# An Efficient Hybrid Algorithm for the Separable Convex Quadratic Knapsack Problem 

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#### Abstract

This article considers the problem of minimizing a convex, separable quadratic function subject to a knapsack constraint and a box constraint. An algorithm called NAPHEAP has been developed to solve this problem. The algorithm solves the Karush-Kuhn-Tucker system using a starting guess to the optimal Lagrange multiplier and updating the guess monotonically in the direction of the solution. The starting guess is computed using the variable fixing method or is supplied by the user. A key innovation in our algorithm is the implementation of a heap data structure for storing the break points of the dual function and computing the solution of the dual problem. Also, a new version of the variable fixing algorithm is developed that is convergent even when the objective Hessian is not strictly positive definite. The hybrid algorithm NAPHEAP that uses a Newton-type method (variable fixing method, secant method, or Newton's method) to bracket a root, followed by a heap-based monotone break point search, can be faster than a Newton-type method by itself, as demonstrated in the numerical experiments.


CCS Concepts: • Mathematics of computing $\rightarrow$ Mathematical analysis; Mathematical optimization; Continuous optimization; Quadratic programming
Additional Key Words and Phrases: Continuous quadratic knapsack, nonlinear programming, convex programming, quadratic programming, separable programming, heap

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## 1. INTRODUCTION

We consider the following separable convex quadratic knapsack problem:

$$
\begin{array}{cl}
\min _{\mathbf{x} \in \mathbb{R}^{n}} \quad q(\mathbf{x}):=\frac{1}{2} \mathbf{x}^{\top} \mathbf{D} \mathbf{x}-\mathbf{y}^{\top} \mathbf{x}  \tag{1}\\
\text { subject to } \quad \ell \leq \mathbf{x} \leq \mathbf{u} \text { and } r \leq \mathbf{a}^{\top} \mathbf{x} \leq s,
\end{array}
$$

[^0]where $\mathbf{a}, \mathbf{y} \in \mathbb{R}^{n}, \ell \in(\mathbb{R} \cup\{-\infty\})^{n}, \mathbf{u} \in(\mathbb{R} \cup\{\infty\})^{n}, r, s \in \mathbb{R}$, and $\mathbf{D} \in \mathbb{R}^{n \times n}$ is a positive semidefinite diagonal matrix with diagonal $\mathbf{d}$. Without loss of generality, we assume that $\ell<\mathbf{u}$. Problem (1) has applications in quadratic resource allocation [Bitran and Hax 1981; Bretthauer and Shetty 1997; Cosares and Hochbaum 1994; Hochbaum and Hong 1995], quasi-Newton updates with bounds [Calamai and Moré 1987], multicapacity network flow problems [Helgason et al. 1980; Nielsen and Zenios 1992; Shetty and Muthukrishnan 1990], continuous formulations of graph partitioning problems [Hager and Hungerford 2015; Hager and Krylyuk 1999], and support vector machines [Dai and Fletcher 2006].

By introducing an auxiliary variable $b \in \mathbb{R}$, we may reformulate Problem (1) as:

$$
\begin{equation*}
\min _{\mathbf{x} \in \mathbb{R}^{n}, b \in \mathbb{R}^{2}}\left\{q(\mathbf{x}): \quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad r \leq b \leq s, \quad \mathbf{a}^{\top} \mathbf{x}=b\right\} \tag{2}
\end{equation*}
$$

Making the substitutions

$$
\mathbf{x} \leftarrow\binom{\mathbf{x}}{b}, \quad \ell \leftarrow\binom{\ell}{r}, \quad \mathbf{u} \leftarrow\binom{\mathbf{u}}{s}, \quad \mathbf{a} \leftarrow\binom{\mathbf{a}}{-1}, \quad n \leftarrow n+1,
$$

and augmenting $\mathbf{y}$ and $\mathbf{d}$ by an additional zero entry, Problem (2) is transformed into the problem

$$
\min _{\mathbf{x} \in \mathbb{R}^{n}}\left\{q(\mathbf{x}): \quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad \mathbf{a}^{\top} \mathbf{x}=0\right\}
$$

Hence, without loss of generality, we may restrict our study to an equality-constrained version of Problem (1):

$$
\begin{equation*}
\min _{\mathbf{x} \in \mathbb{R}^{n}}\left\{q(\mathbf{x}): \quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad \mathbf{a}^{\top} \mathbf{x}=b\right\} \tag{3}
\end{equation*}
$$

where $b \in \mathbb{R}$ is given. The constraints of Problem (3) are continuous analogs of the constraints that arise in the discrete knapsack problem (see Bretthauer et al. [1996]).

In quadratic resource allocation, problem (3) arises with $\mathbf{a}=\mathbf{1}$ and $\ell=\mathbf{0}$, where $\mathbf{1}$ and $\mathbf{0}$ are the vectors whose entries are all 1 and 0 , respectively. The objective is to minimize the total cost of allocating $b$ resources to $n$ projects, where $\frac{1}{2} d_{i} x_{i}^{2}-y_{i} x_{i}$ is the cost function for project $i$. The amount of resources allocated to project $i$ is constrained to lie between $\ell_{i}$ and $u_{i}$.

In other applications [Calamai and Moré 1987; Dai and Fletcher 2006; Hager and Hungerford 2015; Hager and Krylyuk 1999; Helgason et al. 1980; Nielsen and Zenios 1992; Shetty and Muthukrishnan 1990], an objective function $F$ is minimized over a feasible set described by bound constraints and a single linear equality constraint:

$$
\begin{equation*}
\min _{\mathbf{x} \in \mathbb{R}^{n}}\left\{F(\mathbf{x}): \quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad \mathbf{a}^{\top} \mathbf{x}=b\right\} \tag{4}
\end{equation*}
$$

The gradient projection algorithm as formulated in Hager and Zhang [2006] starts with an initial guess $\mathbf{x}_{0} \in \mathbb{R}^{n}$ to a solution of (4), and for each $k \geq 0$, the ( $k+1$ )st iterate is given by

$$
\mathbf{x}_{k+1}=\mathbf{x}_{k}+s_{k} \mathbf{p}_{k}, \quad \text { where } \mathbf{p}_{k}=\operatorname{proj}\left(\mathbf{x}_{k}-\alpha_{k} \nabla F\left(\mathbf{x}_{k}\right)\right) .
$$

Here $s_{k} \geq 0$ is the stepsize, $\alpha_{k} \mathbf{I}$ is an approximation to the inverse Hessian of $F$, and $\operatorname{proj}(\mathbf{z})$ is the unique projection, relative to the 2 -norm, of $\mathbf{z}$ onto the feasible set of (4). That is, $\operatorname{proj}(\mathbf{z})$ is the solution of the problem

$$
\begin{equation*}
\min _{\mathbf{x} \in \mathbb{R}^{n}}\left\{\frac{1}{2}\|\mathbf{x}-\mathbf{z}\|^{2}: \quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad \mathbf{a}^{\top} \mathbf{x}=b\right\} \tag{5}
\end{equation*}
$$

This projection problem is a special case of Problem (3) in which $\mathbf{D}$ is the identity matrix and $\mathbf{y}=\mathbf{x}_{k}-\alpha_{k} \nabla F\left(\mathbf{x}_{k}\right)$.

Specialized algorithms for solving (3) typically assume $\mathbf{D}$ is positive definite and search for a root of the derivative of the dual function, a continuous piecewise linear, monotonic function with at most $2 n$ break points ("kinks" where the slope could change). The first algorithm [Helgason et al. 1980] was based on sorting all the break points and then sequentially marching through the break points until the optimal multiplier was found. The worst case complexity of this algorithm is $O\left(n \log _{2} n\right)$ due to the sort. Starting with Brucker [1984], linear time algorithms were proposed based on a median search rather than a sort. Subsequent developments of the median search method include those in Calamai and Moré [1987], Maculan et al. [2003], and Pardalos and Kovoor [1990].

Bitran and Hax [1981] proposed a method for solving a generalization of Problem (3) in which the objective function is convex and separable but not necessarily quadratic. In their algorithm, which has come to be known as the variable fixing method, each iteration implicitly computes an estimate for the optimal Lagrange multiplier by solving a subproblem in which the box constraints are ignored. Based on the sign of the derivative of the dual function at the multiplier estimate, a non-empty subset of indices is identified for which $x_{i}$ can be optimally fixed at an upper or lower bound. The fixed variables are removed from the problem and the process repeats until all the bound components of an optimal solution have been determined. Subsequent developments of this approach in the context of Problem (3) can be found in Bretthauer et al. [1996], Kiwiel [2008], Michelot [1986], Robinson et al. [1992], Shor [1985], and Ventura [1991]. An efficient and reliable implementation of the variable fixing method for Problem (3), and a thorough convergence analysis, is given by Kiwiel [2008].

Dai and Fletcher [2006] developed a method in which each multiplier estimate is the root of a secant approximation to the derivative of the dual function (with additional modifications for speeding up convergence). In Cominetti et al. [2014], a method is proposed by Cominetti, Mascarenhas, and Silva that uses semi-smooth Newton steps for updating multiplier estimates, with a secant safeguard. Numerical experiments using a standard test set of randomly generated problems showed that the semi-smooth Newton method was faster than the variable fixing method, the secant method, and a median-based method. As we explain in Section 3, Newton's method is not very well suited for problems where the elements of $\mathbf{d}$ are small relative to elements of $\mathbf{y}$ since the derivative of the dual function is nearly piecewise constant, and when a Newton iterate lands on a flat segment, it jumps towards $\pm \infty$. As shown in Cominetti et al. [2014], if $\mathbf{a}>\mathbf{0}$ and $\mathbf{u}=\infty$ or $\ell=-\infty$, then the variable fixing method and Newton's method can generate identical iterates. This result is generalized in Section 2.

In this article we develop an algorithm called NAPHEAP that is built around a monotone break point search implemented using a heap data structure. Given an interval ( $\alpha, \beta$ ) containing an optimal dual multiplier $\lambda^{*}$ associated with the knapsack constraint, the break points on $(\alpha, \beta)$ are arranged in a heap. If there are $m$ break points, then building the heap requires about $m$ comparisons, while updating the heap after removing or adding a break point takes about $\log _{2} m$ comparisons. Hence, if either $\alpha$ or $\beta$ is separated from $\lambda^{*}$ by $l$ break points, then $\lambda^{*}$ can be found using about $m+l \log _{2} m$ comparisons. If $m$ is small, then the heap-based algorithm is fast. One situation where a good starting guess is often available is when the gradient projection algorithm is applied to the nonlinear optimization problem (4). In this situation, we compute the projection $\operatorname{proj}\left(\mathbf{x}_{k}-\alpha_{k} \nabla F\left(\mathbf{x}_{k}\right)\right)$ to obtain the search direction $\mathbf{p}_{k}$. If $\lambda_{k}$ is the Lagrange multiplier associated with the linear constraint in the projection problem $\operatorname{proj}\left(\mathbf{x}_{k}-\alpha_{k} \nabla F\left(\mathbf{x}_{k}\right)\right)$, then it is observed in Fu and Dai [2010] that a good guess for the Lagrange multiplier in the projection problem $\operatorname{proj}\left(\mathbf{x}_{k+1}-\alpha_{k+1} \nabla F\left(\mathbf{x}_{k+1}\right)\right)$ at iteration $k+1$ is $\alpha_{k+1} \lambda_{k} / \alpha_{k}$. In particular, they observe that this starting guess could reduce the computing time by $40 \%$ when compared to the starting guess $\lambda_{k}$ at iteration $k+1$.

Currently, algorithms that are able to take advantage of a good starting guess include the secant-based algorithm and the semi-smooth Newton method. In Cominetti et al. [2014], numerical comparisons were made between algorithms for solving sequences of projection problems arising in Support Vector Machine applications. When hot starts are allowed, the Newton method was shown to be faster than the other methods.

We use the expression "Newton-type method" to refer to the class of methods that includes the variable fixing algorithm, Newton's method, and the secant method. If the knapsack problem is not connected with a convergent algorithm like the gradient projection algorithm, then a few iterations of a Newton-type method may yield a good starting guess. Let $\mathbf{x}^{*}$ denote an optimal solution of (1), assuming it exists. When on the order of $n$ components of $\mathbf{x}^{*}$ satisfy $\ell_{i}<x_{i}^{*}<u_{i}$, the time for an iteration of a Newtontype method is proportional to $n$, even if the iterates start very close to $\lambda^{*}$. On the other hand, the time to pass over a break point in a step of NAPHEAP is proportional to $\log _{2} m$, where $m$ is the number of break points in an interval bracketing a solution. As we will show in the numerical experiments, a hybrid algorithm that uses a Newton-type method to generate a starting guess followed by heap-based search (until convergence) can be much faster than a Newton-type method by itself.

Our article is organized as follows. In Section 2, we give an overview of algorithms for solving (3) when $\mathbf{D}$ is positive definite, and we identify the elements these algorithms have in common. Section 3 studies the complexity of Newton-type methods, both in theory and in experiments. An example is constructed that shows the worst-case running time of a variable fixing algorithm could grow like $n^{2}$. Section 4 develops a generalization of the variable fixing algorithm that can be used even when some components of $\mathbf{d}$ vanish. Section 5 presents our heap-based algorithm for solving Problem (3). Finally, Section 6 compares the performance of our hybrid NAPHEAP algorithm to the Newton method from Cominetti et al. [2014] using randomly generated test problems. The structure of the dual function is investigated in an effort to understand the implications of numerical results based on randomly generated problems.

Notation. 0 and 1 denote vectors whose entries are all 0 and all 1, respectively, the dimensions should be clear from context. If $\mathcal{S}$ is a set, then $|\mathcal{S}|$ denotes the number of elements in $\mathcal{S}$, and $\mathcal{S}^{c}$ is the complement of $\mathcal{S}$. A subscript $k$ is often used to denote the iteration number. Thus $\mathbf{x}_{k}$ is the $k$ th iterate and $x_{k i}$ is the $i$ th component of the $k$ th iterate.

## 2. OVERVIEW OF ALGORITHMS

For ease of exposition, it is assumed throughout the article that $\mathbf{a}>\mathbf{0}$. Note that if $a_{i}=0$, then $x_{i}$ does not appear in the knapsack constraint and the optimal $x_{i}$ is any solution of

$$
\min \left\{.5 x_{i}^{2} d_{i}-y_{i} x_{i}: \ell_{i} \leq x_{i} \leq u_{i}\right\},
$$

and if $a_{i}<0$, then we can make the change of variables $z_{i}=-x_{i}$ to obtain an equivalent problem with $a_{i}>0$. In this section, we also assume that $\mathbf{d}>\mathbf{0}$, while in the next section we take into account vanishing diagonal elements.
The common approach to (3) is to solve the dual problem. Let $\mathcal{L}$ denote the Lagrangian defined by

$$
\mathcal{L}(\mathbf{x}, \lambda)=\frac{1}{2} \mathbf{x}^{\top} \mathbf{D} \mathbf{x}-\mathbf{y}^{\top} \mathbf{x}+\lambda\left(\mathbf{a}^{\top} \mathbf{x}-b\right),
$$

and let $L$ denote the dual function

$$
\begin{equation*}
L(\lambda)=\min \{\mathcal{L}(\mathbf{x}, \lambda): \ell \leq \mathbf{x} \leq \mathbf{u}\} . \tag{6}
\end{equation*}
$$

The dual problem is

$$
\begin{equation*}
\max \{L(\lambda): \lambda \in \mathbb{R}\} . \tag{7}
\end{equation*}
$$

Since $\mathbf{d}>\mathbf{0}$, it follows that $L$ is differentiable (see Clarke [1975, Thm. 2.1] or Danskin [1967]), and

$$
L^{\prime}(\lambda)=\mathbf{a}^{\top} \mathbf{x}(\lambda)-b
$$

where $\mathbf{x}(\lambda)$ is the unique minimizer in (6) given by

$$
\begin{equation*}
x_{i}(\lambda)=\operatorname{mid}\left(\ell_{i},\left(y_{i}-\lambda a_{i}\right) / d_{i}, u_{i}\right), i=1,2, \ldots, n \tag{8}
\end{equation*}
$$

and $\operatorname{mid}(a, b, c)$ denotes the median (or middle) of $a, b$, and $c$. The following result is well known (for example, see Brucker [1984]):

Proposition 2.1. Suppose $\mathbf{d}>\mathbf{0}$ and (3) is feasible.
(1) $L$ is concave and $L^{\prime}$ is non-increasing, continuous, and piecewise linear with break points given by the set

$$
\begin{equation*}
\Lambda:=\bigcup_{1 \leq i \leq n}\left\{\frac{y_{i}-\ell_{i} d_{i}}{a_{i}}, \frac{y_{i}-u_{i} d_{i}}{a_{i}}\right\} \tag{9}
\end{equation*}
$$

(2) (7) has a solution $\lambda^{*} \in \mathbb{R}$ and $L^{\prime}\left(\lambda^{*}\right)=0$.
(3) If $L^{\prime}\left(\lambda^{*}\right)=0$, then $\mathbf{x}\left(\lambda^{*}\right)$ defined in (8) is the unique solution of (3).

Even in the case where $\lambda$ is not an optimal multiplier, one can still use the sign of $L^{\prime}(\lambda)$ to determine some bound components of the optimal solution to Problem (3) (for example, see Ventura [1991, Thm. 6]):

Proposition 2.2. If $\mathbf{d}>\mathbf{0}, \mathbf{a} \geq \mathbf{0}$, and $\mathbf{x}^{*}$ is a solution of (3), then for any $\lambda \in \mathbb{R}$, we have the following:
(1) If $L^{\prime}(\lambda) \geq 0$, then $x_{i}^{*}=\ell_{i}$ for every $i$ such that $x_{i}(\lambda)=\ell_{i}$.
(2) If $L^{\prime}(\lambda) \leq 0$, then $x_{i}^{*}=u_{i}$ for every $i$ such that $x_{i}(\lambda)=u_{i}$.

Proof. If $L^{\prime}(\lambda) \geq 0$, then since $L^{\prime}$ is non-increasing, it follows that $\lambda \leq \lambda^{*}$ for any $\lambda^{*}$ which satisfies $L^{\prime}\left(\lambda^{*}\right)=0$. If $x_{i}(\lambda)=\ell_{i}$, then by the definition of $\mathbf{x}(\lambda)$, we conclude that $\left(y_{i}-\lambda a_{i}\right) / d_{i} \leq \ell_{i}$. And since $a_{i}>0$ and $\lambda^{*} \geq \lambda$,

$$
\frac{y_{i}-\lambda^{*} a_{i}}{d_{i}} \leq \frac{y_{i}-\lambda a_{i}}{d_{i}} \leq \ell_{i} .
$$

This implies that

$$
x_{i}^{*}=x_{i}\left(\lambda^{*}\right)=\operatorname{mid}\left(\ell_{i},\left(y_{i}-\lambda^{*} a_{i}\right) / d_{i}, u_{i}\right)=\ell_{i} .
$$

The case where $L^{\prime}(\lambda) \leq 0$ is treated similarly.
The specialized algorithms that have been developed for Problem (3) have the generic form shown in Algorithm 1. The algorithms start with an initial guess $\lambda_{1}$ for the optimal dual multiplier and update $\lambda_{k}$ until $L^{\prime}\left(\lambda_{k}\right)=0$. The solution $\mathbf{x}^{*}$ to Problem (3) is the vector $\mathbf{x}\left(\lambda_{k}\right)$ constructed using Equation (8). In particular,

$$
x_{i}^{*}=\left\{\begin{array}{cl}
\ell_{i} & \text { if } i \in \mathcal{L}_{k}, \\
u_{i} & \text { if } i \in \mathcal{U}_{k}, \\
\operatorname{mid}\left(\ell_{i},\left(y_{i}-\lambda_{k} a_{i}\right) / d_{i}, u_{i}\right) & \text { otherwise }
\end{array}\right.
$$

In Step 3, Proposition 2.2 may be used to "fix" the values of some components of $\mathbf{x}^{*}$ at each iteration.

```
ALGORITHM 1: Solve a Separable Quadratic Knapsack Problem (3) with \(\mathbf{d}>\mathbf{0}\) and \(\mathbf{a}>\mathbf{0}\).
Input: Vectors \(\mathbf{a}, \mathbf{d}, \mathbf{y}, \ell, \mathbf{u}\) and scalar \(b\) defining the problem (3); initial guess \(\lambda_{1}\)
Output: The solution vector \(\mathbf{x}\)
\(k=1 ; \mathcal{L}_{1}=\emptyset ; \mathcal{U}_{1}=\emptyset ;\)
while true do
    Step 1: Generate \(\lambda_{k}\) (use initial guess if \(k=1\) ).
    Step 2: Stopping Criterion:
    if \(L^{\prime}\left(\lambda_{k}\right)=0\) then
        \(\mathbf{x}^{*}=\mathbf{x}\left(\lambda_{k}\right)\)
        break
    end
    Step 3: Variable fixing (optional):
    Define \(\mathcal{F}_{k}=\left(\mathcal{U}_{k} \cup \mathcal{L}_{k}\right)^{c}\).
    if \(L^{\prime}\left(\lambda_{k}\right)>0\) then
                    \(\mathcal{L}_{k+1}=\mathcal{L}_{k} \cup\left\{i \in \mathcal{F}_{k}: x_{i}\left(\lambda_{k}\right)=\ell_{i}\right\}\) and \(\mathcal{U}_{k+1}=\mathcal{U}_{k}\).
    end
    if \(L^{\prime}\left(\lambda_{k}\right)<0\) then
                    \(\mathcal{U}_{k+1}=\mathcal{U}_{k} \cup\left\{i \in \mathcal{F}_{k}: x_{i}\left(\lambda_{k}\right)=u_{i}\right\}\) and \(\mathcal{L}_{k+1}=\mathcal{L}_{k}\).
    end
    \(k=k+1\)
end
```

Algorithms for solving Problem (3) differ primarily in how they choose $\lambda_{k}$ in Step 1. For example, in a break-point search [Helgason et al. 1980], $\lambda_{k}$ is the next break point between $\lambda_{k-1}$ and $\lambda^{*}$ if such a break point exists. Otherwise, $\lambda_{k}=\lambda^{*}$ is found by linear interpolation. In the median search methods of Brucker [1984], Calamai and Moré [1987], Maculan et al. [2003], and Pardalos and Kovoor [1990], the iteration amounts to selecting $\lambda_{k}$ as the median of the remaining break points. Based on the sign of $L^{\prime}$ at the median, half of the remaining break points can be discarded. In the secant-based algorithm of Dai and Fletcher [2006], $\lambda_{k}$ is the root of a secant based on the value of $L^{\prime}$ at two previous iterates (with additional modifications for speeding up convergence).

In the variable fixing methods [Bitran and Hax 1981; Bretthauer et al. 1996; Kiwiel 2008; Michelot 1986; Robinson et al. 1992; Shor 1985; Ventura 1991], each iteration $k \geq 2$ solves a subproblem over the remaining unfixed variables:

$$
\begin{equation*}
\min \left\{\sum_{i \in \mathcal{F}_{k}} \frac{1}{2} d_{i} x_{i}^{2}-y_{i} x_{i}: \sum_{i \in \mathcal{F}_{k}} a_{i} x_{i}=b_{k}\right\} . \tag{12}
\end{equation*}
$$

Here, $\mathcal{F}_{k}=\mathcal{B}_{k}^{c}$, where $\mathcal{B}_{k}=\mathcal{L}_{k} \cup \mathcal{U}_{k}$, and

$$
\begin{equation*}
b_{k}=b-\sum_{i \in \mathcal{B}_{k}} a_{i} x_{i}^{*} . \tag{13}
\end{equation*}
$$

Note that the constraint $\ell_{i} \leq x_{i} \leq u_{i}$ is dropped in (12). The iterate $\lambda_{k}$ is the optimal multiplier associated with the linear constraint in (12) and is given by

$$
\begin{equation*}
\lambda_{k}^{F}=\frac{-b_{k}+\sum_{i \in \mathcal{F}_{k}} a_{i} y_{i} / d_{i}}{\sum_{i \in \mathcal{F}_{k}} a_{i}^{2} / d_{i}} \tag{14}
\end{equation*}
$$

While the method may begin with an arbitrary guess $\lambda_{1} \in \mathbb{R}$, many implementations compute $\lambda_{1}$ using the formula (14) with $\mathcal{F}_{1}=\{1,2, \ldots, n\}$.

The recent method of Cominetti et al. [2014] updates the multiplier by taking a semi-smooth Newton step in the direction of $\lambda^{*}$ :

$$
\begin{equation*}
\lambda_{k}^{N}=\lambda_{k-1}-\frac{L^{\prime}\left(\lambda_{k-1}\right)}{L_{ \pm}^{\prime \prime}\left(\lambda_{k-1}\right)}, \quad k \geq 2 \tag{15}
\end{equation*}
$$

for some starting guess $\lambda_{1} \in \mathbb{R}$. Here, $L_{ \pm}^{\prime \prime}\left(\lambda_{k-1}\right)$ denotes either the right or left side derivative of $L^{\prime}$ at $\lambda_{k-1}$ (the side is chosen according to the sign of $L^{\prime}$ ). Equivalently, $\lambda_{k}^{N}$ is the maximizer of the second-order Taylor series approximation to $L$ at $\lambda_{k-1}$ in the direction of $\lambda^{*}$ (in the case where $\lambda_{k-1}$ is a break point, the Taylor series approximation depends on the direction in which we expand). If the Newton step is unacceptably large, then $\lambda_{k}$ is either computed using a secant approximation or by simply moving to the next break point in the direction of $\lambda^{*}$. These modifications also prevent the iterates from cycling.

The Newton algorithm [Cominetti et al. 2014] and the variable fixing algorithm [Kiwiel 2008] are closely related. Both algorithms employ the variable fixing operation, and the iterate $\lambda_{k}$ has the property that $\lambda_{k} \in\left(\alpha_{k-1}, \beta_{k-1}\right)$, where

$$
\begin{align*}
\alpha_{j} & =\sup \left\{\lambda_{i}: L^{\prime}\left(\lambda_{i}\right)>0 \text { and } i \leq j\right\} \quad \text { and }  \tag{16}\\
\beta_{j} & =\inf \left\{\lambda_{i}: L^{\prime}\left(\lambda_{i}\right)<0 \text { and } i \leq j\right\},
\end{align*}
$$

where we use the convention that $\sup \emptyset=\infty$ and $\inf \emptyset=-\infty$. For the variable fixing algorithm, this is a consequence of the formula (14)-see Lemma 4.1 in Kiwiel [2008]. For the Newton scheme in Cominetti et al. [2014], this property is ensured by replacing the Newton iterate by a secant iterate when the Newton iterate is not contained in ( $\alpha_{k-1}, \beta_{k-1}$ ).

Any algorithm of the form of the generic Algorithm 1 that includes the variable fixing step and that produces iterates $\lambda_{k} \in\left(\alpha_{k-1}, \beta_{k-1}\right)$, also has the property that

$$
\begin{equation*}
x_{i}\left(\lambda_{k}\right)=x_{i}^{*} \text { for all } i \in \mathcal{B}_{k+1} . \tag{17}
\end{equation*}
$$

In particular, if $i \in \mathcal{B}_{k+1} \backslash \mathcal{B}_{k}$, then by Equations (10) and (11), $x_{i}\left(\lambda_{k}\right)=x_{i}^{*}$. If $i \in \mathcal{B}_{k}$, then it entered $\mathcal{B}_{k}$ at an earlier iteration from either of the updates (10) or (11). To be specific, suppose that for some $j<k$, we have $L^{\prime}\left(\lambda_{j}\right)>0$ and $x_{i}\left(\lambda_{j}\right)=\ell_{i}=x_{i}^{*}$. Since $\alpha_{k-1}$ is expressed as a maximum in Equation (16), it follows that $\lambda_{j} \leq \alpha_{k-1}$. Since $\alpha_{i}>0$, $x_{i}(\lambda)$ is a decreasing function of $\lambda$ with $x_{i}(\lambda) \geq \ell_{i}$ for all $\lambda$. It follows that

$$
x_{i}^{*}=\ell_{i}=x_{i}\left(\lambda_{j}\right) \geq x_{i}\left(\alpha_{k-1}\right) \geq x_{i}\left(\lambda_{k}\right) \geq \ell_{i},
$$

which yields Equation (17). The case $L^{\prime}\left(\lambda_{j}\right)<0$ and $x_{i}\left(\lambda_{j}\right)=u_{i}=x_{i}^{*}$ is similar.
Another interesting connection between these two algorithms, established below, is that both $\lambda_{k}^{F}$ and $\lambda_{k}^{N}$ maximize a quadratic approximation to $L$ of the form

$$
\begin{equation*}
L_{\mathcal{B}}(\lambda)=\min \left\{\mathcal{L}(\mathbf{z}, \lambda): \mathbf{z} \in \mathbb{R}^{n}, z_{i}=x_{i}\left(\lambda_{k-1}\right) \text { for all } i \in \mathcal{B}\right\}, \tag{18}
\end{equation*}
$$

for some choice of $\mathcal{B} \subseteq \mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$, where $\mathcal{A}(\mathbf{x})$ denotes the set of active indices:

$$
\mathcal{A}(\mathbf{x})=\left\{i: x_{i}=\ell_{i} \text { or } x_{i}=u_{i}\right\} .
$$

Lemma 2.3. Assume either that $k=1$ or that $k \geq 2$ and $L^{\prime}\left(\lambda_{k-1}\right) \neq 0$. Then the variable fixing iterate $\lambda_{k}^{F}$ in (14) is the unique maximizer of $L_{\mathcal{B}_{k}}(\lambda)$.

Proof. By (17), $x_{i}\left(\lambda_{k-1}\right)=x_{i}^{*}$ for all $i \in \mathcal{B}_{k}$ (if $k=1$, then $\mathcal{B}_{k}=\emptyset$, so this is vacuously true). Hence, $L_{\mathcal{B}_{k}}$ can be expressed

$$
L_{\mathcal{B}_{k}}(\lambda)=\min \left\{\mathcal{L}(\mathbf{z}, \lambda): \mathbf{z} \in \mathbb{R}^{n}, z_{i}=x_{i}^{*} \text { for all } i \in \mathcal{B}_{k}\right\} .
$$

Consequently, $L_{\mathcal{B}_{k}}$ is the dual function associated with the optimization problem

$$
\begin{equation*}
\min \left\{\frac{1}{2} \mathbf{z}^{\top} \mathbf{D} \mathbf{z}-\mathbf{y}^{\top} \mathbf{z}: \mathbf{a}^{\top} \mathbf{z}=b, z_{i}=x_{i}^{*} \text { for all } i \in \mathcal{B}_{k}\right\} . \tag{19}
\end{equation*}
$$

Since either $k=1$ or $L^{\prime}\left(\lambda_{k-1}\right) \neq 0$, we have that $\mathcal{F}_{k}$ is nonempty; hence, the maximizer of $L_{\mathcal{B}_{k}}$ is unique since it is strongly convex. Since the optimization problem (12) is the same as (19), the maximizer of $L_{\mathcal{B}_{k}}$ is the same as the optimal multiplier associated with the linear constraint of (19), which is the same as the optimal multiplier associated with the linear constraint of (12).

Next, let us relate the Newton iterate to the maximizer of a dual function $L_{\mathcal{B}}$ for some choice of $\mathcal{B}$.

Lemma 2.4. Let $k \geq 2$ and suppose that $L^{\prime}\left(\lambda_{k-1}\right) \neq 0$. Let $\mathcal{A}_{0}$ be defined by

$$
\mathcal{A}_{0}:=\left\{i: \frac{y_{i}-\lambda_{k-1} a_{i}}{d_{i}}=\ell_{i} \text { if } L^{\prime}\left(\lambda_{k-1}\right)<0 \text { and } \frac{y_{i}-\lambda_{k-1} \alpha_{i}}{d_{i}}=u_{i} \text { if } L^{\prime}\left(\lambda_{k-1}\right)>0\right\} .
$$

Then $\lambda_{k}^{N}$ is the unique maximizer of $L_{\mathcal{B}}(\lambda)$, where $\mathcal{B}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right) \backslash \mathcal{A}_{0}$.
Proof. Assume without loss of generality that $L^{\prime}\left(\lambda_{k-1}\right)>0$. By definition, $\lambda_{k}^{N}$ is the maximizer of the second-order Taylor series

$$
\begin{equation*}
L\left(\lambda_{k-1}\right)+L^{\prime}\left(\lambda_{k-1}\right)\left(\lambda-\lambda_{k-1}\right)+\frac{1}{2} L_{+}^{\prime \prime}\left(\lambda_{k-1}\right)\left(\lambda-\lambda_{k-1}\right)^{2} . \tag{20}
\end{equation*}
$$

Hence, we need only show that for $\mathcal{B}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right) \backslash \mathcal{A}_{0}, L_{\mathcal{B}}(\lambda)$ is equivalent to the expression (20). Since $L_{\mathcal{B}}(\lambda)$ is a quadratic function of $\lambda$, we need only show that $L_{\mathcal{B}}\left(\lambda_{k-1}\right)=L\left(\lambda_{k-1}\right), L_{\mathcal{B}}^{\prime}\left(\lambda_{k-1}\right)=L^{\prime}\left(\lambda_{k-1}\right)$, and $L_{\mathcal{B}}^{\prime \prime}\left(\lambda_{k-1}\right)=L_{+}^{\prime \prime}\left(\lambda_{k-1}\right)$.

For each $\lambda$, let $\mathbf{z}(\lambda)$ denote the unique solution to Problem (18). We claim that $\mathbf{z}\left(\lambda_{k-1}\right)=$ $\mathbf{x}\left(\lambda_{k-1}\right)$. If $i \in \mathcal{B}$, then $z_{i}\left(\lambda_{k-1}\right)=x_{i}\left(\lambda_{k-1}\right)$ by the constraint in Problem (18). If $i \in \mathcal{B}^{c}$, then either $i \in \mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)^{c}$ or $i \in \mathcal{A}_{0}$. In either case, it follows from Equation (8) that

$$
x_{i}\left(\lambda_{k-1}\right)=\frac{y_{i}-\lambda_{k-1} a_{i}}{d_{i}} .
$$

By direct substitution, we obtain

$$
\frac{\partial}{\partial x_{i}} \mathcal{L}\left(\mathbf{x}\left(\lambda_{k-1}\right), \lambda_{k-1}\right)=0
$$

Hence, $\mathbf{x}\left(\lambda_{k-1}\right)$ satisfies the first-order optimality conditions for Problem (18) and by the strong convexity of the objective function $\mathbf{z}\left(\lambda_{k-1}\right)=\mathbf{x}\left(\lambda_{k-1}\right)$. It follows immediately that $L_{\mathcal{B}}\left(\lambda_{k-1}\right)=L\left(\lambda_{k-1}\right)$. Moreover, by Clarke [1975, Thm. 2.1], we have

$$
L_{\mathcal{B}}^{\prime}\left(\lambda_{k-1}\right)=\mathbf{a}^{\top} \mathbf{z}\left(\lambda_{k-1}\right)-b=\mathbf{a}^{\top} \mathbf{x}\left(\lambda_{k-1}\right)-b=L^{\prime}\left(\lambda_{k-1}\right)
$$

Next, let us consider the second derivative of $L_{\mathcal{B}}$. We claim that for each $i$,

$$
\begin{equation*}
z_{i}^{\prime}\left(\lambda_{k-1}\right)=x_{i}^{\prime}\left(\lambda_{k-1}^{+}\right) \tag{21}
\end{equation*}
$$

Indeed, if $i \in \mathcal{B}$, then $i \in \mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right) \backslash \mathcal{A}_{0}$, so by definition of $\mathcal{A}_{0}$, we have either

$$
\frac{y_{i}-\lambda_{k-1} a_{i}}{d_{i}} \leq \ell_{i} \text { or } \frac{y_{i}-\lambda_{k-1} a_{i}}{d_{i}}>u_{i} .
$$

So, $x_{i}(\lambda)=u_{i}$ or $\ell_{i}$ if $\lambda \in\left[\lambda_{k}, \lambda_{k}+\epsilon\right]$ with $\epsilon>0$ sufficiently small. Hence, in the case $i \in \mathcal{B}, z_{i}(\lambda)$ is constant on $\mathbb{R}$, and $z_{i}^{\prime}\left(\lambda_{k-1}\right)=0=x_{i}^{\prime}\left(\lambda_{k-1}^{+}\right)$. On the other hand, if $i \in \mathcal{B}^{c}$,
then either $i \in \mathcal{A}_{0}$ or $\ell_{i}<x_{i}\left(\lambda_{k-1}\right)<u_{i}$. In either case, we have

$$
\ell_{i}<x_{i}(\lambda)=\frac{y_{i}-\lambda a_{i}}{d_{i}} \leq u_{i}
$$

for $\lambda \in\left[\lambda_{k-1}, \lambda_{k-1}+\epsilon\right]$ and $\epsilon>0$. Consequently, $z_{i}^{\prime}\left(\lambda_{k-1}\right)=-a_{i} / d_{i}=x_{i}^{\prime}\left(\lambda_{k-1}^{+}\right)$. This completes the proof of the claim (21). Thus, we have

$$
L_{\mathcal{B}}^{\prime \prime}\left(\lambda_{k-1}\right)=\left[\mathbf{a}^{\top} \mathbf{z}\left(\lambda_{k-1}\right)-b\right]^{\prime}=\left[\mathbf{a}^{\top} \mathbf{x}\left(\lambda_{k-1}^{+}\right)-b\right]^{\prime}=L_{+}^{\prime \prime}\left(\lambda_{k-1}\right),
$$

which completes the proof.
The next proposition shows that the iterates $\lambda_{k}^{F}$ and $\lambda_{k}^{N}$ coincide whenever the active components of $\mathbf{x}\left(\lambda_{k-1}\right)$ coincide with $\mathcal{B}_{k}$.

Theorem 2.5. Let $k \geq 2$ and assume that $L^{\prime}\left(\lambda_{k-1}\right) \neq 0$. If $\mathcal{B}_{k}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$, then $\lambda_{k}^{N}=\lambda_{k}^{F}$.

Proof. Assume without loss of generality that $L^{\prime}\left(\lambda_{k-1}\right)>0$. We prove that $\mathcal{A}_{0}=\emptyset$ by contradiction, where $\mathcal{A}_{0}$ is defined in Lemma 2.4. If $i \in \mathcal{A}_{0}$, then $i \in \mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$ and

$$
x_{i}\left(\lambda_{k-1}\right)=\frac{y_{i}-\lambda_{k-1} a_{i}}{d_{i}}=u_{i} .
$$

Since $a_{i}>0$, it follows that $x_{i}(\lambda)<x_{i}\left(\lambda_{k-1}\right)=u_{i}$ for every $\lambda>\lambda_{k-1}$. Hence, $x_{i}^{*}$ cannot have been fixed by the end of iteration $k-1$, and $i \notin \mathcal{B}_{k}$. Since $i \in \mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$ but $i \notin \mathcal{B}_{k}$, we contradict the assumption that $\mathcal{B}_{k}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$. Therefore, $\mathcal{A}_{0}=\emptyset$, and by Lemma 2.4, $\lambda_{k}^{N}$ is the unique maximizer of $\mathcal{L}_{\mathcal{B}}$ for $\mathcal{B}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$. By Lemma 2.3, $\lambda_{k}^{F}$ maximizes $\mathcal{L}_{\mathcal{B}}$ for $\mathcal{B}=\mathcal{B}_{k}$. So since $\mathcal{B}_{k}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right), \lambda_{k}^{N}=\lambda_{k}^{F}$.

Remark 2.6. Theorem 2.5 generalizes Proposition 5.1 in Cominetti et al. [2014] in which the semi-smooth Newton iterates are shown to be equivalent to the variable fixing iterates in the case where $\lambda_{1}$ satisfies (14), $\mathbf{a}>\mathbf{0}$, and $\mathbf{u}=\infty$. In this case, it can be shown that the sequence $\lambda_{k}$ is monotonically increasing towards $\lambda^{*}, x_{i}\left(\lambda_{k}\right)$ is monotonically decreasing, and whenever a variable reaches a lower bound, it is fixed; that is, at each iteration $k$, we have $\mathcal{B}_{k}=\mathcal{A}\left(\mathbf{x}\left(\lambda_{k-1}\right)\right)$. By Theorem $2.5, \lambda_{k}^{F}=\lambda_{k}^{N}$.

## 3. WORST-CASE PERFORMANCE OF NEWTON-TYPE METHODS

In this section, we examine the worst-case performance of the Newton-type methods such as the semi-smooth Newton method, the variable fixing method, and the secant method. All of these methods require the computation of $L^{\prime}\left(\lambda_{k}\right)$ in each iteration. Since this computation involves a sum over the free indices, it follows that if $\Omega(n)$ components of an optimal solution are free, then each iteration of a Newton-type method requires $\Omega(n)$ flops. Here $\Omega(n)$ denotes a number bounded from below by $c n$ for some $c>0$.
Table I shows the performance of the variable fixing algorithm for a randomly selected example from Problem Set 1 of Section 6 of size $n=3,000,000$. For each iteration, we give the CPU time for that iteration in seconds, the value of $\lambda_{k}$, and the size of the sets $\mathcal{F}_{k}=\mathcal{B}_{k}^{c}$ and $\mathcal{B}_{k+1} \backslash \mathcal{B}_{k}$. The algorithm converges after 10 iterations. However, after only 5 iterations, the relative error between $\lambda_{k}$ and $\lambda^{*}$ is already within $0.7 \%$. The time for iteration 10 is smaller than the rest since the stopping condition was satisfied before completing the iteration. For iterations 5 through 9, the time per iteration is about 0.03 s and the number of free variables is on the order of 1 million.

In Table II we solve the same problem associated with Table I, but with the good starting guess $\lambda_{1}=3.84760$, which agrees with the exact multiplier to six significant digits. This starting guess is so good that all the components of the optimal solution that are at the upper bound can be fixed in the first iteration. Nonetheless, the variable

Table I. Statistics for the Variable Fixing Algorithm Applied to an Example from Problem Set 1

| $k$ | CPU (s) | $\lambda_{k}$ | $\left\|\mathcal{F}_{k}\right\|$ | $\left\|\mathcal{B}_{k+1} \backslash \mathcal{B}_{k}\right\|$ |
| :---: | :---: | :---: | :---: | :---: |
| 1 | .190 | 0.2560789 | 2153491 | 846509 |
| 2 | .066 | 1.5665073 | 1895702 | 257789 |
| 3 | .066 | 3.3274053 | 1608655 | 287047 |
| 4 | .054 | 4.4810929 | 1248067 | 360588 |
| 5 | .032 | 3.8210237 | 1179116 | 68951 |
| 6 | .030 | 3.8630132 | 1146762 | 32354 |
| 7 | .029 | 3.8475087 | 1143152 | 3610 |
| 8 | .029 | 3.8476129 | 1142346 | 806 |
| 9 | .028 | 3.8476022 | 1142331 | 15 |
| 10 | .014 | 3.8476022 | 1142331 | 0 |

Table II. Statistics for the Variable Fixing Algorithm Applied to the Same Example Shown in Table I But with a Very Good Starting Guess

| $k$ | CPU (s) | $\lambda_{k}$ | $\left\|\mathcal{F}_{k}\right\|$ | $\left\|\mathcal{B}_{k+1} \backslash \mathcal{B}_{k}\right\|$ |
| :---: | :---: | :---: | :---: | :---: |
| 1 | .127 | 3.8476000 | 2606250 | 393750 |
| 2 | .096 | 0.1763559 | 1775119 | 831131 |
| 3 | .051 | 1.3610706 | 1540751 | 234368 |
| 4 | .045 | 2.6162104 | 1321817 | 218934 |
| 5 | .036 | 3.5998977 | 1176249 | 145568 |
| 6 | .030 | 3.8380911 | 1143634 | 32615 |
| 7 | .029 | 3.8475885 | 1142330 | 1304 |
| 8 | .014 | 3.8476022 | 1142330 | 0 |

fixing algorithm still took eight iterations to reach the optimal solution, and the time for the trailing iterations is still around 0.03 s when the number of free variables is on the order of 1 million.
Although Newton's method, starting from the good guess, would converge in one iteration on the problem of Table I, it still requires $\Omega(n)$ flops per iteration. The good starting guess helps Newton's method by reducing the number of iterations but not the time per iteration. Newton's method may also encounter convergence problems when there are small diagonal elements. To illustrate the effect of small diagonal elements, we consider a series of knapsack problems of the following form:

$$
\begin{equation*}
d_{i}=\delta, \quad a_{i}=1, \quad y_{i}=\operatorname{rand}[-10,10], \quad \ell_{i}=0, \quad u_{i}=1 \tag{22}
\end{equation*}
$$

Here rand $[-10,10]$ denotes a random number between -10 and 10 with a uniform distribution. The series of problems depends on the parameter $\delta$. In Figure 1 we plot $L^{\prime}$ for four different values of $\delta$. When $\delta=4$, the plot is approximately linear, and Newton's method should find the root quickly. However, when $\delta$ decreases to 1 , the plot develops a flat spot, and any Newton iterate landing on this flat spot would be kicked out toward $\pm \infty$. When $\delta$ reaches 0.0625 , the graph is essentially piecewise constant, and, in this case, Newton's method would not work well.
To correct for this poor performance, the authors of Cominetti et al. [2014] implement a safeguard; whenever the Newton iterate lies outside the interval ( $\alpha_{k-1}, \beta_{k-1}$ ), the multiplier update is computed by either moving to the next break point in the direction of $\lambda^{*}$ or using a secant step. In either case, convergence may be quite slow. For example, Figure 2 shows that when the graph of $L^{\prime}$ is essentially piecewise constant, the convergence of a secant iteration based on the function values at $\alpha_{k-1}$ and $\beta_{k-1}$ can be slow. Similar difficulties may be encountered when using a secant method like the


Fig. 1. A plot of $L^{\prime}$ for the problems described in (22) with four different values for $\delta$.


Fig. 2. Potential secant iterates when $L^{\prime}$ is essentially piecewise constant.
one given by Dai and Fletcher [2006]. In their algorithm, significant modifications had to be introduced in order to speed up convergence, particularly for problems where $L^{\prime}$ was nearly piecewise constant.
As the data in Table I indicate, for certain randomly generated problems, Newtontype methods require a small number of iterations to reach the solution. On the other hand, the worst-case complexity could be $O\left(n^{2}\right)$ if the iteration is performed in exact arithmetic. An example demonstrating this worst-case complexity for both Newton's algorithm and the variable fixing algorithm was given in Cominetti et al. [2014], where the authors provide an $L^{\prime}$ for which $n$ Newton iterations are needed to find the root. Here we provide another example that is expressed in terms of the knapsack problem itself. Let us consider the following special case of problem (3):

$$
\begin{equation*}
\min _{\mathbf{x} \in \mathbb{R}^{n}} \frac{1}{2} \mathbf{x}^{\top} \mathbf{x}-\mathbf{1}^{\top} \mathbf{x} \quad \text { subject to } \quad \ell \leq \mathbf{x} \leq \mathbf{1} \text { and } \mathbf{1}^{\top} \mathbf{x}=0 . \tag{23}
\end{equation*}
$$

Consider any lower bound $\ell$ of the following form:

$$
\begin{equation*}
\ell_{1}=\epsilon_{1}, \quad \ell_{i}=\epsilon_{1}+\sum_{j=2}^{i} \frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-1}(n-k)} \text { for all } i \geq 2 \tag{24}
\end{equation*}
$$

where $\epsilon_{i}$ is an arbitrary sequence that satisfies $1>\epsilon_{1}>\epsilon_{2}>\cdots>\epsilon_{n}=0$.
Lemma 3.1. The components of $\ell$ given in (24) satisfy the following:
(1) $1>\ell_{1}>\ell_{2}>\cdots>\ell_{n}$.
(2) For each $i \geq 2$, the following relation holds:

$$
\begin{equation*}
\ell_{i}=\frac{-\sum_{m=1}^{i-1} \ell_{m}}{n-i+1}+\frac{\epsilon_{i}}{\prod_{m=1}^{i-1}(n-m)} \tag{25}
\end{equation*}
$$

Proof.
Part 1. Since the sequence $\epsilon_{i}$ is strictly decreasing and non-negative, it follows that, for any $2 \leq j \leq n$, we have

$$
\epsilon_{j}-(n-j+2) \epsilon_{j-1}<\epsilon_{j}-\epsilon_{j-1}<0
$$

Hence, the terms in the sum in (24) are all negative, which implies that the sequence $\ell_{i}$ is strictly decreasing and $\ell_{1}=\epsilon_{1}<1$.
Part 2. Substituting $\ell_{m}$ using (24), we obtain

$$
\begin{equation*}
\sum_{m=1}^{i-1} \ell_{m}=(i-1) \epsilon_{1}+\sum_{m=1}^{i-1} \sum_{j=2}^{m} \frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-1}(n-k)} \tag{26}
\end{equation*}
$$

Notice that for each $2 \leq j \leq n$, the term $\frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-1}(n-k)}$ appears exactly $(i-j)$ times in the sum of (26). Hence, Equation (26) simplifies to

$$
\sum_{m=1}^{i-1} \ell_{m}=(i-1) \epsilon_{1}+\sum_{j=2}^{i-1}\left(\frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-1}(n-k)}\right)(i-j) .
$$

We make this substitution on the right-hand side of (25), the substitution (24) on the left, and multiply by $n-i+1$ to obtain

$$
\begin{aligned}
& \epsilon_{1}(n-i+1)+\sum_{j=2}^{i}\left(\frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-1}(n-k)}\right)(n-i+1) \\
& \quad=-\epsilon_{1}(i-1)-\sum_{j=2}^{i-1} \frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-1}(n-k)}(i-j)+\frac{\epsilon_{i}}{\prod_{k=1}^{i-2}(n-k)} .
\end{aligned}
$$

In the special case $i=2$, this relation remains valid if the products are treated as 1 and the sums are treated as 0 when the lower limit exceeds the upper limit. We rearrange this relation to get

$$
\begin{equation*}
n \epsilon_{1}+\sum_{j=2}^{i} \frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-2}(n-k)}=\frac{\epsilon_{i}}{\prod_{k=1}^{i-2}(n-k)} . \tag{27}
\end{equation*}
$$

Hence, Equation (25) is equivalent to Equation (27). We prove Equation (27) by induction.

If $i=2$, then the left-hand side of (27) is $n \epsilon_{1}+\epsilon_{2}-n \epsilon_{1}=\epsilon_{2}$, which equals the righthand side. Now suppose that (27) holds for some $i \geq 2$. By the induction hypothesis, we have

$$
\begin{aligned}
n \epsilon_{1}+\sum_{j=2}^{i+1} \frac{\epsilon_{j}-(n-j+2) \epsilon_{j-1}}{\prod_{k=1}^{j-2}(n-k)} & =\frac{\epsilon_{i+1}-(n-i+1) \epsilon_{i}}{\prod_{k=1}^{i-1}(n-k)}+\frac{\epsilon_{i}}{\prod_{k=1}^{i-2}(n-k)} \\
& =\frac{\epsilon_{i+1}}{\prod_{k=1}^{i-1}(n-k)}
\end{aligned}
$$

But this is exactly the assertion made by (27) for $i+1$. Hence, the induction step has been established.

Inserting $i=n$ into Equation (25), we conclude that $\mathbf{1}^{\top} \ell=0$ since $\epsilon_{n}=0$. Hence, $\mathbf{x}=\ell$ is the only feasible point for Problem (23). We now show that if the variable fixing algorithm is applied to the problem (23) with $\ell$ chosen according to Equation (24), then only one variable is fixed at a lower bound in each iteration. Since the optimal solution is $\mathbf{x}^{*}=\ell, n$ iterations are required. Since the time to perform iteration $k$ is $\Omega(n-k)$, the total running time is $\Omega\left(n^{2}\right)$.

Proposition 3.2. At iteration $k$ of the variable fixing algorithm applied to Equation (23) with $\ell$ chosen according to Equation (24), exactly one component of the optimal solution $\mathbf{x}^{*}$ is fixed, namely $x_{k}^{*}=\ell_{k}$.

Proof. The proof is by induction on $k$. For $k=1, \mathcal{F}_{k}=\{1,2, \ldots, n\}, b_{k}=0$, and the solution of the equality constrained reduced problem (12) is $\mathbf{x}_{1}=\mathbf{0}$. Since $\ell_{1}=\epsilon_{1}>0$, the first component of $\mathbf{x}_{1}$ violates the lower bound: $x_{11}=0<\ell_{1}$. On the other hand, we now show that $\ell_{i}<x_{1 i}<u_{i}$ for $i \geq 2$. Since $\epsilon_{2}<\epsilon_{1}$, we have

$$
l_{2}=\epsilon_{1}+\frac{\epsilon_{2}-n \epsilon_{1}}{n-1}<\epsilon_{1}+\frac{\epsilon_{1}-n \epsilon_{1}}{n-1}=\epsilon_{1}-\epsilon_{1}=0=x_{12} .
$$

By part 1 of Lemma 3.1, $\ell_{i}$ is a strictly decreasing function of $i$. Hence, for all $i \geq 2$,

$$
\ell_{i} \leq \ell_{2}<0=x_{1 i}<u_{i}=1 .
$$

This implies that the first component of $\mathbf{x}_{1}$ is the only component that is fixed at iteration 1; moreover, it is fixed at the lower bound. So $\mathcal{B}_{2}=\{1\}$ and $\mathbf{x}_{1}^{*}=\ell_{1}$. This completes the base case.

Proceeding by induction, suppose that for some $k \geq 2, \mathcal{B}_{k}=\{1,2, \ldots, k-1\}$ and $x_{i}^{*}=\ell_{i}$ for all $i<k$. We will show that $\mathcal{B}_{k+1}=\{1,2, \ldots, k\}$ and $x_{k}^{*}=\ell_{k}$. At iteration $k$, the reduced problem is

$$
\begin{equation*}
\min \left\{\sum_{i=k}^{n} \frac{1}{2} x_{i}^{2}-x_{i}: \sum_{i=k}^{n} x_{i}=-\sum_{i=1}^{k-1} \ell_{i}\right\} . \tag{28}
\end{equation*}
$$

The solution is

$$
\begin{equation*}
x_{k i}=\frac{-\sum_{m=1}^{k-1} \ell_{m}}{n-k+1}, \quad i \geq k \tag{29}
\end{equation*}
$$

By Equation (25),

$$
\begin{equation*}
x_{k i}=\ell_{k}-\frac{\epsilon_{k}}{\prod_{m=1}^{k-1}(n-m)}, \quad i \geq k \tag{30}
\end{equation*}
$$

In particular, $x_{k k}<\ell_{k}$.

Table III. Values of $\ell_{i}$ in Double Precision Arithmetic When $n=100$ and $\epsilon_{i}=1-(i / 100)$

| $i$ | $\ell_{i} \times 100$ |
| :---: | :---: |
| 1 | 99.00000000000000000 |
| 2 | -0.010101010101010111 |
| 3 | -1.000103071531643038 |
| 4 | -1.010102072694119869 |
| 5 | -1.010204092701331297 |
| 6 | -1.010205144342275242 |
| 7 | -1.010205155295680612 |
| 8 | -1.010205155410966865 |
| 9 | -1.010205155412193141 |
| 10 | -1.010205155412206325 |
| 11 | -1.010205155412206499 |
| 12 | -1.010205155412206499 |

By the definition (24) of $\ell$, we have

$$
\ell_{k+1}=\ell_{k}+\frac{\epsilon_{k+1}-(n-k+1) \epsilon_{k}}{\prod_{j=1}^{k}(n-j)} .
$$

Substituting for $\ell_{k}$ using Equation (30) yields

$$
\ell_{k+1}=x_{k i}+\frac{\epsilon_{k+1}-\epsilon_{k}}{\prod_{j=1}^{k}(n-j)}<x_{k i} \text { for } i>k .
$$

By Lemma 3.1, $\ell_{i}$ is strictly decreasing. Hence, for every $i>k$ we have

$$
\ell_{i} \leq \ell_{k+1}<x_{k i}=x_{k k}<\ell_{k}<1 .
$$

It follows that $x_{k}^{*}=\ell_{k}$ while $\ell_{i}<x_{k i}<u_{i}$ for $i>k$. So $\mathcal{B}_{k+1}=\mathcal{B}_{k} \cup\{k\}$. This completes the induction step.

Remark 3.3. Proposition 3.2 is a theoretical result in the sense that the computations must be performed with exact arithmetic; in finite precision arithmetic, the $\ell_{i}$ sequence quickly approaches a limit as is seen in Table III.

Remark 3.4. If Newton's method is applied to Problem (23) starting from the $\lambda_{1}$ iterate of the variable fixing algorithm, then it will generate exactly the same iterates as the variable fixing algorithm. The equivalence between Newton's method and the variable fixing method is based on Theorem 2.5. During the proof of Proposition 3.2, we showed that the solution $\mathbf{x}_{k}$ of the subproblem (28) possessed exactly one component $x_{k k}$ that violated the lower bound constraint $x_{k k} \geq \ell_{k}$. If $\lambda_{k}$ is the multiplier associated with the constraint in Problem (28), then

$$
x_{i}\left(\lambda_{k}\right)=x_{k}\left(\lambda_{k}\right)<\ell_{k}<\ell_{i}
$$

for all $i<k$, since $\ell_{i}$ is a decreasing function of $i$ by Lemma 3.1. Consequently, $\mathcal{A}\left(\lambda_{k}\right)=\mathcal{B}_{k+1}$, and, by Theorem 2.5, Newton's method produces the same iterate as the variable fixing method. Hence, Newton's method has complexity $\Omega\left(n^{2}\right)$ when applied to Problem (23).

## 4. VARIABLE FIXING WHEN SOME DIAGONAL ELEMENTS VANISH

When some components of $\mathbf{d}$ vanish, it may not be possible to apply Newton's method to the equation $L^{\prime}(\lambda)=0$ since $L^{\prime \prime}(\lambda)$ could vanish. In this section we develop a version of the variable fixing algorithm that can used even when one or more components of $\mathbf{d}$
vanish. Assuming $L(\lambda)>-\infty$ (that is, $\lambda$ lies in the domain of $L$ ), the set of minimizers $\mathbf{X}(\lambda)$ for the dual function (6) is

$$
\begin{equation*}
\mathbf{X}(\lambda)=\arg \min \{\mathcal{L}(\mathbf{x}, \lambda): \ell \leq \mathbf{x} \leq \mathbf{u}\} . \tag{31}
\end{equation*}
$$

If $d_{i}>0$, then the $i$ th component of $\mathbf{X}(\lambda)$ is unique and is given by (8). The associated break points (9) are $\left(y_{i}-\ell_{i} d_{i}\right) / a_{i}$ and $\left(y_{i}-u_{i} d_{i}\right) / a_{i}$. Between the break points, $X_{i}(\lambda)=$ $\left(y_{i}-\lambda a_{i}\right) / d_{i}$, while outside the break point interval, $X_{i}(\lambda)=\ell_{i}$ or $u_{i}$. If $d_{i}=0$ and $\lambda$ lies in the interior of the domain of the dual function, then

$$
X_{i}(\lambda)=\arg \min \left\{\left(\lambda a_{i}-y_{i}\right) x_{i}: \ell_{i} \leq x_{i} \leq u_{i}\right\}=\left\{\begin{array}{cl}
\ell_{i} & \text { if } a_{i} \lambda>y_{i} \\
{\left[\ell_{i}, u_{i}\right]} & \text { if } a_{i} \lambda=y_{i} \\
u_{i} & \text { if } a_{i} \lambda<y_{i}
\end{array}\right.
$$

The interval $\left[\ell_{i}, u_{i}\right]$ corresponds to the break point $\lambda=y_{i} / a_{i}$. On either side of the break point, $X_{i}(\lambda)$ equals either $\ell_{i}$ or $u_{i}$. Hence, for any $\mathbf{d} \geq \mathbf{0}, \mathbf{X}(\lambda)$ is a linear function of $\lambda$ on the interior of any interval located between break points.

By Clarke [1975, Thm. 2.1] or Danskin [1967], the subdifferential of $L$ can be expressed

$$
\begin{equation*}
\partial L(\lambda)=\left[L^{\prime}\left(\lambda^{+}\right), L^{\prime}\left(\lambda^{-}\right)\right] \tag{32}
\end{equation*}
$$

where

$$
L^{\prime}\left(\lambda^{+}\right)=\min \left\{\mathbf{a}^{\top} \mathbf{x}-b: \mathbf{x} \in \mathbf{X}(\lambda)\right\} \quad \text { and } \quad L^{\prime}\left(\lambda^{-}\right)=\max \left\{\mathbf{a}^{\top} \mathbf{x}-b: \mathbf{x} \in \mathbf{X}(\lambda)\right\}
$$

Since $L$ is concave, its subdifferential is monotone; in particular, if $\lambda_{1}<\lambda_{2}, g_{1} \in \partial L\left(\lambda_{1}\right)$, and $g_{2} \in \partial L\left(\lambda_{2}\right)$, then $g_{1} \geq g_{2}$. The generalization of Proposition 2.1 is the following:

Proposition 4.1. If $\mathbf{d} \geq \mathbf{0}$ and there exists an optimal solution $\mathbf{x}^{*}$ of (3), then there exists a maximizer $\lambda^{*}$ of the dual function, $0 \in \partial L\left(\lambda^{*}\right)$, and $\mathbf{x}^{*} \in \mathbf{X}\left(\lambda^{*}\right)$. Moreover, any $\mathbf{x}^{*} \in \mathbf{X}\left(\lambda^{*}\right)$ with $\mathbf{a}^{\top} \mathbf{x}^{*}=b$ is optimal in (3).

Proof. The existence of a maximizer $\lambda^{*}$ of the dual function along with the optimality conditions for $\lambda^{*}$ and $\mathbf{x}^{*}$ are well-known properties of a concave optimization problem (see Luenberger and Ye [2008] and Rockafellar [1970]).
For any given $\lambda$ and for all sufficiently small $\epsilon>0$, the sets $\mathbf{X}(\lambda+\epsilon)$ and $\mathbf{X}(\lambda-\epsilon)$ are singletons. We define

$$
\mathbf{X}\left(\lambda^{+}\right):=\lim _{\epsilon \rightarrow 0^{+}} \mathbf{X}(\lambda+\epsilon) \quad \text { and } \quad \mathbf{X}\left(\lambda^{-}\right):=\lim _{\epsilon \rightarrow 0^{+}} \mathbf{X}(\lambda-\epsilon) .
$$

The following proposition extends Proposition 2.2 to the case $\mathbf{d} \geq \mathbf{0}$ :
Proposition 4.2. If $\mathbf{d} \geq \mathbf{0}, \mathbf{a}>\mathbf{0}$, and $\mathbf{x}^{*}$ is optimal in (3), then for any $\lambda$ in the domain of $L$, we have the following:
(1) If $L^{\prime}\left(\lambda^{+}\right) \geq 0$, then $x_{i}^{*}=\ell_{i}$ for every $i$ such that $X_{i}\left(\lambda^{+}\right)=\ell_{i}$.
(2) If $L^{\prime}\left(\lambda^{-}\right) \leq 0$, then $x_{i}^{*}=u_{i}$ for every $i$ such that $X_{i}\left(\lambda^{-}\right)=u_{i}$.

Proof. Since the proofs of parts 1 and 2 are similar, we only prove part 1 . Let $\lambda^{*}$ maximize $L$. By Proposition 4.1, $0 \in \partial L\left(\lambda^{*}\right)$. If $L^{\prime}\left(\lambda^{+}\right) \geq 0$, it follows from the monotonicity of $\partial L$, that $\lambda \leq \lambda^{*}$.

Case 1. First suppose that $\lambda=\lambda^{*}$ and $X_{i}\left(\lambda^{+}\right)=\ell_{i}$. Since $0 \in \partial L\left(\lambda^{*}\right)=$ $\partial L(\lambda)=\left[L^{\prime}\left(\lambda^{+}\right), L^{\prime}\left(\lambda^{-}\right)\right]$, we conclude that $L^{\prime}\left(\lambda^{+}\right) \leq 0$. By assumption, $L^{\prime}\left(\lambda^{+}\right) \geq 0$. So it follows that $L^{\prime}\left(\lambda^{+}\right)=0$. If $d_{i}>0$, then $X_{i}(\lambda)=X_{i}\left(\lambda^{+}\right)=\ell_{i}$. So since $\lambda=\lambda^{*}, X_{i}\left(\overline{\lambda^{*}}\right)=\ell_{i}$. And since $\mathbf{x}^{*} \in \mathbf{X}\left(\lambda^{*}\right)$ by Proposition 4.1, $x_{i}^{*}=\ell_{i}$. If $d_{i}=0$, then since $X_{i}\left(\lambda^{+}\right)=\ell_{i}$, it follows that $\lambda a_{i}-y_{i}=\lambda^{*} a_{i}-y_{i} \geq 0$. If $\lambda^{*} a_{i}-y_{i}>0$, then $x_{i}^{*}=X_{i}\left(\lambda^{*}\right)=\ell_{i}$. If $\lambda a_{i}-y_{i}=0$, then $X_{i}(\lambda)=X_{i}\left(\lambda^{*}\right)=\left[\ell_{i}, u_{i}\right]$. Since $\mathbf{x}^{*}$ is optimal in Problem (3), we have $\mathbf{a}^{\top} \mathbf{x}^{*}=b$. If

```
ALGORITHM 2: Solve a Separable Quadratic Knapsack Problem (3) with \(\mathbf{d} \geq \mathbf{0}\) and \(\mathbf{a}>\mathbf{0}\).
Input: Vectors \(\mathbf{a}, \mathbf{d}, \mathbf{y}, \ell, \mathbf{u}\) and scalar \(b\) defining the problem (3); initial guess \(\lambda_{1}\)
Output: The solution vector \(\mathbf{x}\)
\(k=1 ; \mathcal{L}_{1}=\emptyset ; \mathcal{U}_{1}=\emptyset ;\)
while true do
    Step 1: Generate \(\lambda_{k}\) (use initial guess if \(k=1\) ).
    Step 2: Stopping Criterion:
    if \(0 \in \partial L_{k}\left(\lambda_{k}\right)\) then
        find \(\mathbf{x}^{*} \in \mathbf{X}\left(\lambda_{k}\right)\) such that \(\mathbf{a}^{\top} \mathbf{x}^{*}=b\), where
            \(L_{k}(\lambda):=\min \left\{\mathcal{L}(\mathbf{x}, \lambda): \ell \leq \mathbf{x} \leq \mathbf{u}, \quad x_{i}=\ell_{i} \forall i \in \mathcal{L}_{k}, \quad x_{i}=u_{i} \forall i \in \mathcal{U}_{k}\right\}\)
        break
    end
    Step 3: Variable fixing:
    Define \(\mathcal{F}_{k}=\left(\mathcal{U}_{k} \cup \mathcal{L}_{k}\right)^{c}\).
    if \(L_{k}^{\prime}\left(\lambda_{k}^{+}\right)>0\) then
        \(\mathcal{L}_{k+1}=\mathcal{L}_{k} \cup\left\{i \in \mathcal{F}_{k}: X_{i}\left(\lambda_{k}^{+}\right)=\ell_{i}\right\}\) and \(\mathcal{U}_{k+1}=\mathcal{U}_{k}\).
    end
    if \(L_{k}^{\prime}\left(\lambda_{k}^{-}\right)<0\) then
        \(\mathcal{U}_{k+1}=\mathcal{U}_{k} \cup\left\{i \in \mathcal{F}_{k}: X_{i}\left(\lambda_{k}^{-}\right)=u_{i}\right\}\) and \(\mathcal{L}_{k+1}=\mathcal{L}_{k}\).
    end
    \(k=k+1\)
end
```

$x_{i}^{*}>\ell_{i}$ and $a_{i}>0$, then $x_{i}^{*}-\epsilon>\ell_{i}$ for $\epsilon>0$ sufficiently small. Hence, $\mathbf{x}^{*}-\epsilon \mathbf{e}_{i} \in \mathbf{X}\left(\lambda^{*}\right)$, where $\mathbf{e}_{i}$ is the $i$ th column of the identity. Since $a_{i}>0, \mathbf{a}^{\top}\left(\mathbf{x}^{*}-\epsilon \mathbf{e}_{i}\right)-b=-a_{i} \epsilon<0$; it follows from Equation (32) that $-a_{i} \epsilon \in \partial L\left(\lambda^{*}\right)$. This contradicts the fact that $\partial L\left(\lambda^{*}\right)=$ $\left[L^{\prime}\left(\lambda^{*+}\right), L^{\prime}\left(\lambda^{*-}\right)\right]$ where $L^{\prime}\left(\lambda^{*+}\right)=0$. Consequently, $x_{i}^{*}=\ell_{i}$.

Case 2. Suppose that $\lambda<\lambda^{*}$ and $X_{i}\left(\lambda^{+}\right)=\ell_{i}$. If $d_{i}>0$, then by Equation (8), $\frac{y_{i}-\lambda a_{i}}{d_{i}} \leq \ell_{i}$. Since $a_{i}>0$ and $\lambda<\lambda^{*}$, we have

$$
\frac{y_{i}-\lambda^{*} a_{i}}{d_{i}}<\frac{y_{i}-\lambda a_{i}}{d_{i}} \leq \ell_{i} .
$$

Hence, by Equation (8), $x_{i}^{*}=\ell_{i}$. If $d_{i}=0$, then since $X_{i}\left(\lambda^{+}\right)=\ell_{i}$ and $\lambda<\lambda^{*}$, we have

$$
\lambda^{*} a_{i}-y_{i}>\lambda a_{i}-y_{i} \geq 0
$$

since $a_{i}>0$. Again, it follows that $X_{i}\left(\lambda^{*}\right)=\ell_{i}$.
Algorithm 2 is a generic approach for solving Problem (3) in the case where $\mathbf{d} \geq \mathbf{0}$. In Algorithm 2, the sets $\mathcal{L}_{k}$ and $\mathcal{U}_{k}$ store components whose optimal values have been determined to lie at a lower or upper bound, respectively. The minimizing argument of $L_{k}(\lambda)$ in Step 2 is given by

$$
X_{i}^{k}(\lambda)=\left\{\begin{array}{cc}
X_{i}(\lambda) & \text { if } i \in \mathcal{F}_{k}, \\
\ell_{i} & \text { if } i \in \mathcal{L}_{k}, \\
u_{i} & \text { if } i \in \mathcal{U}_{k}
\end{array}\right.
$$

As we will see in Lemma 4.4, $L_{k}$ in Algorithm 2 may be replaced by $L$ without affecting the algorithm. Hence, in the case $\mathbf{d}>\mathbf{0}$, Algorithm 2 is equivalent to Algorithm 1 since then $\partial L(\lambda)=L^{\prime}(\lambda)$ and $\mathbf{X}\left(\lambda^{ \pm}\right)=\mathbf{x}(\lambda)$ for all $\lambda \in \mathbb{R}$.

As in Algorithm 1, the convergence properties of Algorithm 2 depend on how $\lambda_{k}$ is chosen in Step 1. Let us define the set

$$
\mathcal{Z}=\left\{j: d_{j}=0\right\}
$$

When $\mathcal{Z} \neq \emptyset$, the variable fixing iterates $\lambda_{k}^{F}$ defined in (14) and the Newton iterates (15) may be invalid, due to division by zero. However, the variable fixing algorithm has a natural extension to the case $\mathcal{Z} \neq \emptyset$, as we will now show.

Given an instance of Problem (3) with $\mathbf{d} \geq \mathbf{0}$, let $\epsilon>0$ and define $\mathbf{d}^{\epsilon}$ by

$$
d_{i}^{\epsilon}=\left\{\begin{aligned}
\epsilon & \text { if } i \in \mathcal{Z} \\
d_{i} & \text { if } i \notin \mathcal{Z}
\end{aligned}\right.
$$

$i=1,2, \ldots, n$. Consider the perturbed problem

$$
\begin{equation*}
\min _{\mathbf{x} \in \mathbb{R}^{n}}\left\{\frac{1}{2} \mathbf{x}^{\top} \mathbf{D}^{\epsilon} \mathbf{x}-\mathbf{y}^{\top} \mathbf{x}: \quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad \mathbf{a}^{\top} \mathbf{x}=b\right\} \tag{33}
\end{equation*}
$$

where $\mathbf{D}^{\epsilon}$ is the diagonal matrix with $\mathbf{d}^{\epsilon}$ on the diagonal. Since $d_{i}^{\epsilon}>0$ for every $i$, the $k$ th variable fixing iterate for (33) is well defined and is given by

$$
\begin{aligned}
\lambda_{k, \epsilon}^{F} & =\frac{-b_{k}+\sum_{i \in \mathcal{F}_{k} \backslash \mathcal{Z}} \frac{y_{i} a_{i}}{d_{i}}+\sum_{i \in \mathcal{F}_{k} \cap \mathcal{Z}} \frac{y_{i} a_{i}}{\epsilon}}{\sum_{i \in \mathcal{F}_{k} \backslash \mathcal{Z}} \frac{a_{i}^{2}}{d_{i}}+\sum_{i \in \mathcal{F}_{k} \cap \mathcal{Z}} \frac{a_{i}^{2}}{\epsilon}} \\
& =\frac{-\epsilon b_{k}+\epsilon \sum_{i \in \mathcal{F}_{k} \backslash \mathcal{Z}} \frac{y_{i} a_{i}}{d_{i}}+\sum_{i \in \mathcal{F}_{k} \cap \mathcal{Z}} y_{i} a_{i}}{\epsilon \sum_{i \in \mathcal{F}_{k} \backslash \mathcal{Z}} \frac{a_{i}^{2}}{d_{i}}+\sum_{i \in \mathcal{F}_{k} \cap \mathcal{Z}} a_{i}^{2}}
\end{aligned}
$$

Hence when $\mathcal{Z} \cap \mathcal{F}_{k} \neq \emptyset$, we have

$$
\lambda_{k, 0}^{F}:=\lim _{\epsilon \rightarrow 0^{+}} \lambda_{k, \epsilon}^{F}=\frac{\sum_{i \in \mathcal{Z} \cap \mathcal{F}_{k}} y_{i} a_{i}}{\sum_{i \in \mathcal{Z} \cap \mathcal{F}_{k}} a_{i}^{2}}
$$

This motivates the following theorem.
THEOREM 4.3. Let $\mathbf{d} \geq \mathbf{0}, \mathbf{a}>\mathbf{0}$, and suppose that (3) has an optimal solution. For each iteration $k \geq 1$ of $\overline{A l}$ gorithm 2 , let $\lambda_{k}$ be defined by

$$
\lambda_{k}=\left\{\begin{array}{cl}
\lambda_{k, 0}^{F} & \text { if } \mathcal{Z} \cap \mathcal{F}_{k} \neq \emptyset  \tag{34}\\
\lambda_{k}^{F} & \text { if } \mathcal{Z} \cap \mathcal{F}_{k}=\emptyset
\end{array}\right.
$$

Then Algorithm 2 converges in finitely many iterations to some $\lambda^{*}$ such that $0 \in \partial L\left(\lambda^{*}\right)$. Hence, any $\mathbf{x}^{*} \in \mathbf{X}\left(\lambda^{*}\right)$ such that $\mathbf{a}^{\top} \mathbf{x}^{*}=b$ is optimal in Problem (3).

The proof of Theorem 4.3 is based on the following lemma.
Lemma 4.4. Let $\mathbf{d} \geq \mathbf{0}, \mathbf{a}>\mathbf{0}$, and suppose that Problem (3) has an optimal solution. For each iteration $k \geq 1$, define

$$
\begin{aligned}
& \lambda_{k}^{L}:=\sup \left\{\lambda_{i}: 1 \leq i \leq k-1, L_{i}^{\prime}\left(\lambda_{i}^{+}\right)>0\right\} \\
& \lambda_{k}^{R}:=\inf \left\{\lambda_{i}: 1 \leq i \leq k-1, L_{i}^{\prime}\left(\lambda_{i}^{-}\right)<0\right\}
\end{aligned}
$$

where $\sup \emptyset:=-\infty$ and $\inf \emptyset:=\infty$. Then, for every $k \geq 1$ such that $\mathcal{Z} \cap \mathcal{F}_{k} \neq \emptyset$, the following hold:
(1) $\lambda_{k} \in\left(\lambda_{h}^{L}, \lambda_{k}^{R}\right)$.
(2) If $0 \notin \partial L_{k}\left(\lambda_{k}\right)$, then either $\left|\mathcal{L}_{k+1}\right|>\left|\mathcal{L}_{k}\right|$ or $\left|\mathcal{U}_{k+1}\right|>\left|\mathcal{U}_{k}\right|$.
(3) $L_{k}(\lambda)=L(\lambda)$ for every $\lambda \in\left(\lambda_{k}^{L}, \lambda_{k}^{R}\right)$.
(4) For every $\lambda \in \mathbb{R}$ with $0 \in \partial L_{k}(\lambda)$, we have $0 \in \partial L(\lambda)$.

Proof. Parts 1 and 2: Since $\mathcal{Z} \cap \mathcal{F}_{k} \neq \emptyset$, we have

$$
\lambda_{k}=\frac{\sum_{i \in \mathcal{Z} \cap \mathcal{F}_{k}} y_{i} a_{i}}{\sum_{j \in \mathcal{Z} \cap \mathcal{F}_{k}} a_{j}^{2}}=\sum_{i \in \mathcal{Z} \cap \mathcal{F}_{k}} \frac{y_{i}}{a_{i}}\left(\frac{a_{i}^{2}}{\sum_{j \in \mathcal{Z} \cap \mathcal{F}_{k}} a_{j}^{2}}\right)=\sum_{i \in \mathcal{Z} \cap \mathcal{F}_{k}} \alpha_{i, k}\left(\frac{y_{i}}{a_{i}}\right),
$$

where $\alpha_{i, k}=\frac{a_{i}^{2}}{\sum_{j \in \mathcal{Z} \cap \mathcal{F}_{k}} a_{j}^{2}}>0$ and $\sum_{i \in \mathcal{Z} \cap \mathcal{F}_{k}} \alpha_{i, k}=1$. Hence, $\lambda_{k}$ is a convex combination of breakpoints associated with components in $\mathcal{Z} \cap \mathcal{F}_{k}$. In particular, there exist $i, j \in \mathcal{Z} \cap \mathcal{F}_{k}$ such that

$$
\begin{equation*}
\frac{y_{i}}{a_{i}} \leq \lambda_{k} \leq \frac{y_{j}}{a_{j}} \tag{35}
\end{equation*}
$$

Next, we claim that $\lambda_{k}^{L}<\frac{y_{i}}{a_{i}}$. Indeed, suppose by way of contradiction that this is not true. Then, by definition of $\lambda_{k}^{L}$, there is some (smallest) $l \leq k-1$ such that $L_{l}^{\prime}\left(\lambda_{l}^{+}\right)>0$ and $\lambda_{l} \geq \frac{y_{i}}{a_{i}}$. Since $\lambda_{l} \geq \frac{y_{i}}{a_{i}}, X_{i}\left(\lambda_{l}^{+}\right)=\ell_{i}$ by definition of $\mathbf{X}$, and on Step 3 of the $l$ th iteration of Algorithm 2, $i \in \mathcal{L}_{l+1} \backslash \mathcal{L}_{l} \subseteq \mathcal{L}_{k}$, contradicting $i \in \mathcal{F}_{k}$. By a similar argument, $\frac{y_{j}}{a_{j}}<\lambda_{k}^{R}$. Hence, by Equation (35) $\lambda_{k}^{L}<\lambda_{k}<\lambda_{k}^{R}$, which proves Part 1.

Now, suppose that $0 \notin \partial L_{k}\left(\lambda_{k}\right)$. Then by Equation (32) either $L_{k}^{\prime}\left(\lambda_{k}^{+}\right)>0$ or $L_{k}^{\prime}\left(\lambda_{k}^{-}\right)<0$. Notice that (35) implies that $X_{i}\left(\lambda_{k}^{+}\right)=\ell_{i}$ and $X_{j}\left(\lambda_{k}^{-}\right)=u_{j}$, by definition of $\mathbf{X}$. Hence, if $L_{k}^{\prime}\left(\lambda_{k}^{+}\right)>0$, then by Step 3 of the $k$ th iteration of Algorithm 2 we have $i \in \mathcal{L}_{k+1} \backslash \mathcal{L}_{k}$; whereas if $L_{k}^{\prime}\left(\lambda_{k}^{-}\right)<0$, then $j \in \mathcal{U}_{k+1} \backslash \mathcal{U}_{k}$. This completes the proof of Part 2.

Parts 3 and 4: We prove Parts 3 and 4 by induction on $k$. The base cases where $k=1$ are trivial, since $L_{1}=L$. So, suppose that for some $s \geq 1$, we have $\mathcal{Z} \cap \mathcal{F}_{k} \neq \emptyset$ and Parts 3 and 4 hold for every $k \leq s$. Suppose that $\mathcal{Z} \cap \mathcal{F}_{s+1} \neq \emptyset$. We will show that Parts 3 and 4 hold for $k=s+1$.

Since the algorithm did not stop on iteration $s$, we must have $0 \notin \partial L_{s}\left(\lambda_{s}\right)$. We assume without loss of generality that $L_{s}^{\prime}\left(\lambda_{s}^{+}\right)>0$ (the proof in the case $L_{s}^{\prime}\left(\lambda_{s}^{-}\right)<0$ is similar). By the induction hypothesis, $L_{s}(\lambda)=L(\lambda)$ for any $\lambda \in\left(\lambda_{s}^{L}, \lambda_{s}^{R}\right) \supseteq\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right)$. Hence, we will be done with Part 3 when we show that $L_{s+1}(\lambda)=L_{s}(\lambda)$ for every $\lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right)$.

Let $\lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right)$. Since $L_{s}^{\prime}\left(\lambda_{s}^{+}\right)>0$ by assumption, Step 3 of the algorithm gives $\mathcal{U}_{s+1}=\mathcal{U}_{s}$ and $\mathcal{L}_{s} \subset \mathcal{L}_{s+1}$. Hence, $X_{i}^{s+1}(\lambda)=X_{i}^{s}(\lambda)=u_{i}$ for all $i \in \mathcal{U}_{s+1}$ and $X_{i}^{s+1}(\lambda)=$ $X_{i}^{s}(\lambda)=\ell_{i}$ for all $i \in \mathcal{L}_{s}$. Since $\mathcal{F}^{s+1} \subset \mathcal{F}^{s}$, it follows that $X_{i}^{s+1}(\lambda)=X_{i}^{s}(\lambda)$ for all $i \in \mathcal{F}_{s+1}$. Hence, the only indices where $X_{i}^{s+1}(\lambda) X_{i}^{s}(\lambda)$ could differ are those $i \in \mathcal{L}_{s+1} \backslash \mathcal{L}_{s}$. We will now show that for every $i \in \mathcal{L}_{s+1} \backslash \mathcal{L}_{s}, X_{i}^{s}(\lambda)=\ell_{i}=X_{i}^{s+1}(\lambda)$. This will imply $\mathbf{X}^{s+1}(\lambda)=\mathbf{X}^{s}(\lambda)$, and therefore $L_{s+1}(\lambda)=L_{s}(\lambda)$.

Let $i \in \mathcal{L}_{s+1} \backslash \mathcal{L}_{s}$. Since $L_{s}^{\prime}\left(\lambda_{s}^{+}\right)>0$ by assumption, Step 3 of the algorithm implies

$$
\begin{equation*}
X_{i}\left(\lambda_{s}^{+}\right)=\ell_{i} \tag{36}
\end{equation*}
$$

By Part $1, \lambda_{s} \in\left(\lambda_{s}^{L}, \lambda_{s}^{R}\right)$; hence by definition of $\lambda_{s+1}^{L}$,

$$
\begin{equation*}
\lambda_{s+1}^{L}=\lambda_{s} \tag{37}
\end{equation*}
$$

Since $\lambda>\lambda_{s+1}^{L}=\lambda_{s}$ and $X_{i}(\cdot)$ is nonincreasing and bounded below by $\ell_{i}$, we have $X_{i}(\lambda)=\ell_{i}$. Thus, $X_{i}^{s}(\lambda)=X_{i}(\lambda)=\ell_{i}=X_{i}^{s+1}(\lambda)$.

Since $i \in \mathcal{L}_{s+1} \backslash \mathcal{L}_{S}$ was arbitrary, we have shown that $\mathbf{X}^{s+1}(\lambda)=\mathbf{X}^{s}(\lambda)$; hence, $L_{s+1}(\lambda)=$ $L_{s}(\lambda)$. Since $\lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right)$ was arbitrary, we have shown that

$$
\begin{equation*}
L_{s+1}(\lambda)=L_{s}(\lambda)=L(\lambda) \quad \forall \lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right), \tag{38}
\end{equation*}
$$

which proves Part 3 for $k=s+1$.

Now we prove Part 4 holds for $k=s+1$. Suppose that $0 \in \partial L_{s+1}(\lambda)$ for some $\lambda \in \mathbb{R}$. Then by Equation (32),

$$
\begin{equation*}
L_{s+1}^{\prime}\left(\lambda^{+}\right) \leq 0 \leq L_{s+1}^{\prime}\left(\lambda^{-}\right) . \tag{39}
\end{equation*}
$$

We claim that $\lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right)$. Once this claim is proved, it will follow from Equation (38) that $\partial L(\lambda)=\partial L_{s+1}(\lambda)$. Since $0 \in \partial L_{s+1}(\lambda)$, we will have $0 \in \partial L(\lambda)$, and the proof will be complete.

First, we show $\lambda_{s+1}^{L}<\lambda$. If $\lambda_{s+1}^{L}=-\infty$, then we are done. So, suppose $\lambda_{s+1}^{L}>-\infty$. Then

$$
\begin{equation*}
L_{s+1}^{\prime}\left(\lambda_{s+1}^{L+}\right)=L_{s}^{\prime}\left(\lambda_{s+1}^{L+}\right)=L_{s}^{\prime}\left(\lambda_{s}^{+}\right)>0, \tag{40}
\end{equation*}
$$

where the first equality follows from Equation (38), the second follows from Equation (37), and the inequality follows from our assumption. Since $\partial L_{s+1}($.$) is monotone$ nonincreasing, Equation (40) and the first half of Equation (39) imply $\lambda_{s+1}^{L}<\lambda$.

Now, we show $\lambda<\lambda_{s+1}^{R}$. If $\lambda_{s+1}^{R}=\infty$, then we are done. So, suppose $\lambda_{s+1}^{R}<\infty$. Then $\lambda_{s+1}^{R}=\lambda_{j}$ for some $j \leq s$ such that $L_{j}^{\prime}\left(\lambda_{j}^{-}\right)<0$. By Equation (38) and the induction assumption for Part 3,

$$
L_{s+1}(\lambda)=L(\lambda)=L_{j}(\lambda) \quad \forall \lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right) \subseteq\left(\lambda_{j}^{L}, \lambda_{j}^{R}\right) .
$$

Hence,

$$
\begin{equation*}
L_{s+1}^{\prime}\left(\lambda_{s+1}^{R-}\right)=L_{j}^{\prime}\left(\lambda_{s+1}^{R-}\right)=L_{j}^{\prime}\left(\lambda_{j}^{-}\right)<0 . \tag{41}
\end{equation*}
$$

So, since $\partial L(\cdot)$ is monotone nonincreasing, Equation (41) and the second half of Equation (39) imply $\lambda<\lambda_{s+1}^{R}$. Thus, $\lambda \in\left(\lambda_{s+1}^{L}, \lambda_{s+1}^{R}\right)$, and our claim is proved. This completes the proof of Part 4 for $k=s+1$.

Therefore, the proof of Parts 3 and 4 is complete by induction.
Remark 4.5. By Lemma 4.4, we may replace $L_{k}$ with $L$ in Algorithm 2. Hence, in the case where $\mathcal{Z}=\emptyset$, Algorithm 2 is equivalent to Algorithm 1, since then $\partial L(\lambda)=L^{\prime}(\lambda)$ and $\mathbf{X}\left(\lambda^{ \pm}\right)=\mathbf{x}(\lambda)$ for all $\lambda \in \mathbb{R}$.

Proof of Theorem 4.3. By Part 2 of Lemma 4.4, since $\left|\mathcal{L}_{k}\right|$ and $\left|\mathcal{U}_{k}\right|$ are bounded from above by $n$, we eventually must reach an iteration $s \geq 1$ such that either

$$
\begin{equation*}
\left[\mathcal{Z} \cap \mathcal{F}_{s}=\emptyset\right] \quad \text { or } \quad\left[\mathcal{Z} \cap \mathcal{F}_{s} \neq \emptyset \text { and } 0 \in \partial L_{s}\left(\lambda_{s}\right)\right] \tag{42}
\end{equation*}
$$

Let $s$ be the first iteration satisfying Equation (42). In the second case of Equation (42), we have $0 \in \partial L\left(\lambda_{s}\right)$ by Part 3 of Lemma 4.4, and we are done. Now suppose instead that the first case holds. Then for $k \geq s, \lambda_{k}=\lambda_{k}^{F}$ and Algorithm 2 reduces to the Variable Fixing Algorithm of Kiwiel [2008] for the problem

$$
\min \quad q(\mathbf{x})
$$

subject to $\quad \ell \leq \mathbf{x} \leq \mathbf{u}, \quad \mathbf{a}^{\top} \mathbf{x}=b, \quad x_{i}=\ell_{i} \forall i \in \mathcal{L}_{s}, \quad$ and $\quad x_{i}=u_{i} \forall i \in \mathcal{U}_{s}$.
(The formula for $\lambda_{k}^{F}$ (14) comes directly from Kiwiel [2008, page 448]). Therefore, by Kiwiel [2008, Theorem 4.1], $\lambda_{k}$ converges to some $\lambda^{*}$ such that

$$
\partial L_{s}\left(\lambda^{*}\right)=L_{s}^{\prime}\left(\lambda^{*}\right)=\mathbf{a}_{s}^{\top} \mathbf{x}_{s}\left(\lambda^{*}\right)-b_{s}=0,
$$

where $\mathbf{a}_{s}$ and $\mathbf{x}_{s}\left(\lambda^{*}\right)$ represent the vectors formed by only including components in $\mathcal{F}_{s}$, and $b_{s}$ is defined by

$$
\begin{equation*}
b_{s}=b-\sum_{i \in \mathcal{L}_{s}} a_{i} \ell_{i}-\sum_{i \in \mathcal{U}_{s}} a_{i} u_{i} . \tag{43}
\end{equation*}
$$

Therefore, by Part 4 of Lemma 4.4, $0 \in \partial L\left(\lambda^{*}\right)$. This completes the proof.

## 5. NAPHEAP

In many problems, the iterates in a variable fixing method or Newton method exhibit a diminishing returns property: Most of the progress towards finding the optimal multiplier is made in the first several iterations, while subsequent iterations yield steps that are much smaller. We propose an algorithm called NAPHEAP that exploits this property. The algorithm first generates an interval ( $\alpha, \beta$ ) (possibly semi-infinite) that brackets a solution of the dual problem (7) by applying a few iterations of a Newton-type method. Then all the break points of the dual function between $\alpha$ and $\beta$ are arranged in heaps. NAPHEAP then starts from either $\alpha$ or $\beta$ and monotonically visits the break points until reaching an optimal dual solution $\lambda^{*}$. Since $L^{\prime}$ is linear on any interval between break points, a solution of the dual problem can be evaluated by linear interpolation. If there are $m$ break points in ( $\alpha, \beta$ ), then the time to update the heap after visiting a break point is proportional to $\log _{2} m$. Similarly, the time to build a heap with $m$ elements is roughly the time to perform $m$ comparisons. Hence, when $l$ break points are visited on the path to the solution, the total running time of the heap phase is proportional to $m+l \log _{2} m$. For details concerning the construction and updating of heaps, see Cormen et al. [2009]. The steps of the NAPHEAP code are partitioned into three phases that we now discuss in more detail.
NAPHEAP Phase 0. If $a_{i}=0$ for some $i$, then the optimal value for $x_{i}$ can evaluated, and $x_{i}$ can be eliminated from the problem. Hence, without loss of generality, we assume that $a_{i} \neq 0$ for all $i$. If $d_{i}=0$ for some $i$, then we evaluate the expressions

$$
\alpha_{0}=\max \left\{\frac{y_{i}}{a_{i}}: d_{i}=0 \text { and } u_{i}=\infty\right\} \quad \text { and } \quad \beta_{0}=\min \left\{\frac{y_{i}}{a_{i}}: d_{i}=0 \text { and } \ell_{i}=-\infty\right\},
$$

where the maximum and the minimum are defined to be $\infty$ and $-\infty$, respectively, when the arguments are empty. $L$ is finite on the interval $\left[\alpha_{0}, \beta_{0}\right]$, the domain of the dual function, and an optimal dual solution $\lambda^{*}$ lies between $\alpha_{0}$ and $\beta_{0}$. If a starting guess $\lambda_{1}$ was not provided, then we evaluate one by performing a single iteration of the variable fixing algorithm. From the sign of $L^{\prime}\left(\lambda_{1}\right)$, and the values of $\alpha_{0}$ and $\beta_{0}$, we can further refine the bracketing interval $(\alpha, \beta)$.

NAPHEAP Phase 1. The bracketing interval generated in phase 0 is further refined by performing a few iterations of a Newton-type method. The optimal number of iterations depends on the problem. When $d_{i}=0$ for some $i$, Version 2.1 of NAPHEAP currently performs $K$ (a user specified number) Newton-type iterations before switching to the break point search. However, when $\mathbf{d}>\mathbf{0}$, an adaptive strategy is employed based on the following observation: Newton's method often converges monotonically to the solution of the dual problem. Hence, if we multiply the Newton step by a scaling factor $\xi>1$ (default 1.1), then the scaled Newton iterate

$$
\lambda_{k+1}=\lambda_{k}-\xi\left(\frac{L^{\prime}\left(\lambda_{k}\right)}{L_{ \pm}^{\prime \prime}\left(\lambda_{k}\right)}\right)
$$

will typically lie on the opposite side of the root from $\lambda_{k}$ when $\lambda_{k}$ is sufficiently close to a solution. We continue to perform a scaled Newton iteration until we reach the opposite side of the root, and then we terminate phase 1 . When the scaled Newton iterate lands outside the current bracketing interval, then we attempt a secant iteration. If the bracketing interval is semi-infinite and the secant iteration cannot be performed, then we use a variable fixing iteration. We impose an upper bound on the number of Newton-type iterations (default (20)).
NAPHEAP Phase 2. Phase 1 generates either the solution of the knapsack problem or an interval $[\alpha, \beta]$ (possibly semi-infinite) that brackets a solution of the dual problem;
either $\alpha$ or $\beta$ was generated by the last iteration of phase 1 . All the break points between $\alpha$ and $\beta$ are evaluated and arranged in a heap. Starting from the last iterate of phase 1 , we monotonically cross the break points until reaching a solution of the dual problem. The solution of Problem (3) is any $\mathbf{x}^{*} \in \mathbf{X}\left(\lambda^{*}\right)$ for which $\mathbf{a}^{\top} \mathbf{x}^{*}=b$.
The heaps that we employ are binary heaps that are trees where each node in the tree has two children. The nodes correspond to break points that lie between $\alpha$ and $\beta$. If we start from $\alpha$ and move to the right on the $\lambda$ axis, then we employ a min-heap; that is, a heap for which the root is the smallest break point and the break points associated with the children of any node are greater than or equal to the break point of the parent. As we move through the break points from left to right, the current root of the tree always contains the next break point. In a similar manner, if we start from $\beta$ and move left, then the relevant break points are arranged in a max-heap since we need to know the largest of the remaining break points.

In more detail, two separate heaps are maintained. The free heap at the current iterate $\lambda_{k}$ consists of break points for those indices $i$ satisfying $\ell_{i}<x_{i}\left(\lambda_{k}\right)<u_{i}$. These break points are values of $\lambda$ in the direction of $\lambda^{*}$ where $x_{i}(\lambda)$ reaches either $\ell_{i}$ or $u_{i}$. The bound heap corresponds to those $i$ for which $x_{i}\left(\lambda_{k}\right)=\ell_{i}$ or $x_{i}\left(\lambda_{k}\right)=u_{i}$ at the current iterate $\lambda_{k}$. The break points in the bound heap are values of $\lambda$ in the direction of $\lambda^{*}$ where $x_{i}(\lambda)$ leaves a bound and becomes free. When we cross over a break point for an index $i$ in the free heap, $x_{i}(\lambda)$ reaches a bound, and it cannot leave this bound since there are no more break points in the direction of $\lambda^{*}$. Hence, this break point can be deleted from the free heap and the bound value for $x_{i}^{*}$ is known. When we cross over a break point for an index $i$ in the bound heap, $x_{i}(\lambda)$ becomes free; however, when $d_{i}>0$, there is another break point in the direction of $\lambda^{*}$ where $x_{i}(\lambda)$ reaches the opposite bound. This new break point is inserted in the free heap. When $d_{i}=0$, there is only one break point $y_{i} / a_{i}$. If this break point lies between $\lambda_{k}$ and $\lambda^{*}$, then it belongs in the bound heap. When we reach this break point in the search process, $x_{i}(\lambda)$ makes a transition from one bound to the opposite bound. In either case, we continue to march across the break points until reaching a point $\lambda^{*}$ where $0 \in \partial L\left(\lambda^{*}\right)$.

## 6. NUMERICAL RESULTS

We now investigate the performance of the NAPHEAP algorithm using the following set of test problems. In these test problems, a statement of the form $C \in[A, B]$ means that $C$ is randomly chosen from the interval $[A, B]$ with a uniform probability distribution.
(1) $d_{i} \in(0,25], a_{i}$ and $y_{i} \in[-25,25], \ell_{i}$ and $u_{i} \in[-15,15]$.
(2) $a_{i} \in[-25,25], y_{i} \in\left[a_{i}-5, a_{i}+5\right], d_{i} \in\left[.5\left|a_{i}\right|, 1.5\left|a_{i}\right|\right], \ell_{i}, u_{i} \in[-15,15]$.
(3) $a_{i} \in[-25,25], y_{i}=a_{i}+5, d_{i}=\left|a_{i}\right|, \ell_{i}, u_{i} \in[-15,15]$.
(4) $a_{i}=1, y_{i} \in[-10,10], d_{i}=1, \ell_{i}=0, u_{i}=1$.
(5) $a_{i} \in(0,25] \cap \mathbb{Z}, y_{i} \in[-10,10], d_{i}=1, \ell_{i}=0, u_{i}=1$.
(6) $d_{i} \in(0,25], y_{i} \in[-25,25], a_{i}=1, \ell_{i}=0, u_{i}=\infty$.
(7) $d_{i} \in\left(0,10^{-6}\right], y_{i} \in[-25,25], a_{i}=1, \ell_{i}=0, u_{i}=\infty$.

Problem sets 4, 5, and 6 are related to graph partitioning (see Hager and Hungerford [2014] and Hager and Krylyuk [1999]), multilevel graph partitioning [Hager et al. 2014], and quadratic resource allocation [Bitran and Hax 1981; Bretthauer and Shetty 1997; Cosares and Hochbaum 1994; Hochbaum and Hong 1995], respectively. We consider the case $r=s=b$ in (1). In problem sets $1-5$, the scalar $b$ in the constraint $\mathbf{a}^{\top} \mathbf{x}=b$ was chosen randomly in the interval $[A, B]$ where

$$
A=\inf \left\{\mathbf{a}^{\top} \mathbf{x}: \ell \leq \mathbf{x} \leq \mathbf{u}\right\} \quad \text { and } \quad B=\sup \left\{\mathbf{a}^{\top} \mathbf{x}: \ell \leq \mathbf{x} \leq \mathbf{u}\right\} .
$$

In problem sets 6 and $7, b$ was chosen randomly in [1, 100].


Fig. 3. Plots of $L^{\prime}$ for test set 4 through test set 1 (in descending order) are shown in panels (a)-(d), respectively. Also shown is the initial variable fixing iterate and one or two Newton iterates.

One deficiency with randomly generated problems is that the resulting dual function has the property that $L^{\prime}$ is essentially the same for each randomly generated problem instance when $n$ is large. Moreover, for large randomly generated problems, $L^{\prime}$ tends to be relatively linear. In Figure 3, we plot $L^{\prime}$ for problems in the test sets 1-4 and $n=6,250,000$. One factor related to the degree of nonlinearity for $L^{\prime}$ is the size and variability of the elements in d. As $d_{i}$ approaches 0 , the two break points (9) associated with $x_{i}$ coalesce into a single point $y_{i} / a_{i}$, and $x_{i}(\lambda)$ becomes piecewise constant with values equal to either $\ell_{i}$ or $u_{i}$. This causes the graph of $L^{\prime}$ to develop flat regions as seen in Figure 1.

For our test sets, the most linear plot, Figure 3(a), corresponds to test set 4 where $d_{i}=1$ for each $i$. For test set 3 (Figure 3(b)), $d_{i}$ lies between 1 and 26, and there is a very slight bend in the $L^{\prime}$ plot. In test set 2 (Figure 3(c)), $d_{i}$ lies between 0 and 37.5, while in test set 1 (Figure $3(\mathrm{~d})$ ) $d_{i}$ is between 0 and 25 . When $d_{i}$ is chosen randomly in the interval $(0,25]$ or $\left(0,37.5\right.$ ], the smaller values of $d_{i}$ can contribute to nonlinearity. Test set 5 is another version of test set 4 with $d_{i}=1$ and a nearly linear $L^{\prime}$. Test set 6 is similar to test set 1 with respect to the range of $d_{i}$; however, $x_{i}^{*}=0$ for most $i$ due to the choice of the bound constraints. Since most of the components of the solution become fixed at 0 and drop out of the problem, the problems in test set 6 are easier than those in set 1 . Test set 7 is difficult in the sense that the $d_{i}$ nearly vanish, the break points for $x_{i}$ nearly coincide, and $L^{\prime}$ is nearly piecewise constant. Consequently, $L^{\prime \prime}(\lambda)$ is zero for most choices of $\lambda$.

Table IV. CPU Times (Seconds) for NAPHEAP and the Newton Method of Cominetti et al. [2014]

|  | NAPHEAP |  |  | Newton |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Test Set | ave | $\min$ | $\max$ | ave | $\min$ | $\max$ |
| 1 | 0.305 | 0.272 | 0.406 | 0.499 | 0.453 | 0.620 |
| 2 | 0.282 | 0.252 | 0.313 | 0.408 | 0.351 | 0.480 |
| 3 | 0.276 | 0.272 | 0.282 | 0.386 | 0.360 | 0.424 |
| 4 | 0.210 | 0.183 | 0.329 | 0.226 | 0.200 | 0.239 |
| 5 | 0.205 | 0.186 | 0.243 | 0.248 | 0.226 | 0.358 |
| 6 | 0.231 | 0.223 | 0.235 | 0.318 | 0.315 | 0.321 |
| 7 | 0.261 | 0.215 | 0.329 | 4.378 | 1.757 | 7.384 |

The near linearity of $L^{\prime}$ in test set 4 can be explained as follows: Assuming $\lambda$ is not a break point, the second derivative of $L$ is given by

$$
L^{\prime \prime}(\lambda)=-\sum_{i \in \mathcal{F}(\lambda)} a_{i}^{2} / d_{i}, \quad \text { where } \mathcal{F}(\lambda)=\left\{i: \ell_{i}<x_{i}(\lambda)<u_{i}\right\} .
$$

Since $a_{i}=d_{i}=u_{i}=1$ and $\ell_{i}=0$ in test set $4, L^{\prime \prime}(\lambda)=-|\mathcal{F}(\lambda)|$ and

$$
\mathcal{F}(\lambda)=\left\{i: 0<y_{i}-\lambda<1\right\}=\left\{i: y_{i} \in(\lambda, \lambda+1)\right\} .
$$

Since the $y_{i}$ are uniformly distributed on $[-10,10]$, it follows that for $\lambda \in[-10,9]$ and for large $n$, it is highly likely that $|\mathcal{F}(\lambda)|$ is near $n / 20$, a constant independent of $\lambda$. Consequently, $L^{\prime \prime}$ is nearly constant and $L^{\prime}$ is nearly linear on $[-10,9]$.
In Table IV we compare the running time of NAPHEAP to that of the Newton code developed in Cominetti et al. [2014]. As pointed out in the Introduction, the Newton algorithm of Cominetti et al. [2014] was faster than the variable fixing method, the secant method, or a median-based method. The experiments were run on a Dell Precision T7610 Workstation with a Dual Intel Xeon Processor E5-2687W v2 (32 cores, $3.4 \mathrm{GHz}, 25.6 \mathrm{MB}$ cache, 192 GB memory). Only one core was used in the experiments. There were 10 trials of each problem set, each of size $n=6,250,000$. The graph of $L^{\prime}$ is essentially the same for each trial in a test set and the main difference between the problems in a test set is the vertical placement of the horizontal line $L^{\prime}(\lambda)=0$ of Figure 3; the placement of this horizontal line corresponds to the choice of $b$. We report the average, minimum, and maximum running time for the 10 trials.

As seen in Table IV, NAPHEAP and the Newton code perform nearly the same with respect to CPU time for test set 4, where $d_{i}=1$ and $L^{\prime}$ is nearly linear (Figure 3(a)). When the nonlinearity increases, as in Figures 3(b) through 3(d), or when the proportion of free components of $\mathbf{x}^{*}$ increases, the advantage of NAPHEAP over Newton increases. For test sets 1 and 2, the number of components $i$ such that $\ell_{i}<x_{i}^{*}<u_{i}$ is on the order of $n / 3$ to $n / 2$, which implies that the work associated with each Newton iteration near the optimal solution is proportional to $n$. The difference in the run times for test set 7 is connected with how the codes handle the special case $L^{\prime \prime}\left(\lambda_{k}\right)=0$. NAPHEAP uses a secant iteration while the Newton code either moves to a nearby break point where $L^{\prime \prime}$ does not vanish or it uses a secant step.
Table V gives the average, maximum, and minimum number of Newton and secant steps taken by NAPHEAP and the Newton method of Cominetti et al. [2014] for each test set. Neither code performed many secant iterations; in test set 7 where $L^{\prime}$ is nearly piecewise constant, NAPHEAP performed on average about three secant iterations and 1 variable fixing iterating before bracketing the solution and switching to the heapbased break point search. The variable fixing iteration was necessary when the interval bracketing the solution was unbounded. The large number of Newton steps taken by

Table V. Number of Newton and Secant Iterations for NAPHEAP and the Newton Method of Cominetti et al. [2014]

|  | NAPHEAP |  |  |  |  |  | Newton |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Newton Iterations |  |  | Secant Iterations |  |  | Newton Iterations |  |  | Secant Iterations |  |  |
| Test Set | ave | min | max | ave | min | max | ave | min | max | ave | min | max |
| 1 | 3.8 | 3 | 8 | 0.0 | 0 | 0 | 6.4 | 4 | 11 | 0.0 | 0 | 0 |
| 2 | 3.4 | 3 | 5 | 0.0 | 0 | 0 | 5.3 | 3 | 8 | 0.0 | 0 | 0 |
| 3 | 3.0 | 3 | 3 | 0.0 | 0 | 0 | 4.7 | 4 | 6 | 0.0 | 0 | 0 |
| 4 | 2.8 | 1 | 3 | 0.1 | 0 | 1 | 4.0 | 4 | 4 | 0.0 | 0 | 0 |
| 5 | 3.5 | 3 | 8 | 0.0 | 0 | 0 | 4.7 | 4 | 10 | 0.0 | 0 | 0 |
| 6 | 10.6 | 10 | 12 | 0.0 | 0 | 0 | 13.0 | 12 | 14 | 0.0 | 0 | 0 |
| 7 | 0.1 | 0 | 1 | 3.2 | 3 | 4 | 92.7 | 49 | 132 | 0.5 | 0 | 1 |

the Newton method in test set 7 is due to many flat regions on the curve $L^{\prime}$, which are not treated efficiently.

## 7. CONCLUSIONS

We have presented a new hybrid algorithm NAPHEAP for the continuous quadratic knapsack problem (1) with $\mathbf{d} \geq \mathbf{0}$. The algorithm is based on maximizing the dual function $L$ in Problem (6), and computing the optimal multiplier associated with the linear constraint. A Newton-type method is used to generate an interval that brackets an optimal dual multiplier. After arranging the break points inside the bracketing interval in heaps, one heap for the bound variables and one heap for the free variables at the current iterate, NAPHEAP starts from one side of the bracketing interval and moves monotonically across the break points until reaching an optimal dual multiplier. In order to handle the case where $d_{i}=0$ for one or more indices, we developed a new version of the variable fixing algorithm and proved its convergence. For the randomly generated test problems of Section 6, the hybrid algorithm NAPHEAP was faster than a pure Newton implementation of Cominetti et al. [2014].

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