The running time of the algorithm in Theorem E. 2 can be further improved by taking advantage of existing approximation algorithms for the FEEDBACK VERTEX SET problem on directed graphs. [41] have developed a polynomial time approximation algorithm for the FEEDBACK VERTEX SET problem that for a given directed graph $G$, produces an FVS $F$ of size bounded by $c \cdot \tau \log \tau \log \log \tau$ in time $O\left(n^{2} M(n) \log ^{2} n\right)$, where $c$ is a constant, $\tau$ is the size of a minimum FVS for the graph $G$, and $M(n)=O\left(n^{2.376}\right)$ is the complexity of the multiplication of two $n \times n$ matrices. Therefore, for a given instance $(G, k)$ of the DFVS problem, we can first apply the approximation algorithm in [41] to construct an FVS $F$ for the graph $G$. If $|F|>c \cdot k \log k \log \log k$, then we know that the graph $G$ has no FVS of size bounded by $k$. On the other hand, suppose that $|F| \leq c \cdot k \log k \log \log k$. Then we pick a subset $F_{0}$ of arbitrary $k$ vertices in $F$, and let $G_{0}=G-\left(F-F_{0}\right)$. The set $F_{0}$ is an FVS of size $k$ for the graph $G_{0}$. Now we can proceed exactly the same way as we did in Theorem E.2: let $F-F_{0}=\left\{v_{1}, v_{2}, \ldots, v_{h}\right\}$, where $h \leq c \cdot k \log k \log \log k-k$, and let $V_{i}=V_{0} \cup\left\{v_{1}, \ldots, v_{i}\right\}$, and $G_{i}=G\left[V_{i}\right]$, for $i=0,1, \ldots, h$. By repeatedly applying the algorithm in Lemma E.1, we can either stop with a certain index $i$ where the induced subgraph $G_{i+1}$ has no FVS of size bounded by $k$ (thus the original input graph $G$ has no FVS of size bounded by $k$ ), or eventually construct an FVS $F_{h}$ of size bounded by $k$ for the graph $G_{h}=G\left[V_{h}\right]=G$. This process calls for the execution of the algorithm in Lemma E. 1 at most $h=O(k \log k \log \log k)$ times, and each execution takes time $O\left(n^{3} 4^{k} k^{3} k!\right)$. In conclusion, the DFVS problem can be solved in time
$O\left(n^{3} 4^{k} k^{4} k!\log k \log \log k+n^{4.376} \log ^{2} n\right)$, where the second term in the complexity is due to the approximation algorithm given in [41].

We presented a parameterized algorithm of running time $O\left(n^{4} 4^{k} k^{3} k!\right)$ for the DFVS problem, which shows that the problem is fixed-parameter tractable, and resolves an outstanding open problem in parameterized computation and complexity. Before we close the chapter, we give a few remarks on our results and on directions for future research.

There is an edge version of the FEEDBACK SET problem, which is called the FEEDBACK ARC SET problem (briefly, the DFAS problem): given a directed graph $G$ and a parameter $k$, either construct a set of at most $k$ edges in $G$ whose removal leaves a DAG, or report that no such edge set exists. The DFAS problem is also a well-known NP-complete problem [53]. As shown by [41], the DFAS problem and the DFVS problem can be reduced in linear time from one to the other with the same parameter. Therefore, our results also imply an $O\left(n^{4} 4^{k} k^{3} k!\right)$ time algorithm for the DFAS problem.

The techniques developed in this chapter for solving the SKEW SEPARATOR problem seem to be powerful and generally useful in the study of a variety of separator problems. For example, it has been used recently in developing improved algorithms for a multi-cut problem on undirected graphs in which a separator is sought to (uniformly) separate a set of given terminals [23]. It will be interesting to identify the conditions for the multi-cut problems under which these techniques (and their variations and generalizations) are applicable. In particular, it will be interesting to see if the techniques are applicable to derive the fixed-parameter tractability of the feedback vertex set problem on weighted and directed graphs. Note that the fixed-parameter tractability of the problem on weighted and undirected graphs has been derived recently [18].

It will be interesting to develop new techniques that lead to faster parameterized algorithms for the DFVS problem and other related problems. For example, is it possible that the DFVS problem can be solved in time $O\left(c^{k} n^{O(1)}\right)$ for a constant $c$ ? Another direction is to look at the kernelization of the DFVS problem, by which we refer to a polynomial-time algorithm that on an instance $(G, k)$ of the DFVS problem, produces a (smaller) instance $\left(G^{\prime}, k^{\prime}\right)$ of the problem, such that the size of the graph $G^{\prime}$ (the kernel) is bounded by a function $g(k)$ of $k$ (but independent of the size of the original graph $G$ ), that $k^{\prime} \leq k$, and that the graph $G$ has an FVS of size bounded by $k$ if and only if the graph $G^{\prime}$ has an FVS of size bounded by $k^{\prime}$. Since now it is known that the DFVS problem is fixed-parameter tractable, by a general theorem in parameterized complexity theory [37], such a kernelization algorithm exists for the DFVS problem. However, how small can the size of the kernel $G^{\prime}$ be? In particular, can the kernel $G^{\prime}$ have its size bounded by a polynomial of the parameter $k$ ? We note that recently there has been progress in the study of kernelization for the FEEDBACK VERTEX SET problem on undirected graphs. [7] was able to give a kernel of size $O\left(k^{3}\right)$ for the FEEDBACK VERTEX SET problem on undirected graphs, and [8] have shown that the feedback vertex set problem on undirected planar graphs has a kernel of size $O(k)$.

Algorithm $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k\right)$
input: an instance ( $G,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k$ ) of the SKEW SEPARATOR problem.
output: a skew separator of size bounded by $k$ for the pair ( $\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right]$ ), or report "No" (i.e., no such separator exists).

1. if $l=1$ then solve the problem in time $O\left(k n^{2}\right)$;
2. if Rule R2 applies or $k<0$ then return "No";
3. if Rule R1 applies then return $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l-1}\right],\left[T_{1}, \ldots, T_{l-1}\right], k\right)$;
4. if Rule R3 applies on a vertex $w$
then return $\{w\} \cup \operatorname{SMC}\left(G-w,\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{]}\right], k-1\right) ;{ }^{\S}$
5. pick an $S_{l}$-extended vertex $u_{0}$; let $S_{l}^{\prime}=S_{l} \cup\left\{u_{0}\right\}$;
6. let $m$ be the size of a min-cut from $S_{l}$ to $T_{\text {all }}=\bigcup_{i=1}^{l} T_{i}$;
7. if $m>k$ then return "No";
8. let $m^{\prime}$ be the size of a min-cut from $S_{l}^{\prime}$ to $T_{\text {all }}$;
9. if $\left(m=m^{\prime}\right)$
9.1. then return $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T_{l}\right], k\right)$;
9.2. else $X=\left\{u_{0}\right\} \cup \operatorname{SMC}\left(G-u_{0},\left[S_{1}, \ldots, S_{l}\right],\left[T_{1}, \ldots, T_{l}\right], k-1\right)$; if $X \neq$ "No" then return $X$;
9.3. else return $\operatorname{SMC}\left(G,\left[S_{1}, \ldots, S_{l-1}, S_{l}^{\prime}\right],\left[T_{1}, \ldots, T_{l-1}, T\right], k\right)$.
§ For simplicity, we assume that a "No" plus anything gives a "No".

Fig. 8. An algorithm for the sKew separator problem.

Algorithm $\operatorname{DBF}\left(G, D_{1}, D_{2}, k\right)$
input: an instance ( $G, D_{1}, D_{2}, k$ ) of the DAG-BIPARTITION FVS problem.
output: a $D_{1}$-FVS of size bounded by $k$ for $G$, or report "No"
(i.e., no such $D_{1}$-FVS exists).

1. for each topologically sorted order $\pi=\left\{v_{1}, \ldots, v_{h}\right\}$ of the vertices in $G\left[D_{2}\right]$ do
1.1. construct the instance $\left(G_{\pi},\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right], k\right)$ of the SKEW SEPARATOR problem induced by $\left(G, D_{1}, D_{2}, k\right)$ and $\pi$;
1.2. let $X=\mathbf{S M C}\left(G_{\pi},\left[\left\{s_{1}\right\}, \ldots,\left\{s_{h}\right\}\right],\left[\left\{t_{1}\right\}, \ldots,\left\{t_{h}\right\}\right], k\right)$;
1.3. $\quad$ if $X$ is a $D_{1}$-FVS of size bounded by $k$ for $G$ then return $(X)$; stop;
2. return("No").

Fig. 9. An algorithm for the DAG-BIPARTITION FVS problem.

## CHAPTER VI

## MAX-LEAF

In this chapter, we give an fpt-algorithm of running time $O^{*}\left(4^{k}\right)$ for the MAX-LEAF OUT-BRANCHING problem on directed graphs, thus improving the previous best algorithm of running time $O^{*}\left(2^{O\left(k^{3} \log k\right)}\right)$ [12]. The algorithm also works for the MAX-LEAF OUT-BRANCHING problem on undirected graphs, thus improving the previous best algorithm of running time $O^{*}\left(6.75^{k}\right)$ for the MAX-LEAF SPANNING-TREE problem on undirected graphs [14].

As we do in previous chapters, we consider two measures for the MAX-LEAF OUT-BRANCHING problem on directed graphs. Our algorithm extends an out-tree gradually into an out-branching $T_{s}$. During the process of extending an out-tree, we divide the leaves of the out-tree into two subsets $L_{1}$ and $L_{2}$ such that leaves in $L_{1}$ will reach exactly one leaf in $T_{s}$ and the leaves may reach more than one leaves in $T_{s}$. The sizes of $L_{1}$ and $L_{2}$ are the measures for our algorithm. These two measures are indirect measures, since both are not the parameter-the number of leaves in an out-branching-for the MAX-LEAF OUT-BRANCHING problem.

## A. Introduction

The MAX-LEAF SPANNING-TREE problem is to find a spanning tree with the maximum number of leaves in an undirected graph. This problem has a version for directed graphs (i.e., digraphs), which is the max-LEAF OUT-BRANCHING problem. These problems are of importance in many theoretical and practical applications [56, 97, 102]. The problems are NP-hard [53]. In terms of polynomial time approximability, the MAX-LEAF SPANNING-TREE problem is APX-hard, and can be approximated with a ratio 2 in polynomial time [95], while very recent research shows that the

MAX-LEAF OUT-BRANCHING problem can be approximated with a ratio $O(\sqrt{n})$ in polynomial time [40].

These problems have been intensively studied $[1,2,9,11,12,14,43,79,95]$. For the parameterized MAX LEAF problem on undirected graphs, there is a chain of improved algorithms: $O\left(\left(17 k^{4}\right)!\right)$ [9], then $O\left((2 k)^{4 k}\right)$ [36], then $O\left(14.23^{k}\right)$ [43], then $O\left(9.49^{k}\right)$ [11] and finally $O\left(6.75^{k}\right)$ [14]. For the parametrized max leaf problem on digraph, algorithm of time $O\left(2^{O\left(k^{2} \operatorname{logk}\right)}\right)$ was given for strongly connected digraph and acyclic digraph [2], which was improved to $O\left(2^{k \log ^{2} k}\right)$ for strongly connected digraph and $O\left(2^{k l o g k}\right)$ for acyclic digraph. The parameterized max LEAF problem on general digraph was proposed as an open problem in [1], [2], and [58], and was shown to be fixed-parameter tractable with an algorithm of time $O\left(2^{O\left(k^{3} \log k\right)}\right)$ [12].

The more general max-LEAF OUT-BRANCHING problem seems more difficult. The problem has become significantly interesting very recently because of its rich combinatorial structures and challenging algorithmic techniques. In this chapter, we present a simple branch-and-search algorithm of running time $O^{*}\left(4^{k}\right)$ that solves the MAX-LEAF OUT-BRANCHING problem. This result significantly improves the previous best algorithm for the problem that runs in time $O^{*}\left(2^{O(k \log k)}\right)$ [13]. The MAX-LEAF SPANNING-TREE problem can be trivially reduced to the MAX-LEAF OUT-BRANCHING problem if we convert the input undirected graph $G$ into a digraph by replacing each undirected edge $[u v]$ by two directed edges $u v$ and $v u$. Therefore, our algorithm of time $O^{*}\left(4^{k}\right)$ can be used to solve the MAX-LEAF SPANNING-TREE problem, improving the previous best algorithm of running time $O^{*}\left(6.75^{k}\right)$ [14].

## B. Preliminaries

All graphs in our discussion are simple graphs, i.e., there are no multiple edges from a vertex to another vertex, and no self-loops on any vertex. For an edge $x y$ in a digraph $G$, the vertex $x$ is called an in-neighbor of the vertex $y$, and the vertex $y$ is called an out-neighbor of the vertex $x$. The in-degree of a vertex $x$ in $G$ is the number of in-neighbors of $x$, and the out-degree of a vertex $x$ is the number of out-neighbors of $x$.

A path $P$ in a digraph $G$ from a vertex $x_{1}$ to another vertex $x_{q}$ is a sequence of vertices $P=x_{1} \cdots x_{q}$ such that $x_{i} x_{i+1}$ is an edge in $G$ for all $1 \leq i \leq q-1$. The path $P$ is simple if no vertex is repeated in $P$. A vertex $x$ can reach another vertex $y$ (or equivalently, the vertex $y$ is reachable from the vertex $x$ ) in a digraph $G$ if there is a path from $x$ to $y$ in $G$.

An acyclic digraph $T$ is an out-tree (rooted at a vertex $r$ ) if the vertex $r$ has in-degree 0 and for every vertex $x$ in $T$ there is a unique path from $r$ to $x$. The vertex $r$ is the root of the out-tree $T$, and each vertex of out-degree 0 is a leaf of $T$. A vertex $x$ in $T$ is an internal vertex if it is not a leaf. A $k$-out-tree is an out-tree that has at least $k$ leaves. An out-tree in a digraph $G$ is a subgraph of $G$ that is an out-tree. An out-branching of a digraph $G$ is an out-tree of $G$ that contains all vertices of $G$. A $k$-out-branching is an out-branching that has at least $k$ leaves.

Let $T$ be an out-tree in a digraph $G$, and let $x_{0}$ be a leaf of $T$. An out-chain $\pi(l)$ of $T$ is a simple path $x_{0} x_{1} \ldots x_{q}$ in $G$ such that for all $0 \leq i \leq q-1, x_{i+1}$ is the only out-neighbor of $x_{i}$ not in $T \cup\left\{x_{1}, \ldots, x_{i}\right\}$. A $y$-truncated out-chain $\pi_{y}(l)$ of an out-chain $\pi(l)$, where $y \neq l$, is defined as follows: (1) if $y \in \pi(l)$, then $\pi_{y}(l)$ is the path from $l$ to $z$ in $\pi(l)$, where $z$ is the in-neighbor of $y$ in $\pi(l)$; and (2) if $y \notin \pi(l)$, then $\pi_{y}(l)=\pi(l)$. Inductively, $\pi_{y_{1}, \ldots, y_{q}}(l)$ is the $y_{q}$-truncated out-chain of $\pi_{y_{1}, \ldots, y_{q-1}}(l)$.

For a given set $Q$ of out-chains of the out-tree $T$, a leaf $l$ of $T$ is determined if an out-chain $\pi(l)$ is included in $Q$. An ordered set $Q$ of out-chains $\pi\left(l_{1}\right), \ldots, \pi\left(l_{q}\right)$ is consistent if $\pi\left(l_{i}\right)$ and $\pi\left(l_{j}\right)$ contain a common vertex $z$, then the path from $z$ to the last vertex of $w \pi\left(l_{i}\right)$ is also in $\pi\left(l_{j}\right)$ for all $1 \leq i<j \leq q$.

Let $T$ be an out-tree in a digraph $G$ with a given consistent set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}$ of out-chains, and let $k$ be a positive integer. A $k$-out-branching $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ such that (1) $T_{s}$ and $T$ have the same root, (2) $T$ is a subgraph of $T_{s}$, and (3) all vertices reachable in $T_{s}$ from $l_{i}$ are in $\pi\left(l_{i}\right)$. We formulate the following problem.

EXtending max-Leaf: given a digraph $T$, an out-tree $T$ in $G$ with a consistent set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{q}\right)\right\}$ of out-chains, and a parameter $k$, either construct an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$, or report 'No' if no such an out-branching exists.

In the rest of this chapter, $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ also refers to the instance of the EXTENDING MAX-LEAF problem with inputs of an out-tree $T$, a consistent set of out-chains $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}$ for $T$, and an integer $k$.

The following reduction will play an important role in our discussion.

Definition An instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to another instance $\left(T^{\prime} ; \pi^{\prime}\left(l_{1}\right), \ldots, \pi^{\prime}\left(l_{q}\right) ; k^{\prime}\right)$ if
(1) any extended out-branching for $\left(T^{\prime} ; \pi^{\prime}\left(l_{1}\right), \ldots, \pi^{\prime}\left(l_{q}\right) ; k^{\prime}\right)$ is an extended outbranching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$;
(2)if there is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$, there is an extended out-branching for $\left(T^{\prime} ; \pi^{\prime}\left(l_{1}\right), \ldots, \pi^{\prime}\left(l_{q}\right) ; k^{\prime}\right)$.

It is easy to verify that the above reduction is a transitive relation.
The following notations will be used to simplify our discussion.

- $T+x y_{1}+\ldots+x y_{q}$ is the the resulting graph after adding edge $x y_{1}, \ldots, x y_{q}$ to $T$ and including vertex $y_{i}$ in $T+x y$ if $y_{i}$ is not in $T$ for all $1 \leq i \leq q$.
- $T-z y$ is the resulting graph after removing edge $z y$ from $T$ and keeping all vertices of $T$ in $T-z y$.
- $G-T$ is the resulting graph after removing all edges between vertices of $T$ from $G$.

In the rest of this section, we give two lemmas on out-trees and out-branchings, which will be used in our later discussion.

Lemma B. 1 Let $T$ and $T^{\prime}$ be two out-trees in a digraph with the same root such that $T^{\prime}$ is a subgraph of $T$, and let $l$ be a leaf of $T^{\prime}$ and $y \neq l$ be a vertex in $T^{\prime}$. Then $l$ cannot reach $y$ in $T$.

Proof. Since $y$ is a vertex of $T^{\prime}$ and $T^{\prime}$ is an out-tree, there is a unique path $P_{1}$ in $T^{\prime}$ from the root $r$ of $T^{\prime}$ to $y$, which obviously does not pass through the leaf $l$ because $l$ has out-degree 0 in $T^{\prime}$. Since $T^{\prime}$ is a subgraph of $T, P_{1}$ is also a path from $r$ to $y$ in $T$. Now if $l$ can reach $y$ in $T$, then the unique path in $T$ from $r$ to $l$ followed by the path from $l$ to $y$ in $T$ forms a second path $P_{2}$ from $r$ to $y$ in $T$ that passes through the vertex $l$. This contradicts the assumption that $T$ is an out-tree since $r$ is also the root of $T$.

Lemma B. 2 Let $T_{s}$ be an out-branching and $T$ be a $k$-out-tree in a digraph such that $T_{s}$ and $T$ have the same root $r$ and $T$ is a subgraph of $T_{s}$. Then $T_{s}$ is a k-outbranching.

Proof. Let $\left\{l_{1}, \ldots, l_{p}\right\}$ be the set of all leaves of $T$. For each $1 \leq i \leq p$, let $L_{i}$ be the set of leaves of $T_{s}$ that can be reached from $l_{i}$ in $T_{s}$. Note that each set $L_{i}$ contains at least one vertex in $T_{s}$.

We prove that no two sets $L_{i}$ and $L_{j}$ share a common vertex. Suppose that $w \in L_{i} \cap L_{j}$. Then the unique path $P_{i}^{\prime}$ in $T$ (which is also the unique path in $T_{s}$ ) from the root $r$ to $l_{i}$ followed by a path $P_{i}^{\prime \prime}$ from $l_{i}$ to $w$ in $T_{s}$ forms a path $P_{i}$ in $T_{s}$ from the root $r$ to $w$. Similarly, there is a path $P_{j}$ in $T_{s}$ from the root $r$ to $w$ that passes through the vertex $l_{j}$. Since $T_{s}$ is an out-tree, we must have $P_{i}=P_{j}$. In consequence, the path $P_{i}$ in $T_{s}$ contains the vertex $l_{j}$. Since both $l_{i}$ and $l_{j}$ are leaves in $T$, the unique path $P_{i}^{\prime}$ in $T$ from $r$ to $l_{i}$ does not contain the vertex $l_{j}$. Therefore, the vertex $l_{j}$ must be contained in the path $P_{i}^{\prime \prime}$ from $l_{i}$ to $w$, so $l_{i}$ can reach $l_{j}$ in $T_{s}$. However, this is impossible according to Lemma B.1.

Therefore, $T_{s}$ has at least $\left|L_{1}\right|+\cdots+\left|L_{p}\right| \geq p$ leaves. Since $T$ is a $k$-out-tree, $p \geq k$. Thus, $T_{s}$ is a $k$-out-branching.

## C. Extending an Out-tree

Throughout the discussion of this section, we let $T$ be an out-tree with the root $r$ in a digraph $G$, a given consistent set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}$ of out-chains for $T$. Let $x$ be a vertex in $T$ that is not a determined leaf of $T$ and let $y$ be an out-neighbor of $x$ in $G-T$. In subsection 1, we discuss some properties that are useful for reducing the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ to the instance $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$. In subsection 2, we show that the results in subsection 1 can be generalized to the case where many out-neighbors of $x$ in $G-T$ are considered. Furthermore, if $T$ is already a $k$-out-tree, then we can solve the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ in polynomial time, which gives a boundary condition for our algorithm.

## 1. Properties for Extending an Out-tree

We first verify that $T+x y$ is an out-tree in $G$ and that the $y$-truncated out-chains $\pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right)$ are a consistent set of out-chains for $T+x y$. This is necessary before we can discuss the extendability of $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$.

Lemma C. $1 T+x y$ is an out-tree with root $r$ in the digraph $G$.
Proof. Since the vertex $y$ has out-degree 0 in $T+x y$, it cannot help creating any new path in $T+x y$ from the root $r$ to any vertex $w \neq y$ in $T+x y$. Thus, there is a unique path from the root $r$ to $w$ in $T+x y$. Moreover, since $x$ is the only in-neighbor of $y$ in $T+x y$, the only path from the root $r$ to $y$ in $T+x y$ must be the unique path from $r$ to $x$ followed by the edge $x y$. Finally, $y$ is of in-degree one, $r$ is of in-degree zero, and any other vertex $w \neq r$ of $T$ is of in-degree one in $T+x y$. Therefore, $T+x y$ is an out-tree.

Lemma C. 2 The truncated out-chains $\pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right)$ is a consistent set of outchains for $T+x y$.

Proof. Pick any out-chain of $T: \pi\left(l_{i}\right)=\left\{u_{0}, u_{1}, \ldots, u_{q}\right\}$, where $u_{0}=l_{i}$. Then $u_{j}$ has a unique out-neighbor $u_{j+1}$ in $G-\left(T \cup\left\{u_{0}, \ldots, u_{j}\right\}\right)$ for all $0 \leq j \leq q-1$. Thus, and by that $T$ is a subgraph of $T+x y$ and the definition of out-chain, $u_{j+1}$ is the unique out-neighbor in $G-\left(T+x y+u_{0} \cdots u_{j}\right)$ for any $u_{j} \in \pi_{y}\left(l_{i}\right)$ except the last vertex in $\pi_{y}\left(l_{i}\right)$. So $\pi_{y}\left(l_{i}\right)$ is an out-chain of $T+x y$ for all $1 \leq i \leq p$.

It is given that $\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)$ is a consistent set of out-chains for $T$, i.e., if $z$ is in both $\pi\left(l_{i}\right)$ and $\pi\left(l_{j}\right)$ where $i<j$, then the path from $z$ to the last vertex $w$ of $\pi\left(l_{i}\right)$ is in $\pi\left(l_{j}\right)$. Let $z$ be a common vertex of $\pi_{y}\left(l_{i}\right)$ and $\pi_{y}\left(l_{j}\right)$ where $1 \leq i<j \leq q$. If $y \neq w$ is in the path from $w$ to the last vertex of $\pi\left(l_{j}\right)$, then $\pi_{y}\left(l_{i}\right)=\pi\left(l_{i}\right)$ and $\pi_{y}\left(l_{j}\right)$ is the path from $l_{j}$ to the in-neighbour of $y$ in $\pi_{y}\left(l_{j}\right)$, which also includes the path from $z$ to $w$ in $\pi\left(l_{i}\right)$. That is, the path from $z$ to $w$ in $\pi_{y}\left(l_{i}\right)$ is also in $\pi_{y}\left(l_{j}\right)$. If $y$ is
in the path from $z$ to $w$ in $\pi\left(l_{i}\right)$, let $p$ be the in-neighbour of $y$ in $\pi\left(l_{i}\right)$. Then $\pi_{y}\left(l_{i}\right)$ is the path from $l_{i}$ to $p$ in $\pi\left(l_{i}\right)$, and $\pi_{y}\left(l_{j}\right)$ is the path from $l_{j}$ to $p$ in $\pi\left(l_{i}\right)$. Since the path from $z$ to $w$ is in both $\pi\left(l_{i}\right)$ and $\pi\left(l_{j}\right)$, and $y$ is in the path from $z$ to $w$, the path from $z$ to $p$ is in both $\pi_{y}\left(l_{i}\right)$ and $\pi_{y}\left(l_{j}\right)$. If $y \neq z$ is either in the path from $l_{i}$ to $z$ in $\pi\left(l_{i}\right)$ or in the path from $l_{j}$ to $z$ in $\pi\left(l_{j}\right), \pi_{y}\left(l_{i}\right)$ and $\pi_{y}\left(l_{j}\right)$ have no common vertex. Finally, if $y$ is not in $\pi\left(l_{i}\right)$ nor $\pi\left(l_{j}\right), \pi_{y}\left(l_{i}\right)=\pi\left(l_{i}\right)$ and $\pi_{y}\left(l_{j}\right)=\pi\left(l_{j}\right)$. Thus, the path from $z$ to the last vertex of $\pi_{y}\left(l_{i}\right)$ in $\pi_{y}\left(l_{i}\right)$ is in $\pi_{y}\left(l_{j}\right)$. In summary, if $\pi_{y}\left(l_{i}\right)$ and $\pi_{y}\left(l_{j}\right)$ have a common vertex $z$, the path from $z$ to the last vertex of $\pi_{y}\left(l_{i}\right)$ is in $\pi_{y}\left(l_{j}\right)$ for all $1 \leq i \leq j$. Therefore, the set $\left\{\pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right)\right\}$ of truncated out-chains is consistent.

Lemma C. 3 Any extended out-branching for $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$.

Proof. Let $T_{s}$ be an extended out-branching for $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$. By the definition we have: (1) $T_{s}$ is a $k$-out-branching; (2) $T+x y$ and $T_{s}$ has the same root $r$; (3) $T+x y$ is a subgraph of $T_{s}$; and (4) all vertices reachable in $T_{s}$ from $l_{i}$ are in $\pi_{y}\left(l_{i}\right)$. We show $T_{s}$ is also an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$.

By Lemma C.1, $T$ and $T+x y$ have the same root. Thus, by (2), $T$ and $T_{s}$ have the same root. Since $T$ is a subgraph of $T+x y$, by (3) $T$ is also a subgraph of $T_{s}$. Finally, since $\pi_{y}\left(l_{i}\right) \subseteq \pi\left(l_{i}\right)$, by (4), all vertices reachable in $T_{s}$ from $l_{i}$ are in $\pi\left(l_{i}\right)$. In conclusion, $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$.

In the rest of this section, we assume that $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. Let $z$ be the (unique) in-neighbor of $y$ in $T_{s}$. We will prove that $T_{s}-z y+x y$ is an extended out-branching for $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$. We first prove the following properties.

Lemma C. $4 T+x y$ is a subgraph of $T_{s}-z y+x y$.
Proof. Since $T$ is a subgraph of $T_{s}$, and $y \notin T$ so the edge $z y$ is not in $T$, all edges in $T$ are in $T_{s}-z y+x y$. Moreover, the edge $x y$ is contained in $T_{s}-z y+x y$. Therefore, all edges in $T+x y$ are in $T_{s}-z y+x y$ i.e., the out-tree $T+x y$ is a subgraph of $T_{s}-z y+x y$.

Lemma C. $5 T_{s}-z y+x y$ is an out-branching of $G$ with root $r$. Moreover, if $x$ is not a leaf of $T$, then $T_{s}-z y+x y$ is an $k$-out-branching of $G$ with root $r$.

Proof. Since $T_{s}$ is an out-branching of the digraph $G$, by the definition, $T_{s}-z y+x y$ is also rooted at $r$ and contains all vertices in $G$. For any vertex $w$ in $G$, if the unique path $P_{w}$ in $T_{s}$ from the root $r$ to $w$ does not pass through the vertex $y$, then $P_{w}$ is also a path in $T_{s}-z y+x y$. On the other hand, if the path $P_{w}$ from $r$ to $w$ is a concatenation of a path $P_{w}^{\prime}$ from $r$ to $y$ and a path $P_{w}^{\prime \prime}$ from $y$ to $w$, then the path in $T$ from $r$ to $x$ (note that $z y \notin T$ ) followed by the edge $x y$ then by the path $P_{w}^{\prime \prime}$ ( $z y \notin P_{w}$ " since $y$ is of in-degree one) makes a path in $T_{s}-z y+x y$ from $r$ to $w$. The uniqueness of the path in $T_{s}-z y+x y$ from $r$ to $w$ can be easily verified because each vertex in $T_{s}-z y+x y$ has in-degree bounded by 1 .

If $x$ is not a leaf of $T_{s}$, then deleting the edge $z y$ then adding the edge $x y$ can not decrease the number of vertices of degree 0 in the out-branching. Thus, $T_{s}-z y+x y$ has at least as many leaves as that of $T_{s}$. Since $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right), T_{s}$ has at least $k$ leaves. In consequence, $T_{s}-z y+x y$ is a $k$-out-branching.

Lemma C. 6 The vertices reachable in $T_{s}-z y+x y$ from a determined leaf $l_{i}$ are all in $\pi_{y}\left(l_{i}\right)$.

Proof. We first prove that the vertices reachable in $T_{s}-z y+x y$ from the determined leaf $l_{i}$ are all in $\pi\left(l_{i}\right)$. Suppose that $l_{i}$ can reach a vertex $w \notin \pi\left(l_{i}\right)$ in
$T_{s}-z y+x y$ via a path $P$. Then the edge $x y$ must be in $P$. Otherwise, the path $P$ is also in $T_{s}$ and $l_{i}$ can reach $w \notin \pi\left(l_{i}\right)$ in $T_{s}$, contradicting the assumption that $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. Thus, $l_{i}$ can reach $x$ in $T_{s}$, which is impossible by Lemma B. 1 because both $l_{i}$ and $x$ are in $T$, and $l_{i}$ is a determined leaf of $T$ but $x$ is not a determined leaf. This contradiction proves that all vertices reachable in $T_{s}-z y+x y$ from $l_{l}$ must be in $\pi\left(l_{i}\right)$.

Now we prove that the vertices reachable in $T_{s}-z y+x y$ from $l_{i}$ are all in $\pi_{y}\left(l_{i}\right)$. If $y \notin \pi\left(l_{i}\right)$, then $\pi_{y}\left(l_{i}\right)=\pi\left(l_{i}\right)$, and the lemma is proved by the above analysis. Suppose $y \in \pi\left(l_{i}\right)$ and $l_{i}$ can reach a vertex $w \in \pi\left(l_{i}\right)-\pi_{y}\left(l_{i}\right)$ via a path $P^{\prime}$. Since $l_{i}$ is a leaf of $T$, by Lemma B. 1 any vertex $u \neq l_{i}$ of $T$ cannot be in $P^{\prime}$. Thus, $\pi_{y}\left(l_{i}\right) \cup\{y\}$ must be in $P^{\prime}$. However, since (1) both $l_{i}$ and $y$ are in $T+x y,(2) l_{i}$ is a leaf of $T+x y$, (3) $T_{s}-z y+x y$ have the same root as $T$ by Lemmas C. 1 and C.5, and (4) $T+x y$ is a subgraph of $T_{s}-z y+x y$ by Lemma C.4, therefore, $l_{i}$ can not reach $y$ in $T_{s}-z y+x y$ by Lemma B.1. Thus, all vertices reachable from $l_{i}$ in $T_{s}-z y+x y$ are in $\pi_{y}\left(l_{i}\right)$ when $y \in \pi\left(l_{i}\right)$.

## 2. Extending an Out-tree

Now, we are ready to show that we can extend $T+x y$ instead of $T$. Note that $T+x y$ can be found in polynomial time, and this can occur at most $n$ times because $T+x y$ has one more vertex than $T$. Thus, this is an efficient operation without branching.

Lemma C. 7 If $x$ is an internal vertex of $T$, the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)^{1}$.

Proof. Let $T_{s}$ be an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. By

[^5]Lemma C.5, $T_{s}-z y+x y$ is a $k$-out-branching with root $r$ where $z$ is the unique in-neighbour of $y$ in $T_{s}$. Thus $T_{s}-z y+x y$ and $T+x y$ have the same root according to Lemma C.1, and $T+x y$ is a subgraph of $T_{s}-z y+x y$ by Lemma C.4. Finally, all vertices reachable in $T_{s}-z y+x y$ from a determined leaf $l_{i}$ are in $\pi_{y}\left(l_{i}\right)$ by Lemma C.6. In consequence, $T_{s}-z y+x y$ is an extended out-branching for $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$.

By Lemma C.3, any extended out-branching for $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. This completes the proof of the lemma.

Lemma C. 7 can be generalized to the case where we add many out-neighbors of the vertex $x$ in $G-T$ to the out-tree $T$. For this, let $y_{1}, \ldots, y_{q}$ be out-neighbors of $x$ in $G-T$.

Lemma C. 8 If $x$ is an internal vertex of $T$, the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y_{1}+\ldots+x y_{q} ; \pi_{y_{1}, \ldots, y_{q}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{q}}\left(l_{p}\right) ; k\right)$.

Proof. We prove the lemma by induction on $q$. The initial case $q=1$ is proved by Lemma C.7.

Now suppose that when $q=i$, the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y_{1}+\ldots+x y_{i} ; \pi_{y_{1}, \ldots, y_{i}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{i}}\left(l_{p}\right) ; k\right)$. For $q=i+1$, note that $T+x y_{1}+\ldots+x y_{i+1}$ is the out-tree $T+x y_{1}+\ldots+x y_{i}$ plus the edge $x y_{i+1}$ where $y_{i+1} \notin T+x y_{1}+\ldots+x y_{i}$, and the truncated out-chain $\pi_{y_{1}, \ldots, y_{i}, y_{i+1}}\left(l_{j}\right)$ is a $y_{i+1}$-truncated out-chain of the out-chain $\pi_{y_{1}, \ldots, y_{i}}\left(l_{j}\right)$ for the tree $T+x y_{1}+\ldots+x y_{i}$ for all $j$. Moreover, the vertex $x$ is obviously an internal vertex of the tree $T+x y_{1}+\ldots+x y_{i}$. Therefore, by Lemma C. 7 again, the instance $\left(T+x y_{1}+\ldots+x y_{i} ; \pi_{y_{1}, \ldots, y_{i}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{i}}\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y_{1}+\ldots+x y_{i+1} ; \pi_{y_{1}, \ldots, y_{i+1}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{i+1}}\left(l_{p}\right) ; k\right)$. Now the transitivity of the reduction proves that the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$
is reducible to the instance $\left(T+x y_{1}+\ldots+x y_{i+1} ; \pi_{y_{1}, \ldots, y_{i+1}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{i+1}}\left(l_{p}\right) ; k\right)$, and the induction goes through.

The conditions in the previous lemmas can be further relaxed without requiring that the vertex $x$ be an internal vertex of $T$, if the out-tree $T$ is a $k$-out-tree, as shown in the following lemmas.

Lemma C. 9 Suppose that $x$ is not a determined leaf of $T$ and $T$ is a $k$-out-tree. Then $T+x y$ is a $k$-out-tree and the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$.

Proof. By Lemma C.1, $T+x y$ is an out-tree with root $r$. The number of leaves of $T+x y$ cannot be less than that of $T$ because adding the edge $x y$ to $T$ adds a vertex $y$ of out-degree 0 and may change the out-degree of at most one vertex in $T$ (i.e., $x$ ). Since $T$ is a $k$-out-tree, $T+x y$ is also a $k$-out-tree. Next we show that the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$.

By Lemma C.3, any extended out-branching for $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. Now suppose that there is an extended out-branching $T_{s}$ for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. By Lemma C.5, the graph $T_{s}-z y+x y$ is an out-branching. Since $T$ is a subgraph of $T_{s}-z y+x y$ by Lemma C. 4 and $T$ is a $k$-out-tree, $T_{s}-z y+x y$ is a $k$-out-branching by Lemma B.2. Also $T_{s}-z y+x y$ and $T+x y$ have the same root $r$ according to Lemma C. 1 and Lemma C.5. Moreover, $T+x y$ is a subgraph of $T_{s}-z y+x y$ by Lemma C.4. Finally, all vertices reachable in $T_{s}-z y+x y$ by a determined leaf $l_{i}$ are in $\pi_{y}\left(l_{i}\right)$ by Lemma C.6. So $T_{s}-z y+x y$ is a $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$-extended out-branching.

Therefore, $T+x y$ is a $k$-out-tree with root $r$ and the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left(T+x y ; \pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right) ; k\right)$.

Similarly, we can generalize Lemma C. 9 to many out-neighbors of the vertex $x$. For this, again let $y_{1}, \ldots, y_{q}$ be out-neighbors of $x$ in $G-T$.

Lemma C. 10 When $T$ is a $k$-out-tree, the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left.\left(T+x y_{1}+\ldots+x y_{q}\right) ; \pi_{y_{1}, \ldots, y_{q}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{q}}\left(l_{p}\right) ; k\right)$, if $x$ is not a determined leaf of $T$.

Proof. The lemma is proved by induction on $q$. The initial case $q=1$ is proved by Lemma C.9, and the inductive step is proved by Lemma C. 9 and the transitivity of the reduction.

Lemma C. 10 is very useful since it implies that by continuously adding outneighbors of a undetermined leaf, if $T$ is already a $k$-out-tree, we can find an extended out-branching for $\left(T, ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. This will be discussed in the following algorithm and its proof.

Theorem C. 11 For a given $k$-out-tree $T$ in a digraph $G$ of $n$ vertices, with a set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}$ of out-chains for $T$, the algorithm $\operatorname{SpanTree}\left(G, T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ runs in polynomial time and (1) either constructs an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ if such one exists, or (2) reports 'No' otherwise.

Proof. We first prove the correctness of the algorithm by showing that (1) the instance $\left(T_{1} ; \pi_{1}\left(l_{1}\right), \ldots, \pi_{1}\left(l_{q}\right) ; k\right)$ before step 2 is reducible to $\left(T_{2} ; \pi_{2}\left(l_{1}\right), \ldots, \pi_{2}\left(l_{q}\right) ; k\right)$, the instance before step 3, (2) step 3 returns 'No' correctly, (3) step 4 construct an extended out-branching for the instance $\left(T_{4} ; \pi_{4}\left(l_{1}\right), \ldots, \pi_{4}\left(l_{q}\right) ; k\right)$ before step 4 , and (4) step 5 returns an extended out-branching for the instance $\left(T_{4} ; \pi_{4}\left(l_{1}\right), \ldots, \pi_{4}\left(l_{q}\right) ; k\right)$.

Claim $S$ does not contain any determined leaf of $T$. Moreover, any vertex of $T$ which have out-neighbours in $G-T$ is either in $S$ or a determined leaf of $T$.

We prove this Claim by induction on the number of iterations of the while loop in step 2.

```
SpanTree \(\left(G, T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)\)
InPUT: a \(k\)-out-tree \(T\) in a digraph \(G\) with a \(\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}\) of
    out-chains for \(T\), and a parameter \(k\)
OUTPUT: an extended out-branching for \(\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)\) if
    there is such one, or "No" otherwise
1. \(S=T-\left\{l_{1}, \ldots, l_{q}\right\}\);
2. while \(S \neq \emptyset\) do
        pick any vertex \(x \in S\);
2.1 if \(x\) has out-neighbors \(y_{1}, \ldots, y_{q}\) in \(G-T\) then
            for \(i=1\) to \(p\) do \(\pi\left(l_{i}\right)=\pi_{y_{1}, \ldots, y_{q}}\left(l_{i}\right)\);
            \(T=T+x y_{1}+\ldots+x y_{q} ;\)
            \(S=\left(S \cup\left\{y_{1}, \ldots, y_{q}\right\}\right)-\{x\} ;\)
        else \(S=S-\{x\}\);
3. if there is a vertex \(x \notin T \cup\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}\) then return 'No';
4. for \(i=p\) to 1 do
        \(x=l_{i}\);
\(4.1 \quad\) while \(x\) has an out-neighbor in \(\pi\left(l_{i}\right)\) do
            \(T=T+x y ; \quad x=y ;\)
5. return \(T\).
```

Fig. 10. Algorithm for the SPAN $k$-OUT-TREE problem

Before the first iteration of the while loop in step 2, $S$ contains all vertices of $T$ which are not determined leaves of $T$, according to step 1 . Thus $S$ contains all possible vertices of $T$ which are not determined leaves of $T$ and have out-neighbours in $G-T$.

Now assume that $S$ does not contain any determined leaf of $T$, and any vertex of $T$ which has out-neighbours in $G-T$ is either in $S$ or a determined leaf of $T$ before a iteration of the while loop at step 2 . We need to show that this is true before the next iteration of the while loop in step 2.

In step 2.1, only $y_{1}, \ldots, y_{q}$, which are not in $T$, are added to $T$, and no out-chains of $y_{1}, \ldots, y_{q}$ for $T+x y_{1}+\ldots+x y_{q}$ are created. So $S$ does not contain any determined
leaf of $T+x y_{1}+\ldots+x y_{q}$. Moreover, the vertices of $T+x y_{1}+\ldots+x y_{q}$ are the vertices of $T$ plus $y_{1}, \ldots, y_{q}$. By our inductive hypothesis, any vertex in $T+x y_{1}+\ldots+x y_{q}$ which has out-neighbours not in $T+x y+\ldots+x y_{q}$ is either in $S \cup\left\{y_{1}, \ldots, y_{q}\right\}$ or a determined leaf of $T$ which is also a determined leaf of $T+x y+\ldots+x y_{q}$. Since $T+x y+\ldots+x y_{q}$ is the $T$ and $S \cup\left\{y_{1}, \ldots, y_{q}\right\}$ is the $S$ before the next iteration, the Claim is true before the next iteration of the while loop in step 2 .

In step $2.2, x$ has no out-neighbours in $G-T$, according to the condition of the if statement in step 2.1. By out inductive hypothesis, $S-x$ does not contain any determined leaf of $T$, and any vertex of $T$ which have out-neighbours in $G-T$ is either in $S-x$ or a determined leaf of $T$. Since the $T$ before the next iteration is the same $T$ before this iteration, and the $S$ before the next iteration is $S-x$, the Claim is true before the next iteration of the while loop in step 2. This completes the proof of the Claim.

By the Claim, $S$ in step 2.1 does not contain any determined leaf of $T$. Then $x \in S$ in step 2.1 is not a determined leaf of $T$. By Lemma C. 10 and that $k$ is a $k$-outtree, the instance $\left(T ; \pi\left(l_{1}\right), \cdots, \pi\left(l_{q}\right) ; k\right)$ before in step 2.1 is reducible to the instance $\left(T+x y_{1}+\cdots+x y_{q} ; \pi_{y_{1}, \cdots, y_{q}}\left(l_{1}\right), \cdots, \pi_{y_{1}, \cdots, y_{q}}\left(l_{q}\right) ; k\right)$. Therefore, by the transitivity of the reduction, the instance $\left(T_{1} ; \pi_{1}\left(l_{1}\right), \ldots, \pi_{1}\left(l_{q}\right) ; k\right)$ before step 2 is reducible to $\left(T_{2} ; \pi_{2}\left(l_{1}\right), \ldots, \pi_{2}\left(l_{q}\right) ; k\right)$, the instance before step 3 .

By the Claim, only determined leaves of $T_{3}$ can have out-neighbours in $G-T_{3}$. Since any determined leaf $l_{i}$ of $T_{i}$ can only reach vertices in $\pi\left(l_{i}\right)$, step 3 returns 'NO' correctly, if there is a vertex $x$ not in $T$ nor in any out-chain.

After step 3, every vertex not in $T$ is in some out-chain. We prove that step 4 finds an extended out-branching $T_{s}$ for the instance $\left(T_{4} ; \pi_{4}\left(l_{1}\right), \ldots, \pi_{4}\left(l_{q}\right) ; k\right)$ where $T_{s}$ is the $T$ after step 4 by the following loop invariant: before the $t$ th execution of step 4.1, $T$ is an out-tree with the same root as $T_{4}$, the vertices of any out-chain in
the ordered set $\left\{\pi_{4}\left(l_{p-t+2}\right), \ldots, \pi_{4}\left(l_{p-1}\right), \pi_{4}\left(l_{p}\right)\right\}$ are in $T$, and if a vertex $w$ of $\pi_{4}\left(l_{j}\right)$ where $j<p-t+2$ is in $T$, then the vertices in the path from $w$ to the last vertex of $\pi_{4}\left(l_{t}\right)$ are in $T$. Note that $i=p-t+1$ for the $t$ th execution of step 4.1.

The initialization case is when $i=1$, since $T_{4}$ is an out-tree and the ordered set $\left\{\pi_{4}\left(l_{p-t+2}\right), \ldots, \pi_{4}\left(l_{p}\right)\right\}$ is empty $(p-t+2=p+1>p)$. That is, the loop invariant is true for $t=1$. Suppose the loop invariant is true before the $t$ th execution of step 4.1, i.e., before the $t$ th execution of step 4.1, $T$ is an out-tree with the same root as $T_{4}$, the vertices of any out-chain in $\pi_{4}\left(l_{p-t+2}\right), \ldots, \pi_{4}\left(l_{p-1}\right), \pi_{4}\left(l_{p}\right)$ are in $T$, and if a vertex $w$ of $\pi_{4}\left(l_{j}\right)$ where $j<p-t+2$ is in $T$, then the vertices in the path from $w$ to the last vertex of $\pi_{4}\left(l_{t}\right)$ are in $T$.

We show it is still true before the $(t+1)$ th iteration. When the step 4.1 is skipped because all vertices of $\pi\left(l_{p-t+1}\right)$ are already in $T$, then the loop invariant still holds, since nothing changes.

During the $t$ th execution of step 4.1, all remaining vertices of $\pi_{4}\left(l_{p-t+1}\right)$ which are not in $T$ yet are included in $T$ during step 4.1. So all vertices of $\pi_{4}\left(l_{p-t+1}\right)$ are in $T$ after the $t$ th execution of step 4.1. If a vertex $w \notin T$ of $\pi_{4}\left(l_{j}\right)$ where $j<p-t+1$ is visited during the $t$ th execution of step 4.1, then all vertices from $w$ to the last vertex of $\pi_{4}\left(l_{j}\right)$ are in $T$ after the $t$ th execution of step 4.1, since the set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{q}\right)\right\}$ is consistent.

Moreover, the $T$ before the $(t+1)$ th execution of step 4.1, is still an out-tree with the same root as the $T$ before the $t$ th execution, since the $i$ th execution of step 4.1 just adding a path to a leaf of $T$. By our inductive hypothesis, the $T$ before the $t$ th execution of step 4.1 has the same root as $T_{4}$. Thus, the $T$ before the $(t+1)$ th execution of step 4.1, has the same root as $T_{4}$. In consequence, the loop invariant is true before the $(t+1)$ th execution of step 4.1.

When the for loop terminates, $T_{s}=T$ is an out-tree with the same root as $T_{4}$
(by the transitivity of the reduction), and all vertices in $\pi_{4}\left(l_{1}\right), \ldots, \pi_{4}\left(l_{q}\right)$ are in $T$. Thus, $T_{s}$ is an out-branching with the same root as $T_{4}$. Since $T_{4}$ is a subgraph of $T_{s}$ (no edges of $T$ are removed during step 4 ), $T_{s}$ is a $k$-out-branching by Lemma B.2. By step 4.1, $l_{i}$ can only reach vertices in $\pi_{4}\left(l_{i}\right)$ in $T_{s}$ for all $1 \leq i \leq q$. So $T_{s}$ is an extended out-branching for $\left(T_{4} ; \pi_{4}\left(l_{1}\right), \ldots, \pi_{4}\left(l_{q}\right) ; k\right)$.

Step 5 returns $T_{s}$ constructed in step 4, which is an extended out-branching for $\left(T_{4} ; \pi_{4}\left(l_{1}\right), \ldots, \pi_{4}\left(l_{q}\right) ; k\right)$.

Summarizing the above discussion, we conclude that the algorithm SpanTree either constructs an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ if such one exists, or correctly reports 'No' otherwise.

Now we analyze the complexity of the algorithm.
In both step 2.1 and 2.2, a vertex is removed from $S$. Only in step 2.1, vertices are added into $S$. However, only vertices not in $T$ may be added into $S$ in step 2.1, and these vertices are in $T$ once they are added into $S$. So any vertex can be added into $S$ at most one time. In consequence, step 2 can not be executed with more than $n$ iterations. It is obvious that steps 2.1 and 2.2 take polynomial time. Thus, step 2 can be done in polynomial time. Moreover, we can see that steps 3,4 and 5 can be done in polynomial time. Therefor, the algorithm SpanTree runs in polynomial time.
D. The Main Algorithm and Complexity Analysis

Before we solve the max-LEAF OUT-BRANCHING problem, we give an algorithm to solve the problem EXTENDING MAX-LEAF with inputs $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ on a digraph $G$, which is used as a subroutine for the algorithm solving the MAX-LEAF out-branching problem. Consider the algorithm ExtendTree given in Figure 11.

ExtendTree $\left(G, T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$
InPUT: an out-tree $T$ in a digraph $G$, a consistent set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}$ of out-chains for $T$, and a parameter $k$
OUTPUT: an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ if such one exists, or "No" otherwise.

1. if $T$ is an out-branching with fewer than $k$ leaves then return 'No';
2. if $T$ has at least $k$ leaves then return $\operatorname{SpanTree}\left(G, T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$;
3. if an internal vertex $x \in T$ has out-neighbors $y_{1}, \ldots, y_{q}$ in $G-T$ then return ExtendTree $\left(G, T+x y_{1}+\ldots+x y_{q} ; \pi_{y_{1}, \ldots, y_{q}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{q}}\left(l_{p}\right) ; k\right)$;
4. if all leaves of $T$ are determined leaves then return 'No';
5. $\quad x=$ any undetermined leaf $l$ of $T$;
$\pi(l)=x ; \quad x_{s}=x ;$
while $x_{s}$ has a unique out-neighbor $y$ in $G-(T \cup \pi(l))$ do add $y$ to $\pi(l) ; \quad x_{s}=y$;
if $x_{s}$ has no out-neighbors in $G-(T \cup \pi(l))$ then return ExtendTree $\left(G, T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right), \pi(l) ; k\right)$;
7.1 else $T^{\prime}=\operatorname{ExtendTree}\left(G, T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right), \pi(l) ; k\right)$;
if $T^{\prime}$ is not ' No ' then return $T^{\prime}$;
7.2 let $y_{1}, \ldots, y_{q}$ be the out-neighbors of $x_{s}$ in $G-(T \cup \pi(l))$, return ExtendTree $\left(G, T+\pi(l)+x_{s} y_{1}+\ldots+x_{s} y_{q} ; \pi_{y_{1}, \ldots, y_{q}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{q}}\left(l_{p}\right) ; k\right)$

Fig. 11. Algorithm for the EXtending max-leaf problem

Theorem D. 1 For a given out-tree $T$ in a digraph $G$ of $n$ vertices, with a consistent set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right)\right\}$ of out-chains for $T$, the algorithm ExtendTree runs in time $O^{*}\left(4^{k}\right)$ and either constructs an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ if such one exists, or reports ' $N o$ ' otherwise.

Proof. We first verify the correctness of the algorithm.
If the input $T$ is already an out-branching with less than $k$ leaves, then there is no extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. So step 1 reports 'No' correctly. If $T$ is a $k$-out-tree, then by Theorem C.11, the algorithm SpanTree either constructs
an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ if such one exists, or reports 'No' otherwise. Therefore, step 2 is correct. In step 3, $x$ is an internal vertex of $T$. By Lemma C.8, the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to the instance $\left.\left(T+x y_{1}+\ldots+x y_{q}\right) ; \pi_{y_{1}, \ldots, y_{q}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{q}}\left(l_{p}\right) ; k\right)$. Thus, step 3 is correct.

After step 3, only leaves of $T$ may have out-neighbors in $G-T$ and $T$ has fewer than $k$ leaves. Suppose that $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. Then any determined leaf $l_{i}$ of $T$ can only reach vertices of $\pi\left(l_{i}\right)$ in $T_{s}$. By the definition of out-chains, $l_{i}$ can only reach vertices in $\pi\left(l_{i}\right)$ one by one in $T_{s}$. So $l_{i}$ can only reach one leaf in $T_{s}$. If all leaves of $T$ are determined leaves, then $T_{s}$ has fewer than $k$ leaves since $T$ has fewer than $k$ leaves. This contradicts that $T_{s}$ is an extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. Therefore, step 4 returns 'No' correctly.

After step 4, $T$ has fewer than $k$ leaves and some leaves of $T$ are not determined leaves. Step 5 constructs an out-chain $\pi(l)$ for an undetermined leaf $l$ such that the end vertex $x_{s}$ of $\pi(l)$ either has at least two out-neighbors or has no out-neighbors in $G-(T \cup \pi(l))$. Note that $\pi(l)$ can be $\{l\}$. Moreover, for any $1 \leq i \leq p$, if $\pi\left(l_{i}\right)$ and $\pi(l)$ have a common vertex $z$, then the path from $z$ to the last vertex $w$ of $\pi\left(l_{i}\right)$ must be in $\pi(l)$.

Otherwise, let $t$ be the first vertex in the path from $z$ to $w$ in $\pi\left(l_{i}\right)$ which is not in $\pi(l)$, and let $p$ be the in-neighbour of $t$ in $\pi(l)$. Since $z$ is in $\pi(l), t \neq z$ and since $t \neq l, p$ is in $\pi(l)$. Since $p t$ is in $\pi\left(l_{i}\right), p$ has only an out-neighbour $t$ not in $T$, the out-tree for which we construct $\pi\left(l_{i}\right)$. Algorithm ExtendTree does remove edges from $T$, so $T$ is a subgraph of the out-tree $T^{\prime}$, for which $\pi(l)$ is constructed. Thus, $p$ still has at most one out-neighbour $t$ not in $T^{\prime}$. If $p$ has $t$ as the only out-neighbour not in $T^{\prime}$, then $t$ should be added into $\pi(l)$ in step 5 , contradicting the assumption that $t$ is not in $\pi(l)$. If $p$ has no out-neighbours not in $T^{\prime}$, then $t$ is in $T^{\prime}$. Only steps 3 and 7.2 add to $T^{\prime}$ vertices not in $T$. Both steps would change $\pi\left(l_{i}\right)$ to $\pi_{t}\left(l_{i}\right)$, which should
not contain $t$. This contradicts that $\pi\left(l_{i}\right)$ contains $t$. In consequence, the path from $z$ to the last vertex $w$ of $\pi\left(l_{i}\right)$ must be in $\pi(l)$. Therefore, the set $\left\{\pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right), \pi(l)\right\}$ of out-chains is consistent.

If the last vertex $x_{s}$ of $\pi(l)$ has no out-neighbor in $G-(T \cup \pi(l))$, then all vertices reachable from $l$ are in $\pi(l)$ for all possible extended out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$. This means that $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right), \pi(l) ; k\right)$ is the only instance to solve. So step 6 is correct.

Now suppose $x_{s}$ has more than one out-neighbors in $G-(T \cup \pi(l))$. If there is an extended out-branching $T_{s}$ for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right), \pi(l) ; k\right)$, then $T_{s}$ is also a $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$-extended out-branching. So step 7.1 is correct. If step 7.1 can not find an out-branching for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right), \pi(l) ; k\right)$, then for any possible extended out-branching $T_{s}$ for $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$, the vertices reachable in $T_{s}$ from $l$ must contain some vertex $x_{s}$ not in $\pi(l)$. This can only happen when $x_{s}$ is reachable from $l$ and $x_{s}$ is not a leaf in $T_{s}$, i.e, $x_{s}$ is an internal vertex of $T_{s}$. Then all vertices in $\pi(l)$ should be reachable from $l$ in $T_{s}$. By Lemma C.8, $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ is reducible to $\left.\left(T+\pi(l)+x_{x} y_{1}+\ldots+x_{s} y_{q}\right) ; \pi_{y_{1}, \ldots, y_{q}}\left(l_{1}\right), \ldots, \pi_{y_{1}, \ldots, y_{q}}\left(l_{p}\right) ; k\right)$. Thus step 7.2 is correct.

Summarizing the above discussion, we conclude that the algorithm ExtendTree correctly solves the instance $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ of the EXTENDING MAX-LEAF problem.

Now we consider the complexity of the algorithm.
Let $u$ the number of leaves of $T, p$ be the number of determined leaves of $T$, then $a=k-u$ and $b=k-p$. Let $f(a, b)=a+b=2 k-u-p$ be the complexity of the algorithm ExtendTree. We use the search-tree method to analyze the time complexity. Step 1 takes time $O(n)$. Step 2 takes time $O\left(n^{2}\right)$ by Theorem C.11. Steps 3, 4, and 5 can be done in time polynomial of $n$.

At step 6, the leaf $l$ becomes a determined leaf from a undetermined leaf. So $u^{\prime}=u$ and $p^{\prime}=p+1$. Then the complexity changes from $f(a, b)=2 k-u-p$ to $f^{\prime}\left(a^{\prime}, b^{\prime}\right)=2 k-u-p-1=f(a, b)-1$. At step 7.1, the complexity is the same as step 6: $f\left(a_{1}, b_{1}\right)=2 k-u-p-1=f(a, b)-1$. Since $x_{s}$ has at least two out-neighbors in step $7.2, u^{\prime}=u+1$ and $p^{\prime}=p$. So the complexity of step 7.2 is $f\left(a_{2}, b_{2}\right)=2 k-u-p-1=f(a, b)-1$. Since steps 7.1 and 7.2 are two branches, we have the recurrence equation: $f(a, b)=2 f(a, b)-1$.

Initially, $a+b=2 k-u-p \leq 2 k$. When either $a=k-u=0$ or $b=k-p=0, T$ is a $k$-out-tree. By Theorem C.11, Extending $\left(T ; \pi\left(l_{1}\right), \ldots, \pi\left(l_{p}\right) ; k\right)$ can be solved in time of $O\left(n^{2}\right)$, i.e., no branches are needed when $a=0$ or $b=0$. Then when $a+b=0$, no branches are needed for ExtendTree, i.e., $f(0,0)=1$. By $f(a, b)=2 f(a, b)-1$, we have $f(a, b) \leq 2^{a+b}=O\left(4^{k}\right)$, which is an upper bound of the total number of leaves in the search-tree for the algorithm ExtendTree. Since steps 1, 2, and 4 are executed only once, steps 3 and 5 are executed at most $O(n)$ times from the root to a leaf in the search-tree, the depth of the search-tree is $O(n)$, and the running time is polynomial of $n$ for each step, the total running time of SpanTree is $O^{*}\left(4^{k}\right)$.

Now we are ready to present our algorithm for the MAX-LEAF OUT-BRANCHING problem.

Theorem D. 2 There is an algorithm to solve the MAX-LEAF OUT-BRANCHING problem in time $O\left(n^{4} 4^{k}\right)$.

Proof. The algorithm for solving the max-leaf out-branching problem is very simple: for each vertex $x \in G$, we call the algorithm ExtendTree $\left(G, T+x y_{1}+\right.$ $\left.\ldots+x y_{q}\right) ; \emptyset ; k$ ) where (1) $T$ is the single vertex $x ;(2) y_{1}, \ldots, y_{q}$ are all out-neighbors of $x$ in $G$; and (3) $\emptyset$ means that all leaves of $T+x y_{1}+\ldots+x y_{q}$ are not determined. If an extended out-branching $T_{s}$ for $\left(T+x y_{1}+\ldots+x y_{q} ; \emptyset ; k\right)$ is found for some $x$,
then $T_{s}$ is an $k$-out-branching of $G$; otherwise reports 'No'.
To prove that the algorithm is correct, we prove the following claim.
Claim. If there is an $k$-out-branching $T_{s}$ of $G$ whose root is $r$ with out-neighbour $y_{1}, \ldots, y_{q}$, then ExtendTree $\left(G, T+x y_{1}+\ldots+x y_{q} ; \emptyset ; k\right)$ find an extended out-branching for $\left(T+x y_{1}+\ldots+x y ; \emptyset ; k\right)$.

First $r$ must have some out-neighbor $z$ in $T_{s}$, and $z$ must be in $\left\{y_{1}, \ldots, y_{q}\right\}$. Suppose that $z=y_{1}$. Then $T_{s}$ is an extended out-branching for $\left(T+r y_{1} ; \emptyset ; k\right)$ since (1) $T_{s}$ is a $k$-out-branching; (2) $T_{s}$ and $T+r y_{1}$ have the same root $r$; (3) $T+r y_{1}$ is a subgraph of $T_{s}$; and (4) all vertices reachable in $T_{s}$ from a determined leaf $l_{i}$ of $T+r y_{1}$ are in $\pi\left(l_{i}\right)$ since there is no determined leaf in $T+r y_{1}$. Therefore, $T_{s}$ is an extended out-branching for $\left(T+r y_{1} ; \emptyset ; k\right)$, i.e., we only need to solve the instance $\left(T+r y_{1} ; \emptyset ; k\right)$.

Now $r$ is an internal vertex of $T+r y_{1}$, so the instance $\left(T+r y_{1} ; \emptyset ; k\right)$ is reducible to the instance $\left(T+r y_{1}+\ldots+r y_{q}, \emptyset ; k\right)$ by Lemma C.8. Since there is an extended out-branching for $\left(T+r y_{1} ; \emptyset ; k\right)$, there is an extended out-branching for $\left(T+r y_{1}+\right.$ $\left.\ldots+r y_{q} ; \emptyset ; k\right)$. Thus ExtendTree $\left(G, T+r y_{1}+\ldots+r y_{q} ; \emptyset ; k\right)$ will find an extended out-branching for $\left(T+r y_{1}+\ldots+r y ; \emptyset ; k\right)$. This proves the Claim.

If there is a $k$-out-branching $T_{s}$ rooted at $r$ for $G$, the algorithm above for the max-Leaf out-branching problem must call ExtendTree $\left(G, T+r y_{1}+\ldots+\right.$ $\left.\left.r y_{q}\right) ; \emptyset ; k\right)$, where $y_{1}, \ldots, y_{q}$ are the out-neighbors of the vertex $r$, since ExtendTree is executed on every vertex in $G$. By our claim, ExtendTree $\left(G, T+r y_{1}+\ldots+r y_{q} ; \emptyset ; k\right)$ will find an extended out-branching for $\left(T+r y_{1}+\ldots+r y ; \emptyset ; k\right)$, which is a $k$-outbranching by definition.

If no extended out-branching for $\left.\left(T+x y_{1}+\ldots+x y_{q}\right) ; \emptyset ; k\right)$ is found for any $x \in G$, then there is no $k$-out-branching of $G$ by our Claim. In summary, the algorithm for the max-leaf out-branching problem is correct: it either constructs a $k$-out-
branching if such one exists, or reports 'No' otherwise.
The running time of this algorithm is $O(n)$ times the time of ExtendTree, thus is $O^{*}\left(4^{k}\right)$.

## E. Final Remarks

In this chapter, we presented an $O^{*}\left(4^{k}\right)$ time algorithm for the MAX-LEAF OUTBRANCHING problem on digraphs, which significantly improves the previous best algorithm of running time $O^{*}\left(2^{O(k \log k)}\right)$ for the problem. The algorithm can be applied directly to solve the simpler MAX-LEAF SPANNING-TREE problem on undirected graphs, which also gives an improvement over the previous best algorithm of running time $O^{*}\left(6.75^{k}\right)$ for the problem. The improvements are based on a deeper study of the combinatorial structures of digraphs that reveals further relationship between out-trees and out-branchings in digraphs, and on applications of a new algorithmic technique for the design and analysis of parameterized algorithms. In particular, the new algorithmic technique identifies indirected branching measures that bound the number of computational branches in a branch-and-search process that do not reduce the parameter value effectively.

Recently, Kneis et al. independently developed an algorithm of the same running time $O^{*}\left(4^{k}\right)$ [69]. Their algorithm and our algorithm use the same idea of extending an out-tree to an out-branching.

## CHAPTER VII

## SATISFIABILITY

In this chapter, we study the complexity of the well-known SATISFIABILITY problem with respect to the total length $L$ of formulas. We present an exact algorithm of running time $O^{*}\left(1.0652^{L}\right)$ for the SATISFIABILITY problem, which improves the previous best algorithm of running time $O^{*}\left(1.0663^{L}\right)$ [80]. Our algorithm considers both the number of variables and their associated weights as the measures for the SATISFIABILITY problem. By considering such measures, we design reduction rules to design a simple algorithm for the SATISFIABILITY problem. Our result demonstrates that our measure-driven approach is also powerful for designing exact algorithms (without parameters) for NP-hard problems.

## A. Introduction

The satisfiability problem (briefly, SAT: given a CNF Boolean formula, decide if the formula has a satisfying assignment) is perhaps the most famous and most extensively studied NP-complete problem. The problem requires a precise answer Yes/No, and approximation algorithms do not seem to help much. Given the NPcompleteness of the problem [53], it has become natural to develop exponential time algorithms that solve the problem as fast as possible.

There are three popular parameters that have been used in measuring exponential time algorithms for the SAT problem: the number $n$ of variables in the input formula, the number $m$ of clauses in the input formula, and the total length $L$ of the input formula, which is the sum of the clause lengths in the input formula. Note that the parameter $L$ is probably the most precise parameter in terms of standard complexity theory, and both parameters $n$ and $m$ could be sublinear in instance length.

Algorithms for SAT in terms of each of these parameters have been extensively studied. See [55] for a comprehensive review and see [99] for more recent progress on the research in these directions.

In the current chapter, we are focused on algorithms for SAT in terms of the parameter $L$. The research started 20 years ago since the first published algorithm of time $O\left(1.0927^{L}\right)$ [54]. The upper bound was subsequently improved by an impressive list of publications. We summarize the major progress in the following table.

Table III. History of exact algorithms for the satisfiability problem

| Reference | Bound | Year Published |
| :---: | :---: | :---: |
| Van Gelder [54] | $1.0927^{L}$ | 1988 |
| Kullmann et al. $[73]$ | $1.0801^{L}$ | 1997 |
| Hirsh $[62]$ | $1.0758^{L}$ | 1998 |
| Hirsh $[61]$ | $1.074^{L}$ | 2000 |
| Wahlstom $[80]$ | $1.0663^{L}$ | 2005 |

The branch-and-search method has been widely used in the development of SAT algorithms. Given a Boolean formula $\mathcal{F}$, let $\mathcal{F}[x]$ and $\mathcal{F}[\bar{x}]$ be the resulting formula after assigning TRUE and false, respectively, to the variable $x$ in the formula $\mathcal{F}$. The branch-and-search method is based on the fact that $\mathcal{F}$ is satisfiable if and only if at least one of $\mathcal{F}[x]$ and $\mathcal{F}[\bar{x}]$ is satisfiable. Most SAT algorithms are based on this method.

Unfortunately, analysis directly based on the parameter $L$ usually does not give a good upper bound in terms of $L$ for a branch-and-search SAT algorithm. Combinations of the parameter $L$ and other parameters, such as the number $n$ of variables in the input formula, have been used as "measures" in the analysis of SAT algorithms. For example, the measure $L-2 n[81]$ and a more general measure that is a function $f(L, n)$ of the parameters $L$ and $n$ [80] have been used in the analysis of SAT
algorithms whose complexity is measured in terms of the parameter $L$.
In the current chapter, we introduce a new measure, the l-value of a Boolean formula $\mathcal{F}$. Roughly speaking, the measure $l$-value $l(\mathcal{F})$ is defined based on weighted variable frequencies in the input formula $\mathcal{F}$. We develop a branch-and-search algorithm that tries to maximize the decreasing rates in terms of the $l$-value during its branch-and-search process. In particular, by properly choosing the variable frequency weights so that the formula $l$-value is upper bounded by $L / 2$, adopting new reduction rules, and applying the analysis technique of Measure and Conquer recently developed by Fomin et al. [50], we develop a new branch-and-search algorithm for the SAT problem whose running time is bounded by $O\left(1.1346^{l(\mathcal{F}))}\right.$ on an input formula $\mathcal{F}$. Finally, by combining this algorithm with the algorithm in [81] to deal with formulas of lower variable frequencies and converting the measure $l(\mathcal{F})$ into the parameter $L$, we achieve a SAT algorithm of running time $O\left(1.0652^{L}\right)$, improving the previously best SAT algorithm of running time $O\left(1.0663^{L}\right)$ [80].

We remark that although the analysis of our algorithm is lengthy, our algorithm itself is very simple and can be easily implemented. Note that the lengthy analysis needs to be done only once to ensure the correctness of the algorithm, while the simplicity of the algorithm gives its great advantage when it is applied (many times) to determine the satisfiability of CNF Boolean formulas.

## B. Preliminaries

We introduce the notations and terminology that will be used in our discussion.
A (Boolean) variable $x$ can be assigned value either 1 (TRUE) or 0 (FALSE). The variable $x$ has two corresponding literals $x$ and $\bar{x}$. The literal $x$ is satisfied if $x=1$ and the literal $\bar{x}$ is satisfied if $x=0$. Note that exactly one of the literals $x$ and $\bar{x}$
is satisfied for any value assignment to the variable $x$. A clause $C$ is a disjunction of a set of literals, which can be regarded as a set of literals. Therefore, we may write $C_{1}=z C_{2}$ to indicate that the clause $C_{1}$ consists of the literal $z$ plus all literals in the clause $C_{2}$, and use $C_{1} C_{2}$ to denote the clause that consists of all literals that are in either $C_{1}$ or $C_{2}$, or both. The length of a clause $C$, denoted by $|C|$, is the number of literals in $C$. A clause $C$ is satisfied if any literal in $C$ is satisfied. A (CNF Boolean) formula $\mathcal{F}$ is a conjunction of clauses $C_{1}, \ldots, C_{m}$, which can be regarded as a collection of the clauses. The formula $\mathcal{F}$ is satisfied if all clauses in $\mathcal{F}$ are satisfied. The length $L$ of the formula $\mathcal{F}$ is defined as $L=\left|C_{1}\right|+\cdots+\left|C_{m}\right|$.

A literal $z$ is an $i$-literal if $z$ is contained in exactly $i$ clauses, and is an $i^{+}$-literal if $z$ is contained in at least $i$ clauses. An $(i, j)$-literal $z$ is a literal such that exactly $i$ clauses contain $z$ and exactly $j$ clauses contain $\bar{z}$. Note that an $(i, j)$-literal $z$ implies that the literal $\bar{z}$ is a $(j, i)$-literal. We say that a clause $C$ contains a variable $x$ if $C$ contains either the literal $x$ or the literal $\bar{x}$. A variable $x$ is an $i$-variable if exact $i$ of the clauses contain the variable $x$ (in this case, we also say that the degree of $x$ is $i$ ). A clause $C$ is an $i$-clause if $|C|=i$, and is an $i^{+}$-clause if $|C| \geq i$.

Let $x C_{1}, \ldots, x C_{s}$ be all the clauses in the input formula $\mathcal{F}$ that contain the literal $x$, and let $\bar{x} D_{1}, \ldots, \bar{x} D_{t}$ be all the clauses in $\mathcal{F}$ that contain the literal $\bar{x}$. A resolvent on a variable $x$ in $\mathcal{F}$ is a clause of the form $C_{i} D_{j}$ for some $i$ and $j, 1 \leq i \leq s$ and $1 \leq j \leq t$. The resolution on the variable $x$ in the formula $\mathcal{F}$, written as $D P_{x}(\mathcal{F})$, is a formula that is obtained by first removing all clauses containing the variable $x$ from $\mathcal{F}$ and then adding all possible resolvents on the variable $x$ into $\mathcal{F}$.

A branching vector is a tuple of positive real numbers. A branching vector $t=$ $\left(t_{1}, \ldots, t_{r}\right)$ corresponds to a polynomial $1-\sum_{i=1}^{r} x^{-t_{i}}$, which has a unique positive root $\tau(t)$ [20]. We say that a branching vector $t^{\prime}$ is inferior to a branching vector $t^{\prime \prime}$ if $\tau\left(t^{\prime}\right) \geq \tau\left(t^{\prime \prime}\right)$. In particular, if either $t_{1}^{\prime} \leq t_{1}^{\prime \prime}$ and $t_{2}^{\prime} \leq t_{2}^{\prime \prime}$, or $t_{1}^{\prime} \leq t_{2}^{\prime \prime}$ and $t_{2}^{\prime} \leq t_{1}^{\prime \prime}$, then
it can be proved [20] that the branching vector $t^{\prime}=\left(t_{1}^{\prime}, t_{2}^{\prime}\right)$ is inferior to the branching vector $t^{\prime \prime}=\left(t_{1}^{\prime \prime}, t_{2}^{\prime \prime}\right)$.

The execution of a SAT algorithm based on the branch-and-search method can be represented as a search tree $\mathcal{T}$ whose root is labeled by the input formula $\mathcal{F}$. Recursively, if at a node $w_{0}$ labeled by a formula $\mathcal{F}_{0}$ in the search tree $\mathcal{T}$, the algorithm breaks $\mathcal{F}_{0}$, in polynomial time, into $r$ smaller formulas $\mathcal{F}_{1}, \ldots, \mathcal{F}_{r}$, and recursively works on these smaller formulas, then the node $w_{0}$ in $\mathcal{T}$ has $r$ children, labeled by $\mathcal{F}_{1}$, $\ldots, \mathcal{F}_{r}$, respectively. Suppose that we use a measure $\mu(\mathcal{F})$ for a formula $\mathcal{F}$, then the branching vector for this branching, with respect to the measure $\mu$, is $t=\left(t_{1}, \ldots, t_{r}\right)$, where $t_{i}=\mu\left(\mathcal{F}_{0}\right)-\mu\left(\mathcal{F}_{i}\right)$ for all $i$. Finally, suppose that $t^{\prime}$ is a branching vector that is inferior to all branching vectors for any branching in the search tree $\mathcal{T}$, then the complexity of the SAT algorithm is bounded by $O\left(\tau\left(t^{\prime}\right)^{\mu(\mathcal{F})}\right)$ times a polynomial of $L$ [20].

Formally, for a given input formula $\mathcal{F}$, we define the l-value for $\mathcal{F}$ to be $l(\mathcal{F})=$ $\sum_{i \geq 1} w_{i} n_{i}$, where for all $i \geq 1, n_{i}$ is the number of $i$-variables in $\mathcal{F}$, and the frequency weight $w_{i}$ for $i$-variables are set by the following values:

$$
\begin{array}{lll}
w_{0}=0 ; & w_{1}=0.32 ; & w_{2}=0.45  \tag{7.1}\\
w_{3}=0.997, & w_{4}=1.897, & w_{i}=i / 2, \quad \text { for } i \geq 5
\end{array}
$$

Define $\delta_{i}=w_{i}-w_{i-1}$, for $i \geq 1$. Then we can easily verify that

$$
\begin{align*}
& \delta_{i} \geq 0.5, \quad \text { for all } i \geq 3 \\
& \delta_{\min }=\min \left\{\delta_{i} \mid i \geq 1\right\}=\delta_{2}=0.13  \tag{7.2}\\
& \delta_{\max }=\max \left\{\delta_{i} \mid i \geq 1\right\}=\delta_{4}=0.9
\end{align*}
$$

Note that the length $L$ of the formula $\mathcal{F}$ is equal to $\sum_{i \geq 1} i \cdot n_{i}$, and that $i / 5 \leq w_{i} \leq i / 2$
for all $i$. Therefore, we have $L / 5 \leq l(\mathcal{F}) \leq L / 2$.
Given two formulas $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$, by definition we have $l\left(\mathcal{F}_{1}\right)=\sum_{x \in \mathcal{F}_{1}} w(x)$ and $l\left(\mathcal{F}_{2}\right)=\sum_{x \in \mathcal{F}_{2}} w^{\prime}(x)$, where $w(x)$ is the frequency weight of $x$ in $\mathcal{F}_{1}$ and $w^{\prime}(x)$ is the frequency weight of $x$ in $\mathcal{F}_{2}$. The $l$-value reduction from $\mathcal{F}_{1}$ to $\mathcal{F}_{2}$ is $l\left(\mathcal{F}_{1}\right)-$ $l\left(\mathcal{F}_{2}\right)$. The contribution of $x$ to the $l$-value reduction from $\mathcal{F}_{1}$ to $\mathcal{F}_{2}$ is $w(x)-w^{\prime}(x)$. The contribution of a variable set $S$ to the $l$-value reduction from $\mathcal{F}_{1}$ to $\mathcal{F}_{2}$ is the summation of contributions of all variables in $S$.

## C. Reduction Rules

We say that two formulas $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$ are equivalent if $\mathcal{F}_{1}$ is satisfiable if and only if $\mathcal{F}_{2}$ is satisfiable. A literal $z$ in a formula $\mathcal{F}$ is monotone if the literal $\bar{z}$ does not appear in $\mathcal{F}$.

We present in this section a set of reduction rules that reduce a given formula $\mathcal{F}$ to an equivalent formula $\mathcal{F}^{\prime}$ without increasing the $l$-value. Consider the algorithm given in Figure 12.

Lemma C. 1 The algorithm Reduction $\left(\mathcal{F}_{1}\right)$ produces a formula equivalent to the formula $\mathcal{F}_{1}$.

Proof. It suffices to prove that in each of the listed cases, the algorithm Reduction on the formula $\mathcal{F}_{1}$ produces an equivalent formula $\mathcal{F}_{2}$. This can be easily verified for Cases 1, 2, 3, and 5.

The claim holds true for Cases 4 and 9 from the resolution principle [31].
In Case 6 , the clause $z_{1} \overline{z_{2}} C$ in $\mathcal{F}_{1}$ is replaced with the clause $z_{1} C$ in $\mathcal{F}_{2}$. If an assignment $A_{2}$ satisfies $\mathcal{F}_{2}$, then obviously $A_{2}$ also satisfies $\mathcal{F}_{1}$. On the other hand, if an assignment $A_{1}$ satisfies $\mathcal{F}_{1}$ but does not satisfy the clause $z_{1} C$ in $\mathcal{F}_{2}$, then because
of the clause $z_{1} z_{2}$ in $\mathcal{F}_{1}$, we must have $z_{1}=0$ and $z_{2}=1$. Since $A_{1}$ satisfies the clause $z_{1} \overline{z_{2}} C$ in $\mathcal{F}_{1}$, this would derive a contradiction that $A_{1}$ must satisfy $z_{1} C$. Therefore, $A_{1}$ must also satisfy the formula $\mathcal{F}_{2}$.

In Case 7 , the clause $z_{1} z_{2} C_{1}$ in $\mathcal{F}_{1}$ is replaced with the clause $z_{2} C_{1}$ in $\mathcal{F}_{2}$. Again, the satisfiability of $\mathcal{F}_{2}$ trivially implies the satisfiability of $\mathcal{F}_{1}$. For the other direction, let $z_{2} C_{3}$ be the third clause that contains $z_{2}$ (note that $z_{2}$ is a $(2,1)$-literal). If an assignment $A_{1}$ satisfies $\mathcal{F}_{1}$ (in particular satisfies the clause $z_{1} z_{2} C_{1}$ ) but not $\mathcal{F}_{2}$ (i.e., not the clause $z_{2} C_{1}$ ), then we must have $z_{1}=1, z_{2}=0$, and $C_{3}=1$ under $A_{1}$. By replacing the assignment $z_{2}=0$ with $z_{2}=1$ in $A_{1}$, we will obtain an assignment $A_{2}^{\prime}$ that satisfies all $z_{2} C_{1}, z_{1} \overline{z_{2}} C_{2}$, and $z_{2} C_{3}$, thus satisfies the formula $\mathcal{F}_{2}$.

In Case 8 , the formula $\mathcal{F}_{2}$ is obtained from the formula $\mathcal{F}_{1}$ by removing the clause $z_{1} z_{2}$. Thus, the satisfiability of $\mathcal{F}_{1}$ trivially implies the satisfiability of $\mathcal{F}_{2}$. On the other hand, suppose that an assignment $A_{2}$ satisfying $\mathcal{F}_{2}$ does not satisfy $\mathcal{F}_{1}$ (i.e., does not satisfy the clause $z_{1} z_{2}$ ). Then $A_{2}$ must assign $z_{1}=0$ and $z_{2}=0$. We can simply replace $z_{1}=0$ with $z_{1}=1$ in $A_{2}$ and keep the assignment satisfying $\mathcal{F}_{2}$ : this is because $\overline{z_{1} z_{2}} C$ is the only clause in $\mathcal{F}_{2}$ that contains $\overline{z_{1}}$. Moreover, the new assignment now also satisfies $z_{1} z_{2}$, thus satisfies the formula $\mathcal{F}_{1}$.

For Case 10, it is suffice to show that we can always set $z_{1}=\overline{z_{2}}$ in a satisfying assignment for the formula $\mathcal{F}_{1}$. For the subcase where $z_{1}$ is a 1 -literal and $z_{1} z_{2}$ is a 2-clause, if a satisfying assignment $A_{1}$ for $\mathcal{F}_{1}$ assigns $z_{2}=1$ then we can simply let $z_{1}=\overline{z_{2}}=0$ since $z_{1}$ is only contained in the clause $z_{1} z_{2}$. If $A_{1}$ assigns $z_{2}=0$ then because of the 2-clause $z_{1} z_{2}, A_{1}$ must assign $z_{1}=\overline{z_{2}}=1$. For the other subcase, note that the existence of the 2 -clauses $z_{1} z_{2}$ and $\overline{z_{1} z_{2}}$ in the formula $\mathcal{F}_{1}$ trivially requires that every assignment satisfying $\mathcal{F}_{1}$ have $z_{1}=\overline{z_{2}}$.

Finally, in Case 11, the clauses $C D_{1}$ and $C D_{2}$ in $\mathcal{F}_{1}$ are replaced with the clauses $\bar{x} C, x D_{1}$, and $x D_{2}$ in $\mathcal{F}_{2}$. If an assignment $A_{1}$ satisfies $\mathcal{F}_{1}$ (thus satisfies $C D_{1}$ and
$C D_{2}$ ), then in case $C=0$ under $A_{1}$ we assign the new variable $x=0$, and in case $C=1$ under $A_{1}$ we assign the new variable $x=1$. It is not hard to verify that this assignment to the new variable $x$ plus $A_{1}$ will satisfy $\mathcal{F}_{2}$. For the other direction, suppose that an assignment $A_{2}$ satisfies $\mathcal{F}_{2}$. If $A_{2}$ assigns $x=1$ then we have $C=1$ under $A_{2}$ thus the assignment $A_{2}$ also satisfies $C D_{1}$ and $C D_{2}$ thus $\mathcal{F}_{1}$; and if $A_{2}$ assigns $x=0$ then we have $D_{1}=1$ and $D_{2}=1$ under $A_{2}$ and again $A_{2}$ satisfies $\mathcal{F}_{1}$.

Next, we show that the algorithm Reduction always decreases the $l$-value.

Lemma C. 2 Let $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$ be two formulas such that $\mathcal{F}_{2}=\operatorname{Reduction}\left(\mathcal{F}_{1}\right)$, and $\mathcal{F}_{1} \neq \mathcal{F}_{2}$. Then $l\left(\mathcal{F}_{1}\right) \geq l\left(\mathcal{F}_{2}\right)+0.003$.

Proof. Since $\mathcal{F}_{1} \neq \mathcal{F}_{2}$, at least one of the cases in the algorithm Reduction is applicable to the formula $\mathcal{F}_{1}$. Therefore, it suffices to verify that each case in the algorithm Reduction decreases the $l$-value by at least 0.003 .

Cases 1-8 simply remove certain literals, which decrease the degree of certain variables in the formula. Therefore, if any of these cases is applied on the formula $\mathcal{F}_{1}$, then the $l$-value of the formula is decreased by at least $\delta_{\min }=\delta_{2}=0.13$.

Consider Cases 9-11. Note that if we reach these cases then Cases $4-5$ are not applicable, which implies that the formula $\mathcal{F}_{1}$ contains only $3^{+}$-variables.

Case 9. If both $\bar{z}_{1}$ and $\bar{z}_{2}$ are in the same clause, say $C_{1}$, then the resolution $D P_{x}(\mathcal{F})$ after the next application of Case 2 in the algorithm will replace the four clauses $\bar{x} z_{1}, \bar{x} z_{2}, x C_{1}$, and $x C_{2}$ with two 3 -clauses $z_{1} C_{2}$ and $z_{2} C_{2}$, which decreases the $l$-value by $w_{4}$ (because of removing the 4 -variable $x$ ) and on the other hand increases the $l$-value by at most $2 \delta_{\max }$ (because of increasing the degree of the two variables in $C_{2}$ ). Therefore, in this case, the $l$-value is decreased by at least $w_{4}-2 \delta_{\max }=$
$w_{4}-2 \delta_{4}=0.097$. If $\bar{z}_{1}$ and $\bar{z}_{2}$ are not in the same clause of $C_{1}$ and $C_{2}$, say $\bar{z}_{1}$ is in $C_{1}$ and $\bar{z}_{2}$ is in $C_{2}$, then the resolution $D P_{x}(\mathcal{F})$ after the next application of Case 2 in the algorithm will replace the four clauses $\bar{x} z_{1}, \bar{x} z_{2}, x C_{1}$, and $x C_{2}$ with two 3 -clauses $z_{1} C_{2}$ and $z_{2} C_{1}$. In this case, the $l$-value is decreased by exactly $w_{4}=1.897$ because of removing the 4 -variable $x$.

Case 10. Suppose that $z_{1}$ is an $i$-variable and $z_{2}$ is a $j$-variable. Replacing $z_{1}$ by $\bar{z}_{2}$ removes the $i$-variable $z_{1}$ and makes the $j$-variable $z_{2}$ into an $(i+j)$-variable. However, after an application of Case 2 in the algorithm, the clause $z_{1} z_{2}$ in the original formula disappears, thus $z_{2}$ becomes an $(i+j-2)$-variable. Therefore, the total value decreased in the $l$-value is $\left(w_{i}+w_{j}\right)-w_{i+j-2}$. Because of the symmetry, we can assume without loss of generality that $i \leq j$. Note that we always have $i \geq 3$. If $i=3$, then $w_{3}+w_{j}=\delta_{\max }+0.097+w_{j} \geq w_{j+1}+0.097=w_{3+j-2}+0.097$. If $i=4$, then $w_{4}+w_{j}=2 \delta_{\max }+w_{j}+0.097 \geq w_{j+2}+0.097=w_{4+j-2}+0.097$. If $i \geq 5$, then $w_{i}+w_{j}=i / 2+j / 2=(i+j-2) / 2+1=w_{i+j-2}+1$. Therefore, in this case, the $l$-value of the formula is decreased by $\left(w_{i}+w_{j}\right)-w_{i+j-2}$, which is at least 0.097 .

Case 11. Since the clauses $C D_{1}$ and $C D_{2}$ in $\mathcal{F}_{1}$ are replaced with $\bar{x} C, x D_{1}$ and $x D_{2}$, each variable in $C$ has its degree decreased by 1 . Since all variables in $\mathcal{F}_{1}$ are $3^{+}$-variables and $|C| \geq 2$, the degree decrease for the variables in $C$ makes the $l$-value to decrease by at least $2 \cdot \min \left\{\delta_{i} \mid i \geq 3\right\}=1$. On the other hand, the introduction of the new 3 -variable $x$ and the new clauses $\bar{x} C, x D_{1}$ and $x D_{2}$ increases the $l$-value by exactly $w_{3}=0.997$. In consequence, the total $l$-value in this case is decreased by at least $1-0.997=0.003$.

By definition, the $l$-value $l\left(\mathcal{F}_{1}\right)$ of the formula $\mathcal{F}_{1}$ is bounded by $L / 2$, where $L$ is the length $L$ of the formula $\mathcal{F}_{1}$. By the proof of Lemma C.2, each application of a case in the algorithm Reduction takes time polynomial in $L$ and decreases the $l$-value by at least a constant. Therefore, the algorithm Reduction must stop in
polynomial time and produce an equivalent formula $\mathcal{F}_{2}$ for which no cases in the algorithm Reduction are applicable. Such a formula $\mathcal{F}_{2}$ will be called a reduced formula. Reduced formulas have a number of interesting properties, which are given in the following lemmas.

Lemma C. 3 There are no 1-variables or 2-variables in a reduced formula.

Proof. Each 1-variable makes a monotone literal. Each 2-variable either makes a monotone literal or has at most one non-trivial resolvent. Therefore, the Cases 4-5 in the algorithm Reduction ensure that a reduced formula contains no 1-variables and 2-variables.

Lemma C. 4 Let $\mathcal{F}$ be a reduced formula and $x y$ be a clause in $\mathcal{F}$. Then
(1) No other clauses contain $x y$;
(2) No clause contains $x \bar{y}$ or $\bar{x} y$;
(3) At most one $3^{+}$-clause contains $\overline{x y}$. Moreover, if $y$ is a 3 -variable or $x$ is a 1 -literal, then no clause contains $\overline{x y}$.

Proof. (1) Since Case 1 is not applicable, there is no other clauses containing $x y$.
(2) Since Case 6 is not applicable, no clause in $\mathcal{F}$ can contain either $x \bar{y}$ or $\bar{x} y$.
(3) Given a clause $\overline{x y} C$, the clause $C$ cannot be empty since Case 10 is not applicable. If there are two $3^{+}$-clauses containing $\overline{x y}$, Case 11 would be applicable. Therefore, there is at most one $3^{+}$-clause that contains $\overline{x y}$.

If $y$ is a 3 -variable, then $\overline{x y}$ can not be in any clause. Otherwise, the resolution on $y$ would have at most one non-trivial resolvent, and Case 4 would be applicable to $\mathcal{F}$. If $\bar{x}$ is a 1-literal, then $\overline{x y}$ can not be in any clause, since Case 8 is no applicable to $\mathcal{F}$.

Lemma C. 5 Let $\mathcal{F}$ be a reduced formula. Both $x$ and $y$ are variables in $\mathcal{F}$. If only $3^{+}$-clauses contain both variables $x$ and $y$, then for any of $x y, \bar{y}, \bar{y} x$, and $\overline{x y}$, there is at most one clause containing it. Moreover, if $y$ is a 3-variable, then only $\bar{x} y$ can be contained in a $3^{+}$-clause.

Proof. Since $\mathcal{F}$ is a reduced formula, Case 11 is not applicable to $\mathcal{F}$. Then for any of $x y, \bar{y}, \bar{y} x$, and $\overline{x y}$, there is at most one clause containing it.

If $y$ is a 3 -variable, then $\overline{x y}$ can not occur in any clause. Otherwise, the resolution on $y$ would have at most one non-trivial resolvent, and Case 4 would be applicable to $\mathcal{F}$. If $x \bar{y}$ exists, Case 7 would be applicable to $\mathcal{F}$.
D. Main Algorithm

Our main algorithm for the SAT problem is given in Figure 13. The degree $d(\mathcal{F})$ of a formula $\mathcal{F}$ is defined to be the largest degree of a variable in the formula $\mathcal{F}$. Let $\mathcal{F}[x]$ be the resulting formula after removing literal $\bar{x}$ and all clauses containing literal $x$ from $\mathcal{F}$. Similarly, $\mathcal{F}[\bar{x}]$ is the resulting formula after removing literal $x$ and all clauses containing literal $\bar{x}$.

Theorem D. 1 The algorithm $\operatorname{SATSolver}(\mathcal{F})$ solves the $S A T$ problem in time $O\left(1.0652^{L}\right)$, where $L$ is total length of the input formula $\mathcal{F}$.

Proof. It is clear that the algorithm $\operatorname{SATSolver}(\mathcal{F})$ solves the SAT problem. When the degree of $\mathcal{F}$ is 3 , we just apply the algorithm by Wahlström [81] at step 4 in the algorithm $\operatorname{SATSolver}(\mathcal{F})$. The running time of Wahlström's algorithm is $O\left(1.1279^{n}\right)$, where $n$ is the number of variables in $\mathcal{F}$ [81], which is also $O\left(1.1346^{l(\mathcal{F})}\right)$ since $l(\mathcal{F})=w_{3} n$. The proof that the algorithm $\operatorname{SATSolver}(\mathcal{F})$ runs
in time $O\left(1.1346^{l(\mathcal{F})}\right)$ when the degree of $\mathcal{F}$ is greater than 3 is given in the next section. The equality $O\left(1.1346^{l(\mathcal{F})}\right)=O\left(1.0652^{L}\right)$ is because $l(\mathcal{F}) \leq L / 2$.

## E. Analysis of the Main Algorithm

Given a formula $\mathcal{F}$, let $\operatorname{reduced}(\mathcal{F})$ be the output formula of $\operatorname{Reduction}(\mathcal{F})$, and $\operatorname{reduced}_{p}(\mathcal{F})$ be the first formula during the execution of Reduction $(\mathcal{F})$ such that Cases 1-8 are not applicable to the formula. Next we discuss the relationship among the $l$-value of $\mathcal{F}, \operatorname{reduced}_{p}(\mathcal{F})$, and $\operatorname{reduced}(\mathcal{F})$.

Lemma E. $1 \quad l(\mathcal{F}) \geq l\left(\operatorname{reduced}_{p}(\mathcal{F})\right) \geq l(\operatorname{reduced}(\mathcal{F}))$.

Proof. Note that $\operatorname{reduced}(\mathcal{F})=\operatorname{Reduction}\left(\operatorname{reduced}_{p}(\mathcal{F})\right)$. Thus, by lemma C.2, we have $l(\operatorname{reduced}(\mathcal{F})) \leq l\left(\operatorname{reduced}_{p}(\mathcal{F})\right)$. As shown in the proof of lemma C.2, every case in the algorithm $\operatorname{SATSolver}(\mathcal{F})$ decrease the $l$-value of its input formula. So $l(\mathcal{F}) \geq l\left(\operatorname{reduced}_{p}(\mathcal{F})\right)$.

The reason we consider $\operatorname{reduced}_{p}(\mathcal{F})$ is that it is easy to give a bound on the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F})$.

Lemma E. 2 The contributions of any subset of variables is not more than the l-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F})$.

Proof. Let $V$ be the set of all variables in $\mathcal{F}$. Note that $l(\mathcal{F})=\sum_{x \in V} w(x)$ and $l\left(\operatorname{reduced}_{p}(\mathcal{F})\right)=\sum_{x \in V} w^{\prime}(x)$, where $w(x)$ is the frequency weight of $x$ in $\mathcal{F}$ and $w^{\prime}(x)$ is the frequency weight of $x$ in $\operatorname{reduced}_{p}(\mathcal{F})$. By the definition of $\operatorname{reduced}_{p}(\mathcal{F})$, only Cases 1-8 have been applied before we get $\operatorname{reduced}_{p}(\mathcal{F})$. Since Cases 1-8 do not
introduce new variables, we have

$$
l(\mathcal{F})-l\left(\operatorname{reduced}_{p}(\mathcal{F})\right)=\sum_{x \in V}\left(w(x)-w^{\prime}(x)\right)
$$

Therefore, the contribution of $V$ equals to the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F})$. Since Cases 1-8 do not increase the degree of any variable, the contribution of any variable is not negative. Thus the contribution of any subset of variables is not more than the contribution of $V$, i.e., the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F})$.

From now on in this section, let $\mathcal{F}$ be the formula after step 1 in the algorithm $\operatorname{SATSolver}(\mathcal{F})$. Then $\mathcal{F}$ is a reduced formula. In the algorithm $\operatorname{SATSolver}(\mathcal{F})$, we break $\mathcal{F}$ into two formulas $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$ of smaller $l$-values at step 3 or 4 . To give better bound, we are interested in the branching vector from $\mathcal{F}$ to $\operatorname{reduced}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}\left(\mathcal{F}_{2}\right)$, instead of the branching vector from $\mathcal{F}$ to $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$. To give feasible analysis, we focus on the branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)$. This is correct by the following lemma.

Lemma E. 3 The branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)$ is inferior to the branching vector from $\mathcal{F}$ to $\operatorname{reduced}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}\left(\mathcal{F}_{2}\right)$.

Proof. By lemma E.1, $l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)\right) \leq l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)\right)$ and $l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)\right) \leq$ $l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)\right)$. Thus,

$$
\begin{aligned}
& l(\mathcal{F})-l\left(\operatorname{reduced}\left(\mathcal{F}_{1}\right)\right) \geq l(\mathcal{F})-l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)\right), \quad \text { and } \\
& l(\mathcal{F})-l\left(\operatorname{reduced}\left(\mathcal{F}_{2}\right)\right) \geq l(\mathcal{F})-l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)\right) .
\end{aligned}
$$

Note that $\left(l(\mathcal{F})-l\left(\operatorname{reduced}\left(\mathcal{F}_{1}\right)\right), l(\mathcal{F})-l\left(\operatorname{reduced}\left(\mathcal{F}_{2}\right)\right)\right)$ is the branching vector from $\mathcal{F}$ to $\operatorname{reduced}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}\left(\mathcal{F}_{2}\right)$, and $\left(l(\mathcal{F})-l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)\right), l(\mathcal{F})-l\left(\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)\right)\right)$ is the branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)$. By definition, the
branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}_{p}\left(\mathcal{F}_{2}\right)$ is inferior to the branching vector from $\mathcal{F}$ to $\operatorname{reduced}\left(\mathcal{F}_{1}\right)$ and $\operatorname{reduced}\left(\mathcal{F}_{2}\right)$.

Most of the time in the algorithm $\operatorname{SATSolver}(\mathcal{F})$, we break $\mathcal{F}$ into $\mathcal{F}[x]$ and $\mathcal{F}[\bar{x}]$ where $\mathcal{F}$ is of degree $i$ and $x$ is an $i$-variable in $\mathcal{F}$. Let $y$ be a variable contained with variable $x$ in a clause. We can bound the contribution of $y$ to the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$.

Lemma E. 4 Let $y$ be an $i$-variable. The contribution of $y$ to the $l$-value reduction from $\mathcal{F}$ to reduced $_{p}(\mathcal{F}[x])$ is at least
(1) $w_{3}$, if $i=3$ and $y$ is in a clause with literal $x$;
(2) $w_{i}$, if $y$ is contained with literal $\bar{y}$ in a 2-clause;
(3) $\delta_{i}$, if $i>3$ and $y$ is only in one clause with literal $x$, or $2 \delta_{i}$ if $i>3$ and $y$ is in more than one clause with literal $x$.

Proof. (1) The clause containing literal $x$ and $y$ is not in $\mathcal{F}[x]$. Thus $y$ is of degree at most 2 in $\mathcal{F}[x]$, and is not in $\operatorname{reduced}_{p}(\mathcal{F}[x])$ (same proof as in lemma C.3). So the contribution of $y$ is $w_{3}$.
(2) $y$ is a 1 -clause in $\mathcal{F}[x] . y$ is not in $\operatorname{reduced}_{p}(\mathcal{F}[x]$ ) (same proof as in lemma C.3). Thus the contribution of $y$ is $w_{i}$.
(3) If $i>3$ and $y$ is only in one clause with literal $x$, then $y$ is at most a $(i-1)$-variable in $\operatorname{reduced}_{p}(\mathcal{F}[x])$. Thus the contribution of $y$ is at least $\delta_{i}$.

If $i>3$ and $y$ is in more than one clause with literal $x$, then $y$ and literal $x$ are in exactly two clauses in $\mathcal{F}: x y C_{0}$ and $x \bar{y} C_{1}$. Otherwise, Case 9 would be applicable to $\mathcal{F}$, which contradicts that $\mathcal{F}$ is a reduced formula. So the contribution of $y$ is at least $\delta_{i}+\delta_{i-1} \geq 2 \delta_{i}$ when $i>3$.

Let $S$ be the set of variables which are contained with variable $x$ in some clause. We do not include $x$ in $S$. Let $x$ be an $i$-variable in $\mathcal{F}$. Then we can bound the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ with the following calculations:

Step 1: set $c_{x}=w_{i}$ and $c_{y}=0$ for $y \in S$.
Step 2: for each 2-clause $\bar{x} y$ or $\overline{x y}$,
(1) when $i=3$, add $w_{3}$ to $c_{y}$,
(2) when $i>3$, add $w_{i}-\delta_{i}$ to $c_{y}$ if there is a clause $x C$ containing variable $y$, or add $w_{i}$ to $c_{y}$ otherwise.

Step 3: for each clause $x y C$ where $y \in S$,
(1) when $i=3$, add $w_{3}$ to $c_{y}$,
(2) when $i>3$, add $\delta_{i}$ to $c_{y}$.

Step 4: $c=c_{x}+\sum_{y \in S} c_{y}$.
The value $c$ calculated above is the $c$-value from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$. The $c$ value from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$ can be calculated similarly. Next we show that the $c$-value is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F})$.

Lemma E. 5 The c-value is not larger than the contribution of $x$ to the l-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$.

Proof. By lemma E.2, the contribution of $S+x$ is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$. We complete our proof by showing that $c$ is not larger than the contribution of $S+x$. Since $x$ is an $i$-variable in $\mathcal{F}$ and $x$ is not in $\operatorname{reduced}_{p}(\mathcal{F}[x])$, the contribution of $x$ from $\mathcal{F}$ to $\mathcal{F}[x]$ is $w_{i}$. Note that $c=c_{x}+\sum_{y \in S} c_{y}$ at step 4 and $c_{x}=w_{y}$ at step 1 . So we only need to show that $c_{y}$ is not larger than its contribution from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ for all $y \in S$.

At step $1, c_{y}$ is initialized to be 0 .

At step 2 , we only change the $c_{y}$ of variable $y$ which is in 2 -clause $\bar{x} y$ or $\overline{x y}$. For such a variable $y$, there is exactly one clause containing both variable $y$ and literal $\bar{x}$. Then $c_{y}$ is either $w_{i}$ or $w_{i}-\delta_{i}$ after step 2. Moreover, by the rule (2) of lemma E.4, the contribution of $y$ from $\mathcal{F}$ to $\operatorname{reduced}(\mathcal{F}[x])$ is $w_{i}$. So after step2, $c_{y}$ is not larger than the contribution of $y$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$.

At step 3 , we only change the $c_{y}$ of variable $y$ which is in a clause $x C$ where literal $y$ or $\bar{y}$ is in $C$. We consider two cases at step 3 .

Case 1: When $y$ is a 3 -variable, there is no clause $\bar{x} y$ or $\overline{x y}$. So $c_{y}$ is 0 after step 2 , and is $w_{3}$ after step 3 . The contribution of $y$ is at least $w_{3}$ by lemma E.4. Thus $c_{y}$ is not larger than its contribution from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$.

Case 2: When $y$ is an $i$-variable where $i>3$.
If there is a 2 -clause $\bar{x} y$ or $\overline{x y}$, then $c_{y}$ is $w_{i}-\delta_{i}$ after step 2 as shown for step 2 , and there is exactly one clause $x C$ containing variable $y$ and literal $x$ by lemma C.4. Thus $c_{y}$ is $w_{i}$ after step 3. By lemma E.4, the contribution of $y$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ is at least $w_{i}$. So $c_{y}$ is not larger than the contribution of $y$ after step 3.

If there is no 2-clause $\bar{x} y$ or $\overline{x y}$, then $c_{y}$ is 0 after step 2 . For such a variable $y$, there are at most two clauses $x C_{0}$ and $x C_{1}$ containing variable $y$, since $\mathcal{F}$ is a reduced formula, and since Cases 1 and 9 are not applicable to $\mathcal{F}$. Thus $c_{y}$ is not larger than $2 \delta_{i}$ after step 3 , which is not larger than the contribution of $y$ by lemma E.4.

For both cases, we can conclude that $c_{y}$ is not larger than its contribution from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ for all $y \in S$. Then we complete our proof of this lemma.

To give better analysis, some notations are needed.
$n_{1}$ : the number of $3^{+}$-clauses containing literal $x$.
$n_{3}$ : the number of 2 -clauses containing 3 -variables and literal $x$.
$n_{4}$ : the number of 2 -clauses containing 4 -variables and literal $x$. $n_{5}$ : the number of 2 -clauses containing $5^{+}$-variables and literal $x$. $\overline{n_{1}}$ : the number of $3^{+}$-clauses containing literal $\bar{x}$.
$\overline{n_{3}}$ : the number of 2 -clauses containing 3 -variables and literal $\bar{x}$. $\overline{n_{4}}$ : the number of 2-clauses containing 4 -variables and literal $\bar{x}$. $\overline{n_{5}}$ : the number of 2-clauses containing $5^{+}$-variables and literal $\bar{x}$.

$$
\begin{aligned}
& m_{1}=w_{i}+2 n_{1} \delta_{i}+\left(n_{3}+\overline{n_{3}}+\overline{n_{4}}\right) w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} w_{4} . \\
& m_{2}=w_{i}+2 \overline{n_{1}} \delta_{i}+\left(n_{3}+\overline{n_{3}}+n_{4}\right) w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}+n_{5} w_{4} .
\end{aligned}
$$

Recall that $\mathcal{F}$ is a reduced formula, $d(\mathcal{F})=i$ and $x$ is an $i$-variable in $\mathcal{F}$. We have the following lemma:

Lemma E. 6 The value $m_{1}$ is not larger than the l-value reduction from $\mathcal{F}$ to the formula $\operatorname{reduced}_{p}(\mathcal{F}[x])$, and the value $m_{2}$ is not larger than the l-value reduction from $\mathcal{F}$ to the formula reduced ${ }_{p}(\mathcal{F}[\bar{x}])$.

Proof. First we prove that $m_{1}$ is not larger than $c$. By the calculation of the $c$-value, each 2 -clause containing a 3 -variable and literal $\bar{x}$ adds $w_{3}$ to $c$, each 2 -clause containing a 4 -variable and literal $\bar{x}$ adds $w_{3}$ to $c$, and each 2 -clause containing a $i^{+}$-variable and literal adds at least $w_{i}-\delta_{i} \geq w_{4}$ to $c$. where $i \geq 5$. So those 2-clauses containing literal $\bar{x}$ add at least $\left(\overline{n_{3}}+\overline{n_{4}}\right) w_{3}+\overline{n_{5}} w_{4}$. Also each $3^{+}$-clause containing literal $x$ adds at least $2 \delta_{i}$ to $c$, each 2-clause containing a 3 -variable and literal $x$ adds $w_{3}$ to $c$, each 2-clause containing a 4 -variable and literal $x$ adds $\delta_{4}$ to $c$, and each 2 -clause containing a $5^{+}$-variable and literal $x$ adds $\delta_{i} \geq \delta_{5}$ to $c$ since $i \geq 5$. So the clauses containing literal $x$ add at least $2 n_{1} \delta_{i}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}$ to $c$. Thus $c$ is at least $2 n_{1} \delta_{i}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\left(\overline{n_{3}}+\overline{n_{4}}\right) w_{3}+\overline{n_{5}} w_{4}=m_{1}$.

By lemma E.5, $c$ is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$.

So $m_{2}$ is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$. By symmetry, we can prove that $m_{2}$ is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$.

Lemma E. 6 is sufficient for most cases in the following analysis. Sometimes, we may need values better than $m_{1}$. Let
$n_{1,1}$ : the number of 3 -clauses containing literal $x$.
$n_{1,2}$ : the number of $4^{+}$-clauses containing literal $x$.
$\overline{n_{4,1}}$ : the number of 2-clauses containing literal $\bar{x}$ and variable $y$ such that some clause containing both literal $x$ and variable $y$.
$\overline{n_{4,2}}$ : the number of 2-clauses containing literal $\bar{x}$ and variable $y$ such that no clauses containing both literal $x$ and variable $y$.

$$
m_{1}^{\prime}=w_{i}+\left(2 n_{1,1}+3 n_{1,2}\right) \delta_{i}+\left(n_{3}+\overline{n_{3}}+\overline{n_{4,1}}\right) w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{4,2}} w_{4}+\overline{n_{5}} w_{4}
$$

By a proof similar to that for Lemma E.6, we can prove following lemma.

Lemma E. 7 The value $m_{1}^{\prime}$ is not larger than the l-value reduction from $\mathcal{F}$ to the
 from $\mathcal{F}$ to the formula reduced $_{p}(\mathcal{F}[\bar{x}])$.

Now we are ready to analyze the branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ and $\operatorname{reduced}_{p}(\mathcal{F}(\bar{x})$.

Lemma E. 8 Let $\mathcal{F}$ be a reduced formula with $d(\mathcal{F})=i$, and $x$ be an $i$-variable in $\operatorname{reduced}(\mathcal{F})$. Then both $\left(m_{1}, m_{2}\right)$ and $\left(m_{1}^{\prime}, m_{2}\right)$ are inferior to the branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ and $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$.

Proof. By Lemma E.6, $\left(m_{1}, m_{2}\right)$ is inferior to the branching vector from $\mathcal{F}$
to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ and $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$. By Lemma E.7, $\left(m_{1}^{\prime}, m_{2}\right)$ is inferior to the branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ and $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$.

If $x$ is a $(i-1,1)$-literal and no 2 -clause contains $x$, we can have a better branching vector.

Lemma E. 9 Given a reduced formula $\mathcal{F}$ of degree $i$, and an $(i-1,1)$-literal $x$ in $\mathcal{F}$ such that no 2-clause contains $x$, let

$$
m_{1}^{\prime}=w_{i}+2 n_{1} \delta_{i}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, \text { and } m_{2}^{\prime}=w_{i}+3 w_{3}
$$

Then $\left(m_{1}^{\prime}, m_{2}^{\prime}\right)$ is inferior to the branching vector from $\mathcal{F}$ to reduced $_{p}(\mathcal{F}[x])$ and reduced $_{p}\left(\mathcal{F}[\bar{x})\right.$. Moreover, $\overline{n_{1}}=1$ and $\overline{n_{3}}=\overline{n_{4}}=\overline{n_{5}}=0$.

Proof. First we show that $\overline{n_{1}}=1$ and $\overline{n_{3}}=\overline{n_{4}}=\overline{n_{5}}=0$. Let $\bar{x} C$ be the clause containing literal $\bar{x}$. Then we must have $|\mathcal{C}| \geq 2$ : if $|\mathcal{C}|=0$, then Case 5 in the algorithm Reduction would be applicable to $\bar{x}$ in $\mathcal{F}$, and if $|\mathcal{C}|=1$, then Case 10 in the algorithm would be applicable to $\bar{x} \mathcal{C}$ in $\mathcal{F}$. So we must have $|\mathcal{C}| \geq 2$. Since $x$ is an $(i-1,1)$-literal, we have $\overline{n_{1}}=1$, and $\overline{n_{3}}=\overline{n_{4}}=\overline{n_{5}}=0$.

By the definition of inferior vectors, $\left(m_{1}^{\prime}, m_{2}^{\prime}\right)$ is inferior to the branching vector from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ and $\operatorname{reduced}_{p}\left(\mathcal{F}[\bar{x})\right.$, once we show that $m_{1}^{\prime}$ is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ and $m_{2}^{\prime}$ is not larger than the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$.

We first show that $m_{1}^{\prime}$ is not larger than the $l$-value reduction from $\mathcal{F}$ to the
 $n_{5} \delta_{5}=m_{1}^{\prime}$. By lemma E.6, $m_{1}^{\prime}=m_{1}$ is not larger than the $l$-value reduction from $\mathcal{F}$ to reduced $_{p}(\mathcal{F}[x])$.

Now we show that $m_{2}^{\prime}$ is at not larger than the $l$-value reduction from $\mathcal{F}$ to
$\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$.
When $x=0$, then all literals in $\mathcal{C}$ must be false in $\mathcal{F}[\bar{x}]$, since reduced $_{p}(\mathcal{F}[x])$ is not satisfiable. i.e., all variables in $x+C$ are not in $\operatorname{reduced}_{p}(\mathcal{F}[\bar{x}])$. Since $\mathcal{F}$ is a reduced formula, all variables in $x+C$ are $3^{+}$-variables. By lemma E. 2 , the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}(\mathcal{F}[x])$ is at least $w_{i}+|\mathcal{C}| w_{3}$. If $|\mathcal{C}| \geq 3$, then the $l$-value reduction is at least $w_{i}+3 w_{3}$. Otherwise $|\mathcal{C}|=2$. Let $\mathcal{C}=z_{1} z_{2}$. If both $z_{1}$ and $z_{2}$ are $4^{+}$-variables, then the $l$-value reduction is at least $2 w_{4} \geq 3 w_{3}$. Now we only need consider that at least one of $z_{1}$ and $z_{2}$, say $z_{1}$, is a 3 -variable. Then literals $x$ and $\overline{z_{2}}$ can not co-occur with literal $\overline{z_{1}}$ in any clause. Otherwise, Case 4 (on $z_{1}$ ) in the algorithm reduction would be applicable to $\mathcal{F}$. According to lemma C.5, only literal $z_{2}$ can co-occur with literal $\overline{z_{1}}$ in a $3^{+}$-clause. Now we consider two cases.

Case 1: If $\overline{z_{1}} z_{2}$ is in a clause $\overline{z_{1}} z_{2} y C$, then $z_{2}$ must be a $4^{+}$-variable. Otherwise, the clause $\overline{z_{1}} z_{2} y C$ can not exist by lemma C.5. Since $y$ can not be the $(i-1,1)$-literal $x$, the $l$-value reduction is at least $w_{i}+w_{3}+w_{4}+\delta_{i}>w_{i}+3 w_{3}$ when $x=0\left(z_{1}=0\right.$, $z_{2}=0$ and $\left.y=1\right)$.

Case 2: If literal $z_{2}$ does not co-occur with literal $\overline{z_{1}}$ in any clause, then there is a clause $\overline{z_{1}} y_{1} \mathcal{C}_{1}$ such that $y_{1}$ is not in $\left\{x, \bar{x}, z_{2}, \overline{z_{2}}\right\}$. Also there is a clause $\overline{z_{2}} y_{2} \mathcal{C}_{2}$ such that $y_{2}$ is not in $\left\{x, \bar{x}, z_{2}, \overline{z_{2}}\right\}$. Otherwise, there is only a clause $\overline{z_{2}} z_{1}$. The Case 6 is applicable, thus contradicting that $\mathcal{F}$ is a reduced formula. Therefore, the $l$-value reduction is at least $w_{i}+w_{3}+w_{3}+2 \delta_{j}>w_{i}+2 w_{3}+1>w_{i}+3 w_{3}$.

## 1. Analysis for Degree-4 Formulas

Suppose that $d(\mathcal{F})=4$ and $x$ is a 4 -variable, let $d_{1}$ be the degree of literal $x$ and $d_{0}$ be the degree of $\bar{x}$. We then consider all combinations of $d_{1}$ and $d_{0}$ such that $d_{1}+d_{0}=4$. Without loss of generality, we assume that $d_{1} \geq d_{0}$. Since $d(\mathcal{F})=4$, we have $n_{5}=0$ and $\overline{n_{5}}=0$.
a. $\quad d_{0}=1$

Note that $n_{1}=3-n_{3}-n_{4}$. By lemma E. $9, \overline{n_{3}}=\overline{n_{4}}=0$.
Case 1: $n_{3}+n_{4}=0$. Then $n_{1}=3$. By lemma E.9, the branching vector is not inferior to

$$
\left(w_{3}+7 \delta_{4}, 4 w_{3}+\delta_{4}\right)
$$

Case 2: $1 \leq n_{3}+n_{4} \leq 3$. By lemma E.8, the branching vector is not inferior to $\left(m_{1}, m_{2}\right)$, which is:

$$
\begin{aligned}
& \left(w_{3}+6 \delta_{4}, 4 w_{3}+2 \delta_{4}\right) \text { when } n_{3}=0 \text { and } n_{4}=1, \\
& \left(w_{3}+5 \delta_{4}, 5 w_{3}+3 \delta_{4}\right) \text { when } n_{3}=0 \text { and } n_{4}=2, \\
& \left(w_{3}+4 \delta_{4}, 6 w_{3}+4 \delta_{4}\right) \text { when } n_{3}=0 \text { and } n_{4}=3, \\
& \left(2 w_{3}+5 \delta_{4}, 4 w_{3}+\delta_{4}\right) \text { when } n_{3}=1 \text { and } n_{4}=0, \\
& \left(2 w_{3}+4 \delta_{4}, 5 w_{3}+2 \delta_{4}\right) \text { when } n_{3}=1 \text { and } n_{4}=1, \\
& \left(2 w_{3}+3 \delta_{4}, 6 w_{3}+3 \delta_{4}\right) \text { when } n_{3}=1 \text { and } n_{4}=2, \\
& \left(3 w_{3}+3 \delta_{4}, 5 w_{3}+\delta_{4}\right) \text { when } n_{3}=2 \text { and } n_{4}=0, \\
& \left(3 w_{3}+2 \delta_{4}, 6 w_{3}+2 \delta_{4}\right) \text { when } n_{3}=2 \text { and } n_{4}=1, \\
& \left(4 w_{3}+\delta_{4}, 6 w_{3}+\delta_{4}\right) \text { when } n_{3}=3 \text { and } n_{4}=0 .
\end{aligned}
$$

b. $\quad d_{0}=2$

Case 1: Two 2-clauses $\bar{x} y_{1}$ and $\bar{x} y_{2}$ contains literal $\bar{x}$.
Case 1.1: Two $3^{+}$-clauses contain literal $x$, i.e., $n_{3}=2$ and $n_{4}=0$.
If one clause is a $4^{+}$-clauses, i.e., $n_{1,1}=n_{1,2}=1$, then by lemma E.8, the branching vector is not inferior to $\left(m_{1}^{\prime}, m_{2}\right)$, which is

$$
\left(3 w_{3}+6 \delta_{4}, w_{3}+3 \delta_{4}\right)
$$

If both are 3 -clause, i.e., $n_{1}=2$, then we consider following cases:
Subcase 1.1.1: both $y_{1}$ and $y_{2}$ are 4 -variable, i.e., $\overline{n_{4}}=2$ and $\overline{n_{3}}=0$. Then
at most one of $\overline{y_{1}}$ and $\overline{y_{2}}$ co-occur with literal $x$. Otherwise, Case 9 in the algorithm Reduction would be applicable. If one of $y_{1}$ and $y_{2}$ co-occurs with literal $x$, i.e., $\overline{n_{4,1}}=\overline{n_{4,2}}=1$, then by lemma E. 8 , the branching vector is not inferior to

$$
\left(3 w_{3}+6 \delta_{4}, w_{3}+3 \delta_{4}\right)
$$

If none of $y_{1}$ and $y_{2}$ co-occurs with literal $x$, i.e., $\overline{n_{4}}=2$ and $\overline{n_{3}}=0$, then by lemma E.8, the branching vector is not inferior to

$$
\left(w_{3}+7 \delta_{4}, w_{3}+3 \delta_{4}\right)
$$

Subcase 1.1.2 $y_{1}$ is a 4 -variable and $y_{2}$ is a 3 -variable, i.e., $\overline{n_{3}}=\overline{n_{4}}=1$.
If $y_{1}$ does not co-occur with literal $x$, then by lemma E. 8 , the branching vector is not inferior to $\left(m_{1}^{\prime}, m_{2}\right)$, which is

$$
\left(3 w_{3}+6 \delta_{4}, 2 w_{3}+2 \delta_{4}\right) .
$$

Otherwise, literal $y_{1}$ and $x$ are in a clause $x y_{1} z_{1}$ since both clauses containing literal $x$ are 3 -clause in Subcase 1.1.1-1.1.3. Let $x z_{2} z_{3}$ be the other 3 -clause containing literal $x$. Note that $y_{2}$ must be a 2 -literal. Otherwise, Case 10 in the algorithm Reduction would be applied. Let the other two clauses containing variable $y$ be $y_{2} C_{1}$ and $\overline{y_{2}} C_{0}$. By lemma E.9, $\left|C_{0}\right| \geq 2$. Note that variable $x$ and $y_{2}$ are not in $C_{0}$. With these given formulas, step 4.1 in algorithm $\operatorname{SATSolver}(\mathcal{F})$ would be executed.

First we consider the case that $C_{1}$ is a $2^{+}$-clause.
If $C_{0}$ is true, then we can set $y_{2}=1$, then $x=y_{1}$. Clauses $x \overline{y_{1}} z_{1}, \bar{x} y_{1}, x y_{2}$, and $y_{2} C_{1}$ are eliminated and are not in $\mathcal{F}\left[C_{0}=\right.$ true $]$, clause $x z_{1} z_{2}$ in $\mathcal{F}$ becomes $y_{1} z_{1} z_{2}$ in $\mathcal{F}\left[C_{0}=\right.$ true $]$, and clause $\bar{x} C_{0}$ in $\mathcal{F}$ becomes $C_{0}$ in $\mathcal{F}\left[C_{0}=\right.$ true $]$. Note that $y_{1}$ is a 3 -variable in $\mathcal{F}\left[C_{0}=\right.$ true $]$. And it is obvious that only the degrees of $x, y_{2}, y_{1}$ and variables in $C_{1}$ in $\mathcal{F}\left[C_{0}=\right.$ true $]$ are different to those in $\mathcal{F}$. Thus the $l$-value reduction from $\mathcal{F}$ to $\mathcal{F}\left[C_{0}=\right.$ true $]$ is the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{1}$ from $\mathcal{F}$ to $\mathcal{F}\left[C_{0}=\right.$ true $]$. The $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is the summation of the $l$-value reduction from $\mathcal{F}$ to $\mathcal{F}\left[C_{0}=\right.$ true $]$ and the $l$-value
reduction from $\mathcal{F}\left[C_{0}=\operatorname{true}\right]$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\operatorname{true}\right]\right)$. By lemma E.2, the $l$-value reduction from $\mathcal{F}\left[C_{0}=\right.$ true $]$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is at least the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{1}$ from $\mathcal{F}\left[C_{0}=\right.$ true $]$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$. So the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is at least the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{1}$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$. If variable $y_{1}$ is not in $C_{1}$, the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{1}$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is at least $2 w_{3}+4 \delta_{4}$, no matter whether there are 3 -variables or not in $C_{1}$. If variable $y_{1}$ is in $C_{1}$, the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{1}$ is still at least $2 w_{3}+4 \delta_{4}$ since $y_{1}$ is a 4 -variable. For both cases, the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is at least $w_{2}$ (or $\left.2 w_{3}\right)$.

If $C_{0}$ is false, we require that all literals in $C_{0}$ are false, since we work on $\mathcal{F}\left[C_{0}=\right.$ false $]$ only after $\mathcal{F}\left[C_{0}=\right.$ true $]$ is not satisfiable. To satisfy $\mathcal{F}, y_{2}$ must be 0 , which results in $x=0$ from clause $\bar{x} y_{2}$. Hence clauses $\bar{x} y_{1}, \bar{x} y_{2}$ and $\overline{y_{2}} \mathcal{C}_{0}$ are not in $\mathcal{F}\left[C_{0}=\right.$ false $]$. Similarly as above, we can show that the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $\left.]\right)$ is at least the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{0}$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $\left.]\right)$. If variable $y_{1}$ is not in $C_{0}$, the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{0}$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $\left.]\right)$ is at least $4 w_{3}+2 \delta_{4}$, since (1) the contribution of $x$ is $w_{4},(2)$ the contribution of $y_{2}$ is $w_{3},(3)$ the contribution of $C_{0}$ is at least $2 w_{3}$ (recall that $\left|C_{0}\right| \geq 2$ and all variables in $\mathcal{F}$ are $3^{+}$-variables), and (4) the contribution of $y_{1}$ is $\delta_{4}$. If variable $y_{1}$ is in $C_{0}$, the contribution of $\left\{x, y_{2}, y_{1}\right\} \cup C_{0}$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $\left.]\right)$ is still at least $4 w_{3}+2 \delta_{4}$, since (1) the contribution of $x$ and $y_{2}$ is $w_{3}+w_{4}$, and (2) the contribution of $C_{0}$ is at least $w_{4}+w_{3}\left(y_{1}\right.$ is a 4 -variable). For both cases, the $l$-value reduction is at least $4 w_{3}+2 \delta_{4}$.

We conclude that when $C_{1}$ is a $2^{+}$-clause, the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is not inferior to

$$
\left(2 w_{4}+4 \delta_{4}, 4 w_{3}+2 \delta_{4}\right)
$$

Next we consider the case that $\mathcal{C}_{1}$ is a 1-clause $u$.

If $\mathcal{C}_{0}$ is true, then we can prove that the $l$-value reduction from $\mathcal{F}$ to the formula $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ true $\left.]\right)$ is at least $2 w_{3}+3 \delta_{4}$, using similar proof for the first part of the case when $C_{1}$ is a $2^{+}$-clause.

If $\mathcal{C}_{0}$ is false, we require that $C_{0}$ be false. To satisfy $\mathcal{F}$, clause $\overline{y_{2}} u$ requires $y_{2}=0$, which results in $x=0$ because of clause $\bar{x} y_{2}$ and $u=1$ because of clause $y_{2} \mathcal{C}_{1}$. Then clauses $\bar{x} y_{1}, \overline{y_{2}} C_{0}$ and $y_{2} \mathcal{C}_{1}$ are removed from $\mathcal{F}$. Note that only variables in $\left\{x, y_{2}, y_{1}, u\right\} \cup C_{0}$ change their degrees from $\mathcal{F}$ to $\mathcal{F}\left[C_{0}=\right.$ false $]$. Similarly as above, we can show that the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=f a l s e\right]\right)$ is at least the contribution of $\left\{x, y_{2}, y_{1}, u\right\} \cup C_{0}$ from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $)$. The contribution of $x$ and $y_{2}$ is $w_{4}+w_{3}=2 w_{3}+\delta_{4}$. Note that $u$ is not in $C_{0}$ by lemma C. 4 , since $y_{2}$ is a 3 -variable. Also $u$ can not be $x$ or $y_{2}$. If variable $y_{1}$ is not in $\{u\} \cup C_{0}$, then the contribution of $\left\{y_{1}, u\right\} \cup C_{0}$ is at least $3 w_{3}+\delta_{4}$ since $\left|C_{0}\right| \geq 2$ and $y_{1}$ is a 4 -variable. If variable $y_{1}$ is in $\{u\} \cup C_{0}$, the contribution of $\left\{y_{1}, u\right\} \cup C_{0}$ is at least $w_{4}+2 w_{3}=3 w_{3}+\delta_{4}$ since $y_{1}$ is a 4 -variable. For both cases, the $l$-value reduction from $\mathcal{F}$ to $\operatorname{reduced}_{p}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $\left.]\right)$ is at least $2 w_{3}+\delta_{4}+3 w_{3}+\delta_{4}=5 w_{3}+4 \delta_{4}$.

We conclude that when $C_{1}$ is a 1 -clause, the $l$-value reduction from $\mathcal{F}$ to the formula $^{\operatorname{reduced}_{p}}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $)$ is not inferior to

$$
\left(2 w_{3}+4 \delta_{4}, 5 w_{3}+4 \delta_{4}\right) .
$$

Subcase 1.1.3: both $y_{1}$ and $y_{2}$ are 3 -variables, i.e., $\overline{n_{3}}=2$ and $\overline{n_{4}}=0$. By lemma E.8, the branching vector is not inferior to

$$
\left(3 w_{3}+5 \delta_{4}, 3 w_{3}+\delta_{4}\right)
$$

Case 1.2: One $3^{+}$-clause contains literal $x$, i.e., $n_{1}=1$ and $n_{3}+n_{4}=1$. Let $x z$ be the 2-clause.

If $z$ is a 3 -variable, $n_{3}=1$ and $n_{4}=0$. By lemma E. 8 , the branching vector is not inferior to $\left(m_{1}, m_{2}\right)=\left(w_{4}+2 n_{1} \delta_{4}+\left(n_{3}+\overline{n_{3}}+\overline{n_{4}}\right) w_{3}+n_{4} \delta_{4}, w_{4}+2 \overline{n_{1}} \delta_{4}+\left(n_{3}+\right.\right.$ $\left.\left.\overline{n_{3}}+n_{4}\right) w_{3}+\overline{n_{4}} \delta_{4}\right)=\left(w_{4}+2 \delta_{4}+3 w_{3}, w_{4}+\left(1+\overline{n_{3}}\right) w_{3}+\overline{n_{4}} \delta_{4}\right)$, which is not inferior to

$$
\left(4 w_{3}+3 \delta_{4}, 2 w_{3}+3 \delta_{4}\right) \text { when } \overline{n_{3}}=0 \text { and } \overline{n_{4}}=2 .
$$

If $z$ is a 4 -variable, $n_{3}=0$ and $n_{4}=1$. By lemma C.4, variable $z$ can not be neither variable $y_{1}$ nor $y_{2}$. Then $n_{4,1}=0, n_{4,2}=1$ and $n_{3}+n_{4,1}=1$. By lemma E.8, the branching vector is not inferior to $\left(m_{1}, m_{2}^{\prime}\right)=$ $\left(w_{4}+2 n_{1} \delta_{4}+\left(n_{3}+\overline{n_{3}}+\overline{n_{4}}\right) w_{3}+n_{4} \delta_{4}, w_{4}+\left(2 \overline{n_{1,1}}+3 \overline{n_{1,2}}\right) \delta_{4}+\left(n_{3}+\overline{n_{3}}+n_{4,1}\right) w_{3}+\right.$ $\left.\overline{n_{4}} \delta_{4}+n_{4,2} w_{4}\right)=\left(w_{4}+3 \delta_{4}+2 w_{3}, 2 w_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} w_{4}\right)$, which is not inferior to

$$
\left(3 w_{3}+4 \delta_{4}, 2 w_{3}+4 \delta_{4}\right) \text { when } \overline{n_{3}}=0 \text { and } \overline{n_{4}}=2 .
$$

Case 1.3: Two 2-clauses contain literal $x$, i.e., $n_{3}=2$ and $n_{4}=0$.
By lemma E.8, the branching vector is not inferior to
$\left(w_{4}+\left(n_{3}+\overline{n_{3}}+\overline{n_{4}}\right) w_{3}+n_{4} \delta_{4}, w_{4}+\left(n_{3}+\overline{n_{3}}+n_{4}\right) w_{3}+\overline{n_{4}} \delta_{4}\right)=\left(w_{4}+\left(n_{3}+2\right) w_{3}+\right.$ $\left.n_{4} \delta_{4}, w_{4}+\left(2+\overline{n_{3}}\right) w_{3}+\overline{n_{4}} \delta_{4}\right)$, which is not inferior to

$$
\left(3 w_{3}+3 \delta_{4}, 3 w_{3}+3 \delta_{4}\right) \text { when } n_{3}=\overline{n_{3}}=0
$$

Case 2: One 2-clauses $\bar{x} y_{1}$ and a $3^{+}$-clause $\overline{x y_{2} y_{3}} \mathcal{C}_{1}$ with $\left|\mathcal{C}_{1}\right| \geq 0$ contain literal $\bar{x}$, i.e., $\overline{n_{1}}=1$ and $\overline{n_{3}}+\overline{n_{4}}=1$.

Case 2.1: Two $3^{+}$-clauses contain literal $x$, i.e., $n_{3}+n_{4}=2$.
By lemma E.8, the branching vector is not inferior to

$$
\begin{aligned}
& \left(2 w_{3}+5 \delta_{4}, w_{3}+4 \delta_{4}\right) \text { when } \overline{n_{3}}=0, \\
& \left(2 w_{3}+5 \delta_{4}, 2 w_{3}+3 \delta_{4}\right) \text { when } \overline{n_{3}}=1 .
\end{aligned}
$$

Case 2.2: One $3^{+}$-clause contains literal $x$.
We have $n_{3}+n_{4}=1$ and $\overline{n_{3}}+\overline{n_{4}}=1$. By lemma E.8, the branching vector is not inferior to $t=\left(w_{4}+2 \delta_{4}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}}\left(w_{4}-\delta_{4}\right), w_{4}+2 \delta_{4}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+\right.$ $\left.n_{4}\left(w_{4}-\delta_{4}\right)+\overline{n_{4}} \delta_{4}\right)=\left(w_{4}+2 \delta_{4}+w_{3}+n_{3} w_{3}+n_{4} \delta_{4}, w_{4}+2 \delta_{4}+w_{3}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}\right)$. It is clear that $t$ is not inferior to

$$
\left(2 w_{3}+4 \delta_{4}, 2 w_{3}+4 \delta_{4}\right)\left(\text { by setting } n_{3}=\overline{n_{3}}=1\right)
$$

Case 2.3: Two 2-clauses contain literal $x$.
This case is symmetric to Case 1.2.

Case 3: Two $3^{+}$-clauses $\bar{x} \mathcal{C}_{1}$ and $\bar{x} \mathcal{C}_{2}$ contain literal $\bar{x}$, i.e., $\overline{n_{3}}=\overline{n_{4}}=0$ and $\overline{n_{1}}=2$.
Case 3.1: Two $3^{+}$-clauses contain literal $x$.
We have $n_{1}=2, n_{3}=n_{4}=0$. By lemma E. 8 , the branching vector is not inferior to

$$
\left(w_{3}+5 \delta_{4}, w_{3}+5 \delta_{4}\right)
$$

Case 3.2: One $3^{+}$-clause contains literal $x$.
This case is symmetric to Case 2.1.
Case 3.3: Two 2-clauses contain literal $x$.
This case is symmetric to Case 1.1.

## 2. Analysis for Degree-5 Formulas

Suppose variable $x$ is of degree 5 , let $d_{1}$ be the degree of literal $x$ and $d_{0}$ be the degree of $\bar{x}$. We then consider all combinations of $d_{1}$ and $d_{0}$ such that $d_{1}+d_{0}=5$. W.l.o.g, it can be assumed that $d_{1} \geq d_{0}$.
a. $\quad d_{0}=1$
$3^{+}$-clause $\bar{x} y_{1} y_{2} \mathcal{C}$ contains literal $\bar{x}$ where $|\mathcal{C}| \geq 0$. We have $4 \geq n_{3}+n_{4}+n_{5} \geq 0$, $n_{1}=5-n_{3}-n_{4}-n_{5}$.

Case $1 n_{3}+n_{4}+n_{5}=0 . x$ is a $(4,1)$-literal. By lemma E.9, the branching vector is not inferior to

$$
\left(w_{5}+8 \delta_{5}, w_{5}+3 w_{3}\right)
$$

Case $24 \geq n_{3}+n_{4}+n_{5}=i \geq 1$. By lemma E.8, the branching vector is not inferior to $t=\left(w_{5}+2 \delta_{5}\left(4-n_{3}-n_{4}-n_{5}\right)+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+2 w_{3}+n_{3} w_{3}+\right.$ $n_{4}\left(w_{4}-\delta_{4}\right)+n_{5}\left(w_{5}-\delta_{5}\right)$. Note that $t$ is not inferior to $t^{\prime}=\left(w_{5}+\delta_{5}\left(8-2\left(n_{3}+n_{4}\right)-\right.\right.$ $\left.\left.n_{5}\right)+\left(n_{3}+n_{4}\right) \delta_{4}, w_{5}+2 w_{3}+\left(n_{3}+n_{4}+n_{5}\right) w_{3}+n_{5} \delta_{4}\right)$.

When $n_{5}=0, t^{\prime}$ is

$$
\begin{aligned}
& \left(w_{5}+6 \delta_{5}+\delta_{4}, w_{5}+3 w_{3}\right) \text { when } n_{3}+n_{4}+n_{5}=1, \\
& \left(w_{5}+4 \delta_{5}+2 \delta_{4}, w_{5}+4 w_{3}\right) \text { when } n_{3}+n_{4}+n_{5}=2, \\
& \left(w_{5}+2 \delta_{5}+3 \delta_{4}, w_{5}+5 w_{3}\right) \text { when } n_{3}+n_{4}+n_{5}=3, \\
& \left(w_{5}+4 \delta_{4}, w_{5}+6 w_{3}\right) \text { when } n_{3}+n_{4}+n_{5}=4 .
\end{aligned}
$$

When $n_{5}=1, t^{\prime}$ is

$$
\begin{aligned}
& \left(w_{5}+7 \delta_{5}, w_{5}+3 w_{3}+\delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=1, \\
& \left(w_{5}+5 \delta_{5}+\delta_{4}, w_{5}+4 w_{3}+\delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=2, \\
& \left(w_{5}+3 \delta_{5}+2 \delta_{4}, w_{5}+5 w_{3}+\delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=3, \\
& \left(w_{5}+3 \delta_{4}, w_{5}+6 w_{3}\right)+\delta_{4} \text { when } n_{3}+n_{4}+n_{5}=4 .
\end{aligned}
$$

When $n_{5}=2, t^{\prime}$ is

$$
\begin{aligned}
& \left(w_{5}+6 \delta_{5}, w_{5}+3 w_{3}+2 \delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=2, \\
& \left(w_{5}+4 \delta_{5}+\delta_{4}, w_{5}+4 w_{3}+2 \delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=3, \\
& \left(w_{5}+2 \delta_{5}+2 \delta_{4}, w_{5}+5 w_{3}+2 \delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=4 .
\end{aligned}
$$

When $n_{5}=3, t^{\prime}$ is

$$
\begin{aligned}
& \left(w_{5}+5 \delta_{5}, w_{5}+3 w_{3}+3 \delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=3 \\
& \left(w_{5}+3 \delta_{5}+\delta_{4}, w_{5}+4 w_{3}+3 \delta_{4}\right) \text { when } n_{3}+n_{4}+n_{5}=4 .
\end{aligned}
$$

When $n_{5}=4, n_{3}=n_{4}=0 . t^{\prime}$ is

$$
\left(w_{5}+4 \delta_{5}, w_{5}+6 w_{3}+4 \delta_{4}\right) .
$$

b. $\quad d_{0}=2$

Case 1: Two 2-clauses $\bar{x} y_{1}$ and $\bar{x} y_{2}$ contain literal $\bar{x}$, i.e., $\overline{n_{1}}=0$ and $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=2$.
Case 1.1: Three $3^{+}$-clauses contains literal $x$, i.e., $n_{1}=3, n_{3}=n_{4}=n_{5}=0$.
By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+6 \delta_{5}+\overline{n_{3}} w_{3}+\overline{n_{4}}\left(w_{4}-\delta_{4}\right)+\overline{n_{5}}\left(w_{5}-\delta_{5}\right), w_{5}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$. Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=2, t=\left(w_{5}+6 \delta_{5}+2 w_{3}+\overline{n_{5}} \delta_{4}, w_{5}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$ and is not inferior to
$t^{\prime}=\left(w_{5}+6 \delta_{5}+2 w_{3}+\overline{n_{5}} \delta_{4}, w_{5}+\left(\overline{n_{3}}+\overline{n_{4}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$, which is:

$$
\begin{aligned}
& \left(w_{5}+6 \delta_{5}+2 w_{3}, w_{5}+2 \delta_{4}\right) \text { when } \overline{n_{5}}=0, \\
& \left(w_{5}+6 \delta_{5}+2 w_{3}+\delta_{4}, w_{5}+\delta_{4}+\delta_{5}\right) \text { when } \overline{n_{5}}=1, \\
& \left(w_{5}+6 \delta_{5}+2 w_{3}+2 \delta_{4}, w_{5}+2 \delta_{5}\right) \text { when } \overline{n_{5}}=2 .
\end{aligned}
$$

Case 1.2: two $3^{+}$-clauses and one 2 -clause $x z_{1}$ contain literal $x$, i.e., $n_{1}=2$ and $n_{3}+n_{4}+n_{5}=1$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+4 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}} w_{3}+n_{5} \delta_{5}+\overline{n_{5}} w_{4}, w_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} w_{3}+\right.$ $\left.\overline{n_{4}} \delta_{4}+n_{5} w_{4}+\overline{n_{5}} \delta_{5}\right)$. Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=2$ and $n_{3}+n_{4}+n_{5}=1$, we have
$t=\left(w_{5}+4 \delta_{5}+2 w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+w_{3}+n_{5} \delta_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$.
Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+4 \delta_{5}+2 w_{3}+\left(1-n_{5}\right) \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+w_{3}+n_{5} \delta_{4}+\left(2-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$,
which is

$$
\begin{aligned}
& \left(w_{5}+4 \delta_{5}+2 w_{3}+\delta_{4}, w_{5}+w_{3}+2 \delta_{4}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+4 \delta_{5}+2 w_{3}+2 \delta_{4}, w_{5}+w_{3}+\delta_{4}+\delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+4 \delta_{5}+2 w_{3}+3 \delta_{4}, w_{5}+w_{3}+2 \delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=2, \\
& \left(w_{5}+5 \delta_{5}+2 w_{3}, w_{5}+w_{3}+3 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+5 \delta_{5}+2 w_{3}+\delta_{4}, w_{5}+w_{3}+2 \delta_{4}+\delta_{5}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+5 \delta_{5}+2 w_{3}+2 \delta_{4}, w_{5}+w_{3}+\delta_{4}+3 \delta_{5}\right) \text { when } n_{5}=3 \text { and } \overline{n_{5}}=2 .
\end{aligned}
$$

Case 1.3: one $3^{+}$-clauses and two 2-clauses $x z_{1}, x z_{2}$ contain literal $x$, i.e., $n_{1}=1$ and $n_{3}+n_{4}+n_{5}=2$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+2 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}} w_{3}+n_{5} \delta_{5}+\overline{n_{5}} w_{4}, w_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} w_{3}+\right.$ $\left.\overline{n_{4}} \delta_{4}+n_{5} w_{4}+\overline{n_{5}} \delta_{5}\right)$. Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=2$ and $n_{3}+n_{4}+n_{5}=2$, we have $t=\left(w_{5}+2 \delta_{5}+2 w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 w_{3}+n_{5} \delta_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$.

Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+2 \delta_{5}+2 w_{3}+\left(2-n_{5}\right) \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 w_{3}+n_{5} \delta_{4}+\left(2-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$,
which is

$$
\begin{aligned}
& \left(w_{5}+2 \delta_{5}+2 w_{3}+2 \delta_{4}, w_{5}+2 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+2 \delta_{5}+2 w_{3}+3 \delta_{4}, w_{5}+2 w_{3}+\delta_{4}+\delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+2 \delta_{5}+2 w_{3}+4 \delta_{4}, w_{5}+2 w_{3}+2 \delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=2, \\
& \left(w_{5}+3 \delta_{5}+2 w_{3}+\delta_{4}, w_{5}+2 w_{3}+3 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+3 \delta_{5}+2 w_{3}+2 \delta_{4}, w_{5}+2 w_{3}+2 \delta_{4}+\delta_{5}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+3 \delta_{5}+2 w_{3}+3 \delta_{4}, w_{5}+2 w_{3}+\delta_{4}+2 \delta_{5}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=2, \\
& \left(w_{5}+4 \delta_{5}+2 w_{3}, w_{5}+2 w_{3}+4 \delta_{4}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+4 \delta_{5}+2 w_{3}+\delta_{4}, w_{5}+2 w_{3}+3 \delta_{4}+\delta_{5}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+4 \delta_{5}+2 w_{3}+2 \delta_{4}, w_{5}+2 w_{3}+2 \delta_{4}+2 \delta_{5}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=2 .
\end{aligned}
$$

Case 1.4: Three 2-clauses $x z_{1}, x z_{2}$ and $x z_{3}$ contain literal $x$, i.e., $n_{1}=0$ and $n_{3}+n_{4}+n_{5}=3$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}} w_{3}+n_{5} \delta_{5}+\overline{n_{5}} w_{4}, w_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} w_{3}+\overline{n_{4}} \delta_{4}+\right.$ $\left.n_{5} w_{4}+\overline{n_{5}} \delta_{5}\right)$. Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=2$ and $n_{3}+n_{4}+n_{5}=3$, we have
$t=\left(w_{5}+2 w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+3 w_{3}+n_{5} \delta_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$.
Note that $t$ is not inferior to
$t=\left(w_{5}+2 w_{3}+\left(3-n_{5}\right) \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+3 w_{3}+n_{5} \delta_{4}+\left(2-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$, which is

$$
\left(w_{5}+2 w_{3}+3 \delta_{4}, w_{5}+3 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=0,
$$

$$
\left(w_{5}+2 w_{3}+4 \delta_{4}, w_{5}+3 w_{3}+\delta_{4}+\delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=1,
$$

$$
\left(w_{5}+2 w_{3}+5 \delta_{4}+2 \delta_{5}, w_{5}+3 w_{3}+2 \delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=2,
$$

$$
\left(w_{5}+2 w_{3}+2 \delta_{4}+\delta_{5}, w_{5}+3 w_{3}+3 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=0
$$

$$
\left(w_{5}+2 w_{3}+3 \delta_{4}+\delta_{5}, w_{5}+3 w_{3}+2 \delta_{4}+\delta_{5}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=1
$$

$$
\left(w_{5}+2 w_{3}+3 \delta_{4}+\delta_{5}, w_{5}+3 w_{3}+\delta_{4}+2 \delta_{5}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=2,
$$

$$
\left(w_{5}+2 w_{3}+\delta_{4}+2 \delta_{5}, w_{5}+3 w_{3}+4 \delta_{4}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=0
$$

$$
\begin{aligned}
& \left(w_{5}+2 w_{3}+2 \delta_{4}+2 \delta_{5}, w_{5}+3 w_{3}+3 \delta_{4}+\delta_{5}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+2 w_{3}+3 \delta_{4}+2 \delta_{5}, w_{5}+3 w_{3}+2 \delta_{4}+2 \delta_{5}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=2, \\
& \left(w_{5}+2 w_{3}+3 \delta_{5}, w_{5}+3 w_{3}+5 \delta_{4}\right) \text { when } n_{5}=3 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+2 w_{3}+3 \delta_{5}+\delta_{4}, w_{5}+3 w_{3}+4 \delta_{4}+\delta_{5}\right) \text { when } n_{5}=3 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+2 w_{3}+3 \delta_{5}+2 \delta_{4}, w_{5}+3 w_{3}+3 \delta_{4}+2 \delta_{5}\right) \text { when } n_{5}=3 \text { and } \overline{n_{5}}=2 .
\end{aligned}
$$

Case 2: One 2-clause $\bar{x} y_{1}$ and a $3^{+}$-clause $\bar{x} y_{2} y_{3} \mathcal{C}_{1}$ with $\left|\mathcal{C}_{1}\right| \geq 0$ contain literal $\bar{x}$, i.e., $\overline{n_{1}}=1$ and $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=1$.

Case 2.1: Three $3^{+}$-clauses contains literal $x$, i.e., $n_{1}=3, n_{3}=n_{4}=n_{5}=0$.
By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+6 \delta_{5}+\overline{n_{3}} w_{3}+\overline{n_{4}} w_{3}+\overline{n_{5}} w_{4}, w_{5}+2 \delta_{5}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$. Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=1$, $t=\left(w_{5}+6 \delta_{5}+w_{3}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$, not inferior to $t^{\prime}=\left(w_{5}+6 \delta_{5}+2 w_{3}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+\left(1-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right) . t^{\prime}$ is:

$$
\begin{aligned}
& \left(w_{5}+6 \delta_{5}+w_{3}, w_{5}+2 \delta_{5}+\delta_{4}\right) \text { when } \overline{n_{5}}=0 \\
& \left(w_{5}+6 \delta_{5}+w_{3}+\delta_{4}, w_{5}+3 \delta_{5}\right) \text { when } \overline{n_{5}}=1
\end{aligned}
$$

Case 2.2: two $3^{+}$-clauses and one 2-clause $x z_{1}$ contain literal $x$, i.e., $n_{1}=2$ and $n_{3}+n_{4}+n_{5}=1$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+4 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}} w_{3}+n_{5} \delta_{5}+\overline{n_{5}} w_{4}, w_{5}+2 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+\right.$ $n_{4} w_{3}+\overline{n_{4}} \delta_{4}+n_{5} w_{4}+\overline{n_{5}} \delta_{5}$ ). Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=1$ and $n_{3}+n_{4}+n_{5}=1$, we have $t=\left(w_{5}+4 \delta_{5}+w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+w_{3}+n_{5} \delta_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$.

Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+4 \delta_{5}+w_{3}+\left(1-n_{5}\right) \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+w_{3}+n_{5} \delta_{4}+\left(1-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$, which is

$$
\begin{aligned}
& \left(w_{5}+4 \delta_{5}+w_{3}+\delta_{4}, w_{5}+\delta_{5}+w_{3}+2 \delta_{4}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+4 \delta_{5}+w_{3}+2 \delta_{4}, w_{5}+2 \delta_{5}+w_{3}+\delta_{4}+\delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+5 \delta_{5}+w_{3}, w_{5}+\delta_{5}+w_{3}+3 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=0,
\end{aligned}
$$

$$
\left(w_{5}+5 \delta_{5}+w_{3}+\delta_{4}, w_{5}+2 \delta_{5}+w_{3}+2 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=1
$$

Case 2.3: one $3^{+}$-clauses and two 2 -clauses $x z_{1}, x z_{2}$ contain literal $x$, i.e., $n_{1}=1$ and $n_{3}+n_{4}+n_{5}=2$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+2 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}} w_{3}+n_{5} \delta_{5}+\overline{n_{5}} w_{4}, w_{5}+2 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+\right.$ $\left.n_{4} w_{3}+\overline{n_{4}} \delta_{4}+n_{5} w_{4}+\overline{n_{5}} \delta_{5}\right)$. Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=1$ and $n_{3}+n_{4}+n_{5}=2$, we have $t=\left(w_{5}+2 \delta_{5}+w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+w_{3}+n_{5} \delta_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$.

Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+2 \delta_{5}+w_{3}+\left(2-n_{5}\right) \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+2 w_{3}+n_{5} \delta_{4}+\left(1-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$, which is

$$
\begin{aligned}
& \left(w_{5}+2 \delta_{5}+w_{3}+2 \delta_{4}, w_{5}+2 \delta_{5}+2 w_{3}+\delta_{4}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+2 \delta_{5}+w_{3}+3 \delta_{4}, w_{5}+2 \delta_{5}+2 w_{3}+\delta_{5}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+3 \delta_{5}+w_{3}+\delta_{4}, w_{5}+2 \delta_{5}++2 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+3 \delta_{5}+w_{3}+2 \delta_{4}, w_{5}+3 \delta_{5}+2 w_{3}+\delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+4 \delta_{5}+w_{3}, w_{5}+2 \delta_{5}+2 w_{3}+3 \delta_{4}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+4 \delta_{5}+w_{3}+\delta_{4}, w_{5}+3 \delta_{5}+2 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=1 .
\end{aligned}
$$

Case 2.4: Three 2-clauses $x z_{1}, x z_{2}$ and $x z_{3}$ contain literal $x$, i.e., $n_{1}=0$ and $n_{3}+n_{4}+n_{5}=3$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} \delta_{4}+\overline{n_{4}} w_{3}+n_{5} \delta_{5}+\overline{n_{5}} w_{4}, w_{5}+2 \delta_{5}+\left(n_{3}+\overline{n_{3}}\right) w_{3}+n_{4} w_{3}+\right.$ $\overline{n_{4}} \delta_{4}+n_{5} w_{4}+\overline{n_{5}} \delta_{5}$ ). Since $\overline{n_{3}}+\overline{n_{4}}+\overline{n_{5}}=1$ and $n_{3}+n_{4}+n_{5}=3$, we have $t=\left(w_{5}+w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+3 w_{3}+n_{5} \delta_{4}+\overline{n_{3}} w_{3}+\overline{n_{4}} \delta_{4}+\overline{n_{5}} \delta_{5}\right)$. Note that $t$ is not inferior to
$t=\left(w_{5}+w_{3}+\left(3-n_{5}\right) \delta_{4}+n_{5} \delta_{5}+\overline{n_{5}} \delta_{4}, w_{5}+2 \delta_{5}+3 w_{3}+n_{5} \delta_{4}+\left(1-\overline{n_{5}}\right) \delta_{4}+\overline{n_{5}} \delta_{5}\right)$, which is

$$
\left(w_{5}+w_{3}+3 \delta_{4}, w_{5}+2 \delta_{5}+3 w_{3}+\delta_{4}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=0
$$

$$
\begin{aligned}
& \left(w_{5}+w_{3}+4 \delta_{4}, w_{5}+3 \delta_{5}+3 w_{3}\right) \text { when } n_{5}=0 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+w_{3}+2 \delta_{4}+\delta_{5}, w_{5}+2 \delta_{5}+3 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+w_{3}+3 \delta_{4}+\delta_{5}, w_{5}+3 \delta_{5}+3 w_{3}+\delta_{4}\right) \text { when } n_{5}=1 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+w_{3}+\delta_{4}+2 \delta_{5}, w_{5}+2 \delta_{5}+3 w_{3}+3 \delta_{4}\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+w_{3}+2 \delta_{4}+2 \delta_{5}, w_{5}+3 \delta_{5}+3 w_{3}+2 \delta_{4}+\right) \text { when } n_{5}=2 \text { and } \overline{n_{5}}=1, \\
& \left(w_{5}+w_{3}+3 \delta_{5}, w_{5}+2 \delta_{5}+3 w_{3}+4 \delta_{4}\right) \text { when } n_{5}=3 \text { and } \overline{n_{5}}=0, \\
& \left(w_{5}+w_{3}+3 \delta_{5}+\delta_{4}, w_{5}+3 \delta_{5}+3 w_{3}+3 \delta_{4}\right) \text { when } n_{5}=3 \text { and } \overline{n_{5}}=1 .
\end{aligned}
$$

Case 3: Two $3^{+}$-clauses $\bar{x} \mathcal{C}_{1}$ and $\bar{x} \mathcal{C}_{2}$ contain literal $\bar{x}$, i.e., $\overline{n_{1}}=2$ and $\overline{n_{3}}=\overline{n_{4}}=$ $\overline{n_{5}}=0$.

Case 3.1: Three $3^{+}$-clauses contains literal $x$, i.e., $n_{1}=3, n_{3}=n_{4}=n_{5}=0$.
By lemma E.8, the branching vector is not inferior to

$$
\left(w_{5}+6 \delta_{5}, w_{5}+4 \delta_{5}\right)
$$

Case 3.2: two $3^{+}$-clauses and one 2 -clause $x z_{1}$ contain literal $x$, i.e., $n_{1}=2$ and $n_{3}+n_{4}+n_{5}=1$.

By lemma E.8, the branching vector is not inferior to $t=\left(w_{5}+4 \delta_{5}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+n_{3} w_{3}+n_{4} w_{3}+n_{5} w_{4}\right)$. Since $n_{3}+n_{4}+n_{5}=1$, $t=\left(w_{5}+4 \delta_{5}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+w_{3}+n_{5} \delta_{4}\right)$. Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+4 \delta_{5}+\left(1-n_{5}\right) \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+w_{3}+n_{5} \delta_{4}\right)$, which is

$$
\begin{aligned}
& \left(w_{5}+5 \delta_{5}+\delta_{4}, w_{5}+4 \delta_{5}+w_{3}\right) \text { when } n_{5}=0 \\
& \left(w_{5}+5 \delta_{5}, w_{5}+\delta_{5}+w_{3}+\delta_{4}\right) \text { when } n_{5}=1
\end{aligned}
$$

Case 3.3: one $3^{+}$-clauses and two 2-clauses $x z_{1}, x z_{2}$ contain literal $x$, i.e., $n_{1}=1$ and $n_{3}+n_{4}+n_{5}=2$.

By lemma E.8, the branching vector is not inferior to $t=\left(w_{5}+2 \delta_{5}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+2 \delta_{5}+n_{3} w_{3}+n_{4} w_{3}+n_{5} w_{4}\right)$. Since $n_{3}+n_{4}+n_{5}=2$,
$t=\left(w_{5}+2 \delta_{5}+2 w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+n_{5} \delta_{4}\right)$. Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+2 \delta_{5}+w_{3}+\left(2-n_{5}\right) \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+n_{5} \delta_{4}\right)$, which is

$$
\begin{aligned}
& \left(w_{5}+2 \delta_{5}+w_{3}+2 \delta_{4}, w_{5}+4 \delta_{5}+2 w_{3}\right) \text { when } n_{5}=0, \\
& \left(w_{5}+3 \delta_{5}+w_{3}+\delta_{4}, w_{5}+4 \delta_{5}+2 w_{3}+\delta_{4}\right) \text { when } n_{5}=1, \\
& \left(w_{5}+4 \delta_{5}+w_{3}, w_{5}+4 \delta_{5}+2 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=2 .
\end{aligned}
$$

Case 3.4: Three 2-clauses $x z_{1}, x z_{2}$ and $x z_{3}$ contain literal $x$, i.e., $n_{1}=0$ and $n_{3}+n_{4}+n_{5}=3$.

By lemma E.8, the branching vector is not inferior to
$t=\left(w_{5}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+n_{3} w_{3}+n_{4} w_{3}+n_{5} w_{4}\right)$. Since $n_{3}+n_{4}+n_{5}=3$, $t=\left(w_{5}+2 w_{3}+n_{3} w_{3}+n_{4} \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+n_{5} \delta_{4}\right)$. Note that $t$ is not inferior to
$t^{\prime}=\left(w_{5}+w_{3}+\left(3-n_{5}\right) \delta_{4}+n_{5} \delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+n_{5} \delta_{4}\right)$, which is

$$
\begin{aligned}
& \left(w_{5}+w_{3}+3 \delta_{4}, w_{5}+4 \delta_{5}+2 w_{3}\right) \text { when } n_{5}=0, \\
& \left(w_{5}+w_{3}+2 \delta_{4}+\delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+\delta_{4}\right) \text { when } n_{5}=1, \\
& \left(w_{5}+w_{3}+\delta_{4}+2 \delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+2 \delta_{4}\right) \text { when } n_{5}=2, \\
& \left(w_{5}+w_{3}+3 \delta_{5}, w_{5}+4 \delta_{5}+2 w_{3}+3 \delta_{4}\right) \text { when } n_{5}=3 .
\end{aligned}
$$

## 3. Analysis for Formulas of Degree Larger Than 5

Let $x$ be a $\left(d_{1}, d_{0}\right)$-literal. Then $d=d_{1}+d_{0} \geq 6$. Suppose there are $s_{1} 2$-clauses containing literal $x$, and $s_{0} 2$-clauses containing literal $\bar{x}$. Let $l_{1}$ be the $l$-value reduction and $c_{1}$ be the value of $c$ from $\mathcal{F}$ to $\left.\mathcal{F}[x]\right)$. Let $l_{0}$ be the $l$-value reduction and $c_{0}$ be the $c$-value from $\mathcal{F}$ to $\mathcal{F}[\bar{x}]$. By lemma E.5, $c_{1}\left(c_{0}\right)$ is not larger than $l_{1}\left(l_{0}\right)$. Thus $l_{1}+l_{2} \geq c_{1}+c_{2}$. It can be verified that each 2-clause (containing either $x$ or $\bar{x}$ ) adds at least 0.5 to both $c_{1}$ and $c_{0}$ by the calculation of $c$-value. Thus the $s_{1}+s_{2} 2$-clauses add at least $s_{1}+s_{0}$ to $c_{1}+c_{2}$. Moreover, the $d_{1}-s_{1} 3^{+}$-clauses containing literal $x$ add at
least $2\left(d_{1}-s_{1}\right) \delta_{i} \geq d_{1}-s_{1}$ to $c_{1}$ since $i \geq 3$ (all variables in $\mathcal{F}$ are 3 -variables), and the $d_{0}-s_{0} 3^{+}$-clauses containing literal $\bar{x}$ add at least $d_{0}-s_{0}$ to $c_{0}$. Thus the $3^{+}$-clauses add at least $\left(d_{1}-s_{1}\right)+\left(d_{0}-s_{0}\right)$ to $c_{1}+c_{2}$. Finally, $x$ adds $w_{d}=0.5 d \geq 3$ to both $c_{1}$ and $c_{0}$ since $x$ is a $d$-variable where $d \geq 6$. Thus $x$ add at least $2 w_{d} \geq 6$ to $c_{1}+c_{0}$. Therefore, we have that $l_{1}+l_{0} \geq c_{1}+c_{0} \geq\left(s_{1}+s_{0}\right)+\left(d_{1}-s_{1}\right)+\left(d_{0}-s_{0}\right)+6 \geq 12$.

Next we prove that both of $l_{1}$ and $l_{0}$ is greater than $0.5+d / 2=3.5$. As shown above, the $s_{1}+s_{0} 2$-clauses add at least $0.5\left(s_{1}+s_{0}\right)$ to both $c_{1}$ and $c_{0}$, the $\left(d_{1}-s_{1}\right)$ $3^{+}$-clauses add at least $d_{1}-s_{1}$ to $c_{1}$, the $\left(d_{0}-s_{0}\right) 3^{+}$-clauses add at least $d_{0}-s_{0}$ to $c_{0}$, and $x$ add at least 3 to both $c_{1}$ and $c_{0}$. Thus $c_{1} \geq 0.5\left(s_{1}+s_{0}\right)+\left(d_{1}-s_{1}\right)+3$ and $c_{0} \geq 0.5\left(s_{1}+s_{0}\right)+\left(d_{0}-s_{0}\right)+3$. Note that $d_{0}-s_{1} \geq 0$ and $d_{0}-s_{0} \geq 0$. If $s_{1}+s_{0}=1$, then both $c_{1}$ and $c_{0}$ are not less than 3.5. If $s_{1}+s_{0}=0$, then $s_{1}=s_{0}=0$. Since both $d_{1}$ and $d_{0}$ are not less than 1 , we have that both $c_{1}$ and $c_{0}$ are not less than 4 . By lemma E.5, both $l_{1}$ and $l_{0}$ are not less than 3.5.

So the branching vector in this case is at least $\left(l_{1}, l_{2}\right)$, not inferior to $(3.5,12-$ $3.5)=(3.5,8.5)$, which leads to $O\left(1.1313^{l(\mathcal{F})}\right)$.

## 4. Branching Vector for the Main Algorithm

Summarizing all the above discussion, we can verify that the worst case occurs in the Subcase 1.1.3 in subsection E.1. The branching vector for this worst case is $t_{0}=\left(3 w_{3}+5 \delta_{4}, 3 w_{3}+\delta_{4}\right)=(7.491,3.891)$, which is inferior to the branching vectors for all other cases. The root of the polynomial corresponding to this worst branching vector is $\tau\left(t_{0}\right) \leq 1.1346$. In conclusion, we derive that the time complexity of the algorithm SATSolver is bounded by $O\left(1.1346^{l(\mathcal{F})}\right)$ on an input formula $\mathcal{F}$. Let $L$ be the total length of the formula $\mathcal{F}$, and observe that $l(\mathcal{F}) \leq L / 2$, we finally conclude the running time of the algorithm SATSolver bounded by $O\left(1.0652^{L}\right)$, which completes the proof of Theorem D.1.

## F. Final Remarks

Our main algorithm is very simple. Though our algorithm has detailed analysis of numerous cases, its analysis is quite straightforward and simple. The previous algorithm by Wahlstöm is also simple. But its analysis is quite complicated. It is interesting that both our algorithm and the algorithm by Wahlstöm deal with formulas with low degree and formulas with large degree differently.

## Algorithm Reduction $(\mathcal{F})$

INPUT: a non-empty formula $\mathcal{F}$
OUTPUT: an equivalent formula on which no further reduction is applicable
change $=$ true;
while change do
Case 1. a clause $C$ is a subset of a clause $D$ : remove $D$;
Case 2. a clause $C$ contains both $x$ and $\bar{x}$ : remove $C$;
Case 3. a clause $C$ contains multiple copies of a literal $z$ :
remove all but one $z$ in $C$;
Case 4. there is a variable $x$ with at most one non-trivial resolvent: $\mathcal{F} \leftarrow D P_{x}(\mathcal{F}) ;$
Case 5. there is a 1-clause $(z)$ or a monotone literal $z: \mathcal{F} \leftarrow \mathcal{F}[z]$;
Case 6. there exist a 2-clause $z_{1} z_{2}$ and a clause $z_{1} \overline{z_{2}} C$ :
remove $\overline{z_{2}}$ from the clause $z_{1} \overline{z_{2}} C$;
Case 7. there are clauses $z_{1} z_{2} C_{1}$ and $z_{1} \overline{z_{2}} C_{2}$ and $z_{2}$ is a (2,1)-literal: remove $z_{1}$ from the clause $z_{1} z_{2} C_{1}$;
Case 8. there are clauses $z_{1} z_{2}$ and $\overline{z_{1} z_{2}} C$ such that literal $\overline{z_{1}}$ is a 1-literal: remove the clause $z_{1} z_{2}$;
Case 9. there is a (2,2)-variable $x$ with clauses $\bar{x} z_{1}, \bar{x} z_{2}$ and two 3 -clauses $x C_{1}$ and $x C_{2}$ such that both $\overline{z_{1}}$ and $\overline{z_{2}}$ are 4 -variables in either $C_{1}$ or $C_{2}: \mathcal{F} \leftarrow D P_{x}(\mathcal{F})$. Apply Case 2, if possible.
Case 10. there is a 2 -clause $z_{1} z_{2}$ where $z_{1}$ is a 1 -literal, or there are two

2-clauses $z_{1} z_{2}$ and $\overline{z_{1} z_{2}}$ :
replace $z_{1}$ with $\overline{z_{2}}$. Apply Case 2, if possible;
Case 11. there are two clauses $C D_{1}$ and $C D_{2}$ with $|C|>1$ :
replace $C D_{1}$ and $C D_{2}$ with $\bar{x} C, x D_{1}$, and $x D_{2}$, where $x$ is a new variable. Apply Case 2 on variables in $C$, if possible.
default: change $=$ false;

Fig. 12. The reduction algorithm

## Algorithm SATSolver $(\mathcal{F})$

InPUT: a CNF formula $\mathcal{F}$
output: a report whether $\mathcal{F}$ is satisfiable

1. Reduction $(\mathcal{F})$;
2. pick a $d(\mathcal{F})$-variable $x$;
3. if $d(\mathcal{F})>5$ then
return SATSolver $(\mathcal{F}[x]) \vee$ SATSolver $(\mathcal{F}[\bar{x}])$;
4. else if $d(\mathcal{F})>3$ then
$4.1 \quad$ if $x$ is a $(2,2)$-variable with clauses $x \bar{y}_{1} z_{1}, x z_{2} z_{3}, \bar{x} y_{1}$, and $\bar{x} y_{2}$ such that $y_{1}$ is a 4 -variable and $y_{2}$ is a 3 -variable then
let $\bar{y}_{2} C_{0}$ be the clause containing $\bar{y}_{2}$;
return $\operatorname{SATSolver}\left(\mathcal{F}\left[C_{0}=\operatorname{true}\right]\right) \vee \operatorname{SATSolver}\left(\mathcal{F}\left[C_{0}=\right.\right.$ false $\left.]\right)$;
4.2 if both $x$ and $\bar{x}$ are $2^{+}$-literals then
return SATSolver $(\mathcal{F}[x]) \vee$ SATSolver $(\mathcal{F}[\bar{x}])$;
4.3 else (* assume that $\bar{x}$ occurs in a single clause $\left.\left(\bar{x} \vee z_{1} \vee \cdots \vee z_{h}\right)^{*}\right)$
return $\operatorname{SATSolver}(\mathcal{F}[x]) \vee \operatorname{SATSolver}\left(\mathcal{F}\left[\bar{x}, \bar{z}_{1}, \ldots, \bar{z}_{h}\right]\right)$;
5. else if $d(\mathcal{F})=3$ then

Apply the algorithm by Wahlström [81];
6. else return true;

Fig. 13. Algorithm for the satisfiability problem

## CHAPTER VIII

## SUMMARY AND FUTURE RESEARCH

## A. Dissertation Summary

In this dissertation, we study a new approach - measure driven algorithm design and analysis. In Chapters II to VI, we present improved fpt-algorithms for several NPhard problems. In Chapter VII, we present an improved exact algorithm for the well-known satisfiability problem. For those problems, we pick multiple measures and then find properties which help to design better algorithms. The traditional choice of single measure for a problem often restricts the application of structural properties of that problem. The discussions in previous chapters reveal that proper choice of measures do allow more structural properties to be applied to design better algorithms.

In Chapter II, we consider two measures for the 3-D matching and 3-SET PACKING problems: the number of colors needed in the coloring step and the number of elements in the dynamic programming step. With the choice of these measures, we find it is better to search for a matching (packing) $M^{\prime}$ of size $i+1$ when a matching (packing) $M$ of size $i$ is given, according to the following property-every tuple in $M$ contains at least two symbols in $M^{\prime}$. This property reduces both the number of colors needed in the coloring step and the number of elements in the dynamic programming step, thus resulting in an improved deterministic algorithm of running time $O^{*}\left(4.61^{3 k}\right)$ for both the 3-D matching and 3 -Set packing problems. Moreover, the 3-D matching problem can be solved with elements in two columns of $M^{\prime}$ in the dynamic programming step. By taking this advantage of the 3-D matching problem, we can further reduce the number of colors needed in the coloring step and
the number of elements in the dynamic programming step. The further reduction results in an improved deterministic algorithm of running time $O^{*}\left(2.32^{3 k}\right)$ for the 3-D MATCHING problem.

In Chapter III, we study the multiway cut problem. Besides considering the size of the multiway cut to search for as a measure, we also take as a measure the minimum cut from a terminal to the other terminals. These two measures are effective because of three properties: (1) there is a vertex which either decreases the size of the multiway cut to search for or increase the minimum cut from a terminal to the other terminals, (2) the multiway cut problem can be answered negatively when the minimum cut from a terminal to the other terminals is larger than the size of the multiway cut to search for, and (3) the mULTiWAY CUT problem can be solved in polynomial time when the size of the multiway cut to search for is zero. These properties and measures lead to an faster algorithm for the MULTIWAY CUT problem.

In Chapter IV, the FEEDBACK VERTEX SET problem on undirected graphs is studied. We apply the iterative compression approach to this problem. It is more effective to consider as measures both the size of the feedback vertex set (fvs) to search for and the number of connected components ( $n c c$ ) in another feedback vertex set which does not intersect with fvs. It turns out there is a vertex which either decrease fvs or ncc. Moreover, the feedback vertex set problem can be solved when $n c c=1$ or $|f v s|=0$. With these properties and measures, we design an fpt-algorithm of running time $O^{*}\left(5^{k}\right)$ for the FEEDBACK VERTEX SET problem on undirected graphs. Further investigations result in an fpt-algorithms of running time $O^{*}\left(5^{k}\right)$ for the weighted FEEDBACK VERTEX SET problem on undirected graphs.

In Chapter V, we study the FEEDBACK VERTEX SET problem on directed graphs, which had been an important open problem for 16 years before our algorithm. Similarly as to the FEEDBACK VERTEX SET problem on undirected graphs, we still apply
the iterative compression approach to the FEEDBACK VERTEX SET problem on directed graphs. However, the concept of connected components makes no sense in directed graphs. Thus the measures for the FEEDBACK VERTEX SET problem on undirected graphs can not be used for the FEEDBACK VERTEX SET problem on directed graphs. However, the feedback vertex set problem on directed graphs can be transformed into $O(k!)$ Skew separator problems. For the skew separaTOR problem, we consider two measures: the size of the skew separator cut to search for and the minimum cut from the last source to all the sinks. Properties similar to those for the MULTIWAY CUT problem can be proved for the SKEW SEPARATOR problem. An fpt-algorithm of running time $O^{*}\left(4^{k}\right)$ follows from these properties and measures for the SKEW SEPARATOR. Then an fpt-algorithm of running time $O^{*}\left(k!4^{k}\right)$ for the FEEDBACK VERTEX SET problem on digraphs follows from the algorithm for the SKEW SEPARATOR problem and the transformation from the FEEDBACK VERTEX SET problem to the $O(k!)$ SKEW SEPARATOR problems.

In Chapter VI, the max leaf problem on directed graphs is studied. We focus on a special case of this problem - the root $r$ of an out-branching $\mathcal{T}$ with at least $k$ vertices is already given. For this special case, we try to extend an out-tree $\mathcal{T}^{\prime}$ rooted at $r$ to $\mathcal{T}$. During the processing of extending $\mathcal{T}^{\prime}$ to $\mathcal{T}$, there are two important measures: (1) a subset $L_{1}$ of leaves in $\mathcal{T}^{\prime}$ such that every leaf in $L_{1}$ can reach exactly one leaf in $\mathcal{T}$, and (2) the subset $L_{2}$ of the remaining leaf in $\mathcal{T}^{\prime}$. A leaf in $\mathcal{T}^{\prime}$ either belong to $L_{1}$ or $L_{2}$. Moreover, a leaf of $\mathcal{T}^{\prime}$ with only one out-neighbour in $G-\mathcal{T}^{\prime}$ must be in $L_{1}$. Then a leaf either is in $L_{1}$ or extends $\mathcal{T}^{\prime}$ to a new out-tree which has one more leaf than $\mathcal{T}^{\prime}$. The max leaf problem can be solved when $\left|L_{1}\right|=k$ or $\left|L_{1}\right|+\left|L_{2}\right| \geq k$. With these measures and properties, we design an improved fptalgorithm of running time $O^{*}\left(4^{k}\right)$ for the Max leaf problem on directed graphs. This algorithm can also be easily applied to solve the MAX LEAF problem on undirected
graphs.
In Chapter VII, we study the well know satisfiability problem by considering its time complexity related to the total length $L$ of input formulas. Most previous algorithms study this problem by analyzing how the length $L$ changes for various cases. The previous best algorithm considers a complicated function of $L$ and other factors. The algorithm is fast. But its analysis is daunting because of the complicated function. Our algorithm takes a simple function of the number $n_{i}$ of variables of degree $i$ and the weights $w_{i}$ associated with variables of degree $i$. Our function does not depend on $L$ directly. Instead, it is bounded by $L / 2$ by careful choices of $w_{i}$ 's. With this function, we design reduction rules which make our main algorithm for the SATISFIABILITY problem quite simple.

## B. Future Work

There are several interesting questions related to the new approach and the problems studied in this dissertation: further study of this new approach, randomized and algebraic algorithms, and kernelization.

## 1. Further Study of the New Approach

Normally, it is difficult to design fpt-algorithms of upper bound better than $O^{*}\left(4^{k}\right)$ with our new approach. For example, we present $O^{*}\left(4^{k}\right)$ algorithms for the multiway CUT problem, the FEEDBACK VERTEX SET problem, the SKEW SEPARATOR problem, and the max leaf problem in this dissertation. These problems are studied with the consideration of two measures $m_{1}$ and $m_{2}$. The changes of $m_{1}$ and $m_{2}$ for these problems are not identical. For example, either $m_{1}$ or $m_{2}$ increases by 1 for the MULTIWAY CUT problem, while either $m_{1}$ increases by 1 and $m_{2}$ decreases by 1 or
$m_{2}$ increases by 1 for the max leaf problem. However, analyses show that the algorithms for both problems are of running time $O^{*}\left(4^{k}\right)$.

To have better algorithms, we can study our new approach in two directions: (1) find new structural properties which ensure larger measure changes, or (2) design better measures for known properties. For example, if some property requires that $m_{1}$ increase by at least 2 , we can have much faster algorithms for the MULTIWAY CUT problem and the max leaf problem. On the other hand, new measures may also improve analyses and result in better algorithms. Overall, we should consider both measures and properties together. This consideration allows more opportunities for better algorithms.

## 2. Randomized and Algebraic Algorithms

There are faster randomized algorithms than our algorithms for the FEEDBACK VERTEX SET problem on undirected graphs [6], the 3-D matching and 3-SET PaCKing problems [72]. However, there are no randomized algorithms faster than our algorithms for the MULTIWAY CUT problem, the FEEDBACK VERTEX SET problem on directed graphs, and the mAx LEAF problem. Simple randomization of our deterministic algorithms do not lead to faster randomized algorithms for the MULTIWAY CUT problem, the FEEDBACK VERTEX SET problem on directed graphs, and the MAX LEAF problem. Can we speed up our deterministic algorithms by randomizing our algorithms for those problems? A natural and interesting question is what properties of a problem result in randomized algorithms faster than deterministic algorithms.

The best algorithms for the 3-D matching and 3-SET PACKING problems are randomized algebraic algorithms [72]. Can these randomized algebraic algorithms be derandomized? Can we have algebraic algorithms faster than our deterministic algorithms for the MULTIWAY CUT problem, or for the FEEDBACK VERTEX SET problems?

## 3. Kernelization

Fpt-algorithms are of running time $O\left(f(k) n^{O(1)}\right)$, which can be very practical when both $n$ and $k$ are small. There are applications whose parameter $k$ is small. But their input size $n$ can be very large. It is of practical interests to reduce their input size $n$ significantly. Formally, given any instance $I$ of a problem, we in polynomial time reduced the the instance to another instance $I^{\prime}$ such that (1) $I$ has a solution if and only if $I^{\prime}$ has a solution, (2) we can construct a solution to $I$ from a solution to $I^{\prime}$ in polynomial time, and (3) If the size of $I^{\prime}$ is larger than some function $g(k)$, then we can find a solution for $I^{\prime}$. Such reduction is called kernelization and we say the problem has a $g(k)$ kernel. Theoretically, a parameterized problem has a kernel if and only if it is fixed-parameter tractable [39].

Currently, the 3-D matching and the 3-set packing problems have $O\left(k^{3}\right)$ kernels [44], and the FEEDBACK VERTEX SET problem on undirected graphs has an $O\left(k^{2}\right)$ kernel [98]. It is challenging to have better kernels for these problems. Moreover, can we find a polynomial kernel for the multiway cut problem or for the feedback VERTEX SET problem on directed graphs?

While the max leaf problem has no polynomial kernel, it has polynomial number of polynomial kernels [51]. That is, any instance of this problem can be reduced to $O\left(n^{O(1)}\right)$ number of smaller instances whose size is bounded by $O\left(k^{O(1)}\right)$. Can we find many polynomial kernels for the MULTIWAY CUT problem or for the FEEDBACK VERTEX SET problem on directed graphs?

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[^0]:    ${ }^{1}$ Formally, $O^{*}(f)$ refers to $O\left(f n^{O(1)}\right)$ where $n$ is the input size.

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[^2]:    ${ }^{2}$ To simplify the expression, we suppose that "No" plus any vertex set gives a "No". Therefore, step 2 will return a "No" if $\mathbf{N M C}\left(G-w,\left\{T_{1}, \ldots, T_{l}\right\}, k-1\right)$ returns a "No".

[^3]:    *Reprinted with permission from "Improved algorithms for feedback vertex set problems", by J. Chen, F. V. Fomin, Y. Liu, S. Liu, and Y. Villanger, 2008, Journal of Computer and System Sciences, volume 74, pages 1188-1198, Copyright [2008] by Elsevier Inc.

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[^5]:    ${ }^{1}$ by lemma C.2, $\left\{\pi_{y}\left(l_{1}\right), \ldots, \pi_{y}\left(l_{p}\right)\right\}$ is a consistent set of out-chains. The requirement on a consistent set of out-chains can be verified to be true for instances discussed in this chapter. We ignore this in later discussions.

