

# Computing Optimal Experimental Designs via Interior Point Method \*

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October 12, 2012 (Revised: March 17, June 25, August 16, 2013)

## Abstract

In this paper, we study optimal experimental design problems with a broad class of smooth convex optimality criteria, including the classical A-, D- and  $p$ th mean criterion. In particular, we propose an interior point (IP) method for them and establish its global convergence. Further, by exploiting the structure of the Hessian matrix of the optimality criteria, we derive an explicit formula for computing its rank. Using this result, we then demonstrate that the Newton direction arising in the IP method can be computed efficiently via Sherman-Morrison-Woodbury formula when the size of the moment matrix is small relative to the size of design space. Finally, we compare our IP method with the widely used multiplicative algorithm introduced by Silvey et al. [33] and the standard IP solver SDPT3 [36, 40]. The computational results show that our IP method generally outperforms these two methods in both speed and solution quality.

**Key words:** Optimal experimental design, A-criterion, c-criterion, D-criterion,  $p$ th mean criterion, interior point method

## 1 Introduction

In this paper, we consider the optimal experimental design problems on a given finite design space  $\chi = \{x_1, \dots, x_n\} \subseteq \mathfrak{R}^m$ . In this setting, we consider a coefficient matrix  $K \in \mathfrak{R}^{m \times k}$  of full column rank and the moment matrix defined as

$$\mathcal{M}(w) = \sum_{i=1}^n w_i A_i$$

for  $w \in \Omega := \{w : w_i \geq 0, \sum_{i=1}^n w_i = 1\}$ , where  $A_i$  is a real symmetric positive semidefinite matrix related to  $x_i$ ,  $i = 1, \dots, n$  (see, for example, [29, Section 1.24]). As in [45], throughout this paper we assume that  $A_i$ 's are  $m \times m$  real symmetric positive semidefinite matrices and that there exists a  $w \in \Omega$  such that  $\mathcal{M}(w)$  is positive definite. This in particular implies that  $\mathcal{M}(w)$  is positive definite for all positive  $w \in \Omega$ . The optimal experimental design problem can then be formulated as the following minimization problem (see [29, Section 7.10]):

$$\begin{aligned} f^* &:= \inf_w \Phi(\mathcal{M}(w)) := \Psi(\mathcal{C}_K(\mathcal{M}(w))) \\ \text{s.t. } & w \in \Omega, \text{ Range}(K) \subseteq \text{Range}(\mathcal{M}(w)), \end{aligned} \tag{1}$$

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\*This work was supported in part by an NSERC Discovery Grant.

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where  $\text{Range}(X)$  denotes the range space of a matrix  $X$ ;  $\Psi$  is a function defined on the set of positive definite matrices; and  $\mathcal{C}_K(\mathcal{M}(w))$  is the information matrix defined by  $\mathcal{C}_K(\mathcal{M}(w)) := (K^T(\mathcal{M}(w))^\dagger K)^{-1}$ . Here  $A^\dagger$  denotes the Moore-Penrose pseudoinverse of a matrix  $A$ . The well-definedness of  $\mathcal{C}_K(\mathcal{M}(w))$  is guaranteed by the range inclusion condition in the constraint of (1) and the fact that  $K$  has full column rank [29, Chapter 3]. The range inclusion condition has been used in the literature and it is a required condition for estimability (see, for example, [29, Sections 1.18 and 3.4] for details). The function  $\Phi$  in the objective is commonly referred to as an “optimality criterion”. Some classical optimality criteria include (see [29, Chapter 6]):

- (i) A-criterion  $\Phi(X) := \text{tr}(K^T X^\dagger K)$ ;
- (ii) c-criterion  $\Phi(X) := c^T X^\dagger c$ ;
- (iii) D-criterion  $\Phi(X) := \log \det(K^T X^\dagger K)$ ;
- (iv)  $p$ th mean criterion  $\Phi(X) := \text{tr}((K^T X^\dagger K)^{-p})$ .

for some  $p < 0$ ,  $c \in \mathfrak{R}^m$  and  $K \in \mathfrak{R}^{m \times k}$  of full column rank.

It is easy to observe that c-criterion is just a special case of A-criterion with  $K = c$  and A-criterion is a special case of  $p$ th mean criterion with  $p = -1$ . We shall also mention that  $p$ th mean criterion can be defined more generally to include D-criterion as a special case (see [29, Chapter 6] for details). Furthermore, it can be shown that the constraint set of (1) is convex [29, Section 3.3], and the criteria (i)-(iv) are convex functions in the constraint set (by using [29, Theorem 5.14] and [29, Theorem 6.13], or [28, Proposition IV.14] and [28, Proposition IV.15]). Hence, problem (1) with these criteria is a convex optimization problem. Indeed, it is known that (1) with the above criteria can be reformulated as (possibly nonlinear) semidefinite programming (SDP) problems (see, for example, [14, 8, 10, 27]). Furthermore, problem (1) with c-criterion above can be reformulated and solved as a linear programming problem (see, for example, [19]).

The optimal design problems (1) with the aforementioned criteria usually do not have closed form solutions. Numerous procedures have thus been proposed to solve (1) (see, for example, [13, 44, 5, 6, 43, 7, 20, 9, 28, 3, 38, 1, 12, 30, 39, 32]). Among them, the multiplicative algorithm introduced in [33] has been widely explored. For example, Titterton [34], Pázman [28], Dette et al. [12] and Harman and Trnovská [21] studied the multiplicative algorithm for D-criterion. In addition, Fellman [15] and Torsney [37] considered the multiplicative algorithm for A-criterion under the assumption that all  $A_i$ 's are rank-one. Recently, Yu [45] studied the multiplicative algorithm for a class of convex optimality criteria and proved its global convergence under some assumptions. Nevertheless, for several commonly used optimality criteria, some of those assumptions may not hold and hence there is no theoretical guarantee for its convergence. Indeed, as observed in [45, Section 5], one of the assumptions does not hold for  $p$ th mean criterion with  $p = -2$ . Moreover, for such a criterion, our numerical experiments in Section 5 demonstrate that the multiplicative algorithm appears not to converge when  $p < -1$ . More details about the multiplicative algorithm for solving (1) are given in Section 2.

In this paper, we consider an alternative approach to solve problem (1). In particular, we propose an interior point (IP) method for (1) and establish its global convergence. The method is a Newton-type method that can be efficiently applied to solve problem (1) with a broad class of convex optimality criteria and moderate-sized matrices  $A_i$ 's. By exploiting the structure of the Hessian matrix of these optimality criteria, we derive an explicit formula for its rank. This formula applies to the classical A-, D- and  $p$ th mean criterion. Using this result, we further show that the Newton direction arising in the IP method for (1) with the aforementioned classical optimality criteria can be computed efficiently via Sherman-Morrison-Woodbury formula when  $n \gg m^2$ , i.e., when the size of  $A_i$ 's is small relative to the size of design space. We finally compare the IP method with the multiplicative algorithm and also with the standard IP solver SDPT3 [36, 40] which solves linear/log-determinant SDP reformulations of the problems via a primal-dual IP method. The computational results show that our IP method usually outperforms these two methods in both speed and solution quality.

The rest of this paper is organized as follows. In Subsection 1.1, we introduce the notations that are used throughout the paper. In Section 2, we review the multiplicative algorithm and address its convergence. In Section 3, we propose an IP method for solving problem (1) with a large class of convex optimality criteria and address its convergence. In Section 4, we discuss how the IP method can be applied to solve problem (1) with criteria (i)–(iv) and demonstrate how the Newton direction can be computed efficiently when  $n \gg m^2$ . In Section 5, we conduct numerical experiments to test the performance of the method and compare it with the multiplicative algorithm and the standard IP solver SDPT3. Finally, we present some concluding remarks in Section 6.

## 1.1 Notations

In this paper, the symbol  $\mathfrak{R}_{++}$  denotes the set of all positive real numbers and  $\mathfrak{R}^n$  denotes the  $n$ -dimensional Euclidean space. For a vector  $x \in \mathfrak{R}^n$ ,  $\|x\|$  denotes the Euclidean norm of  $x$  and  $\mathcal{D}(x)$  denotes the diagonal matrix whose  $i$ th diagonal entry is  $x_i$  for all  $i$ . For  $\alpha \in \mathfrak{R}$  and a vector  $x \in \mathfrak{R}^n$  with positive entries,  $x^\alpha$  denotes the vector whose  $i$ th entry is  $x_i^\alpha$  for all  $i$ . For  $x, y \in \mathfrak{R}^n$ ,  $x \circ y$  denotes the Hadamard (entry-wise) product of  $x$  and  $y$ . The letter  $e$  denotes the vector of all ones, whose dimension should be clear from the context. The set of all  $m \times n$  matrices with real entries is denoted by  $\mathfrak{R}^{m \times n}$ . Given  $A, B \in \mathfrak{R}^{m \times n}$ ,  $A \otimes B$  denotes the Kronecker product of  $A$  and  $B$ , while  $A \circ B$  denotes the Hadamard (entry-wise) product of  $A$  and  $B$ . The trace of a real square matrix  $A$  is denoted by  $\text{tr}(A)$ . We denote by  $I$  the identity matrix, whose dimension should be clear from the context.

The space of  $n \times n$  symmetric matrices is denoted by  $\mathcal{S}^n$ . If  $A \in \mathcal{S}^n$  is positive semidefinite (resp., definite), we write  $A \succeq 0$  (resp.,  $A \succ 0$ ). The cone of positive semidefinite (resp., definite) matrices is denoted by  $\mathcal{S}_+^n$  (resp.,  $\mathcal{S}_{++}^n$ ). For  $A, B \in \mathcal{S}^n$ ,  $A \succeq B$  (resp.,  $A \succ B$ ) means  $A - B \succeq 0$  (resp.,  $A - B \succ 0$ ). For any  $A \in \mathcal{S}^n$ , we define the vectors  $\mathbf{svec}(A) \in \mathfrak{R}^{n(n+1)/2}$  and  $\mathbf{svec}_0(A) \in \mathfrak{R}^{n(n+1)/2}$  as

$$\begin{aligned} \mathbf{svec}(A) &= (a_{11}, \sqrt{2}a_{21}, \dots, \sqrt{2}a_{n1}, a_{22}, \sqrt{2}a_{32}, \dots, \sqrt{2}a_{n2}, \dots, a_{nn})^T, \\ \mathbf{svec}_0(A) &= (a_{11}, a_{21}, \dots, a_{n1}, a_{22}, a_{32}, \dots, a_{n2}, \dots, a_{nn})^T. \end{aligned}$$

Notice that  $\mathbf{svec}$  is an isometry between  $\mathcal{S}^n$  and  $\mathfrak{R}^{n(n+1)/2}$  and moreover,

$$\text{tr}(AB) = \mathbf{svec}(A)^T \mathbf{svec}(B) \quad \forall A, B \in \mathcal{S}^n. \quad (2)$$

We denote the inverse map of  $\mathbf{svec}$  by  $\mathbf{smat}$ . Clearly, they are adjoint of each other, that is,

$$a^T \mathbf{svec}(B) = \text{tr}(\mathbf{smat}(a)B) \quad \forall a \in \mathfrak{R}^{n(n+1)/2}, B \in \mathcal{S}^n.$$

The symmetric Kronecker product of any two (not necessarily symmetric) matrices  $C, D \in \mathfrak{R}^{n \times n}$  is a square matrix of order  $n(n+1)/2$  such that

$$(C \otimes_s D) \mathbf{svec}(A) = \frac{1}{2} \mathbf{svec}(CAD^T + DAC^T) \quad \forall A \in \mathcal{S}^n. \quad (3)$$

As mentioned in [35],  $C \otimes_s D$  can be expressed in terms of the standard Kronecker product of  $C$  and  $D$  as follows:

$$C \otimes_s D = \frac{1}{2} Q(C \otimes D + D \otimes C)Q^T,$$

where  $Q \in \mathfrak{R}^{n(n+1)/2 \times n^2}$  is such that

$$Q \mathbf{vec}(A) = \mathbf{svec}(A), \quad Q^T \mathbf{svec}(A) = \mathbf{vec}(A) \quad \forall A \in \mathcal{S}^n. \quad (4)$$

It is easy to observe that the above  $Q$  exists and is unique. Moreover,  $QQ^T = I$ .

A function  $f : \mathcal{S}^n \rightarrow \mathfrak{R}$  is said to be increasing (resp., decreasing) if for any  $A \succeq B$ , it holds that

$$f(A) \geq f(B) \quad (\text{resp.}, f(A) \leq f(B)).$$

## 2 The multiplicative algorithm

In this section we review the multiplicative algorithm introduced in [33] for solving problem (1) and discuss its convergence. In particular, we first describe the multiplicative algorithm as follows, which is specified through a power parameter  $\lambda \in (0, 1]$ .

### Multiplicative Algorithm:

1. **Start:** Let a positive  $w^0 \in \Omega$  and  $\lambda \in (0, 1]$  be given.
2. **For**  $t = 0, 1, \dots$

$$w_i^{t+1} = w_i^t \frac{(d_i(w^t))^\lambda}{\sum_{j=1}^n w_j^t (d_j(w^t))^\lambda}, \quad i = 1, \dots, n, \quad (5)$$

where  $d_i(w) = -\text{tr}(\nabla\Phi(\mathcal{M}(w))A_i)$  and  $\nabla\Phi(\mathcal{M}(w))$  is the gradient of  $\Phi$  at  $\mathcal{M}(w)$ .

**End** (for)

**Remark 2.1.** *The above algorithm is the same as the one described in [45], in the sense that both algorithms generate exactly the same sequence  $\{w^t\}$  provided the initial points  $w^0$  are identical. ■*

We now state a global convergence result recently established by Yu [45, Theorem 2] for the multiplicative algorithm when applied to solve the following problem, which is closely related to (1):

$$\begin{aligned} \text{val} := \sup_w & -\Phi(\mathcal{M}(w)) \\ \text{s.t.} & w \in \Omega, \mathcal{M}(w) \succ 0. \end{aligned} \quad (6)$$

Observe that (1) and (6) are equivalent (i.e., the optimal value being negative of each other) if there exists an optimal solution  $w^*$  of (1) with  $\mathcal{M}(w^*) \succ 0$ , or if  $\Phi$  is convex in

$$\mathcal{S}_+^m(K) := \{X \in \mathcal{S}_+^m : \text{Range}(K) \subseteq \text{Range}(X)\}.$$

and (6) has an optimal solution.

**Proposition 2.1.** *Let  $\{w^t\}$  be the sequence generated from the above multiplicative algorithm. Suppose the following assumptions hold:*

- (a) *for any feasible point  $w$  of (6),  $\nabla\Phi(\mathcal{M}(w)) \preceq 0$  and  $\nabla\Phi(\mathcal{M}(w))A_i \neq 0$  for  $i = 1, \dots, n$ ;*
- (b) *for any feasible point  $w$  of (6), if  $T(w) \neq w$ , then  $\Phi(\mathcal{M}(T(w))) < \Phi(\mathcal{M}(w))$ , where*

$$[T(w)]_i := w_i \frac{(d_i(w))^\lambda}{\sum_{j=1}^n w_j (d_j(w))^\lambda}, \quad i = 1, \dots, n;$$

- (c)  *$\Phi$  is strictly convex and  $\nabla\Phi$  is continuous in  $\mathcal{S}_{++}^m$ ;*
- (d) *for any  $\{X^t\} \subset \mathcal{S}_{++}^m$ , if  $X^t \rightarrow X^*$  and  $\{\Phi(X^t)\}$  is decreasing, then  $X^* \succ 0$ .*

*Then  $\Phi(\mathcal{M}(w^t)) \rightarrow -\text{val}$  monotonically, and moreover, any accumulation point of  $\{w^t\}$  is an optimal solution of (6).*

**Remark 2.2.** *Notice that the assumptions in the above proposition imply that any accumulation point  $w^*$  of  $\{w^t\}$  satisfies  $\mathcal{M}(w^*) \succ 0$ . Hence, if the assumptions in Proposition 2.1 hold and  $\Phi$  is convex in  $\mathcal{S}_+^m(K)$ , then (1) is equivalent to (6) and any accumulation point of the sequence  $\{w^t\}$  generated from the above multiplicative algorithm solves (1). ■*

Using Proposition 2.1 and some technical results developed in [45], one can establish the convergence of the above multiplicative algorithm when applied to problem (1) with A-, D- and  $p$ th mean criterion for  $p \in (-1, 0)$  and  $K = I$ , which is summarized as follows.

**Corollary 2.1.** *Assume that  $K = I$  and  $A_i \neq 0$  for  $i = 1, \dots, n$ . Then the multiplicative algorithm converges for any  $\lambda \in (0, 1]$  when applied to problem (1) with D- and  $p$ th mean criterion for  $p \in (-1, 0)$ . Also, it converges for A-criterion when  $\lambda \in (0, 1)$ .*

As seen from Proposition 2.1 and Corollary 2.1, the multiplicative algorithm converges for a large class of optimality criteria  $\Phi$ . Nevertheless, for some interesting convex optimality criteria, the assumptions stated in Proposition 2.1 may not hold and hence there is no theoretical guarantee for its convergence. Indeed, as observed in [45, Section 5], the assumption (b) with  $\lambda = 1$  does not hold for  $p$ th mean criterion with  $p = -2$ . Moreover, for such a criterion, our numerical experiments in Section 5 demonstrate that the multiplicative algorithm appears not to converge when  $p < -1$ .

Due to the aforementioned potential drawbacks of the multiplicative algorithm, we will propose an IP method for solving problem (1) with a broad class of optimality criteria  $\Phi$  including A-, D- and  $p$ th mean criterion in subsequent sections.

### 3 IP method for a class of convex optimality criteria

In this section, we propose an IP method for solving (1) with a class of convex optimality criteria  $\Phi = \Psi \circ \mathcal{C}_K$ , where  $\circ$  denotes composition of functions. We make the following assumption on  $\Psi$  throughout this paper.

**Assumption 3.1.** *The function  $\Psi$  is convex, decreasing, twice continuously differentiable and bounded below on any bounded subset of  $\mathcal{S}_{++}^k$ . Moreover, for any bounded sequence  $\{X^t\} \subseteq \mathcal{S}_{++}^k$  with  $\lambda_{\min}(X^t) \rightarrow 0$ , one has  $\Psi(X^t) \rightarrow \infty$ .*

**Remark 3.1.** *We now make some brief comments on the above assumptions.*

- (a) *Assumption 3.1 is fairly reasonable. Indeed, all optimality criteria described in Section 1 satisfy this assumption.*
- (b) *Since the feasible set is not necessarily closed, problem (1) with a general convex optimality criterion may not have an optimal solution. However, when the optimality criterion satisfies Assumption 3.1, it must have an optimal solution as shown in Theorem 3.1(a). We refer the readers to [29, Chapter 5] for more discussion on conditions guaranteeing existence of solutions for problem (1).*
- (c) *In contrast to Proposition 2.1, we do not require the existence of a positive definite optimal moment matrix  $\mathcal{M}(w^*)$ . Indeed, Assumption 3.1 may hold even when problem (1) does not have a positive definite optimal moment matrix. For instance, the design problem*

$$\begin{aligned} \min_{w, X} \quad & \begin{pmatrix} 1 \\ 0 \end{pmatrix}^T X^\dagger \begin{pmatrix} 1 \\ 0 \end{pmatrix} \\ \text{s.t.} \quad & X = \begin{pmatrix} w_1 & 0 \\ 0 & w_2 \end{pmatrix}, w_1 + w_2 = 1, w_1, w_2 \geq 0, \\ & \begin{pmatrix} 1 \\ 0 \end{pmatrix} \in \text{Range}(X), \end{aligned}$$

*has a unique optimal solution at  $(w_1, w_2) = (1, 0)$ . The corresponding optimal moment matrix is not positive definite; thus, the assumption (d) of Proposition 2.1 does not hold. However, it is easy to check that Assumption 3.1 is satisfied for this design problem (with  $\Psi(t) = 1/t$ ). In general, the assumption (d) of Proposition 2.1 is likely not satisfied when  $K$  is not invertible, while our Assumption 3.1 is independent of  $K$ .*

■

Under Assumption 3.1, it is not hard to show that the function  $\Phi(\mathcal{M}(\cdot))$  is bounded below on the feasible set of (1). Also, it is routine to show that the function  $\Phi$  is twice continuously differentiable

in  $\mathcal{S}_{++}^m$ . Furthermore, it can be shown that  $\Phi$  is convex in  $\mathcal{S}_+^m(K)$  by considering suitable Schur complements (see, for example, [27, Section 6]). We include a short proof below for the convenience of readers. Before proceeding, we state the following well-known fact, which concerns the Schur complement of a positive semidefinite submatrix (see, for example, [29, Lemma 3.12]).

**Lemma 3.1.** *Let  $A \in \mathcal{S}^k$ ,  $B \in \mathfrak{R}^{m \times k}$  and  $C \in \mathcal{S}^m$ . Then the matrix  $\begin{pmatrix} A & B^T \\ B & C \end{pmatrix}$  is positive semidefinite if and only if  $A \succeq B^T C^\dagger B$ ,  $C \succeq 0$  and  $\text{Range}(B) \subseteq \text{Range}(C)$ .*

**Proposition 3.1.** *The optimality criterion  $\Phi$  is convex in  $\mathcal{S}_+^m(K)$ .*

*Proof.* First of all, it can be shown that the set  $\mathcal{S}_+^m(K)$  is convex (see, for example, [29, Section 3.3]). In addition, notice that for any  $X \in \mathcal{S}_+^m(K)$ , we have

$$\begin{aligned} \Phi(X) &= \Psi((K^T X^\dagger K)^{-1}) = \inf_U \{ \Psi(U) : (K^T X^\dagger K)^{-1} \succeq U \succ 0 \} \\ &= \inf_U \{ \Psi(U) : U^{-1} \succeq K^T X^\dagger K, U \succ 0 \} \\ &= \inf_U \left\{ \Psi(U) : \begin{pmatrix} U^{-1} & K^T \\ K & X \end{pmatrix} \succeq 0, U \succ 0 \right\} \\ &= \inf_U \{ \Psi(U) : X \succeq K U K^T, U \succ 0 \}, \end{aligned} \tag{7}$$

where the second equality follows from the fact that  $\Psi$  is decreasing, the fourth and last equalities follow from Lemma 3.1, while the third equality holds because  $K^T X^\dagger K$  is invertible for  $X \in \mathcal{S}_+^m(K)$  when  $K$  has full column rank. Convexity of  $\Phi$  in  $\mathcal{S}_+^m(K)$  now follows from [31, Theorem 5.7]. ■

Observe that  $\mathcal{M}(w) \succ 0$  whenever  $w > 0$ . Thus, under Assumption 3.1, the function  $\Phi$  is twice continuously differentiable for any positive  $w \in \Omega$ . It is hence natural to develop an IP method to solve (1) since such a method keeps all iterates in the relative interior of  $\Omega$  until convergence. To proceed, we first reformulate the problem by eliminating the equality constraint. The resulting equivalent problem is given by

$$\begin{aligned} f^* &= \inf_{\tilde{w}} f(\tilde{w}) := \Phi(\mathcal{M}(P\tilde{w} + q)) \\ \text{s.t.} & \quad e^T \tilde{w} \leq 1, \tilde{w} \geq 0, \\ & \quad \text{Range}(K) \subseteq \text{Range}(\mathcal{M}(P\tilde{w} + q)), \end{aligned} \tag{8}$$

where  $P \in \mathfrak{R}^{n \times (n-1)}$  and  $q \in \mathfrak{R}^n$  are such that

$$P\tilde{w} + q = \begin{pmatrix} \tilde{w} \\ 1 - e^T \tilde{w} \end{pmatrix} \quad \forall \tilde{w} \in \mathfrak{R}^{n-1}. \tag{9}$$

We next develop an IP method for solving problem (8) instead. First, we need to build a suitable barrier function. Given any  $\tilde{w} > 0$  satisfying  $e^T \tilde{w} < 1$ , one can observe that  $P\tilde{w} + q > 0$  and hence  $\mathcal{M}(P\tilde{w} + q) \succ 0$ , which leads to  $\text{Range}(K) \subseteq \text{Range}(\mathcal{M}(P\tilde{w} + q))$ . This implies that any barrier function that takes into account the first two inequality constraints of (8) is sufficient for the development of an IP method. Here we naturally choose the logarithmic barrier function and then solve the barrier subproblem in the form of

$$\min_{\tilde{w}} f_\mu(\tilde{w}) := f(\tilde{w}) - \mu \sum_{i=1}^{n-1} \log(\tilde{w}_i) - \mu \log(1 - e^T \tilde{w}) \tag{10}$$

for a sequence of parameters  $\mu \downarrow 0$ . In view of Assumption 3.1, we see that any level set of  $f_\mu$  is compact. Moreover,  $f_\mu$  is strictly convex. Thus, there exists a unique minimizer to (10) for any  $\mu > 0$ . Furthermore, it follows from Assumption 3.1 that  $f_\mu$  is twice continuously differentiable and its Hessian is positive definite in its domain. Therefore, problem (10) can be suitably solved by the Newton's method with a line search whose stepsize is chosen by Armijo rule.

We are now ready to present our IP method for solving problem (8).

**IP Method:**

1. **Start:** Let a strictly feasible  $\tilde{w}^0$ ,  $0 < \beta, \gamma, \eta, \sigma < 1$  and  $\mu_1 > 0$  be given. Let  $\epsilon(\mu)$  be an increasing function of  $\mu$  so that  $\lim_{\mu \downarrow 0} \epsilon(\mu) = 0$ . Set  $\tilde{w} = \tilde{w}^0$  and  $t = 1$ .

2. **While**  $\|\nabla f_{\mu_t}(\tilde{w})\| > \epsilon(\mu_t)$  **do**

(a) Compute the Newton direction

$$d := -(\nabla^2 f_{\mu_t}(\tilde{w}))^{-1} \nabla f_{\mu_t}(\tilde{w}). \quad (11)$$

(b) Let  $\alpha_{\max}(\tilde{w}) := \max\{\alpha : \tilde{w}[\alpha] \geq 0, e^T \tilde{w}[\alpha] \leq 1\}$ , where  $\tilde{w}[\alpha] := \tilde{w} + \alpha d$ .

(c) Let  $\alpha$  be the largest element of  $\{\bar{\alpha}(\tilde{w}), \beta \bar{\alpha}(\tilde{w}), \beta^2 \bar{\alpha}(\tilde{w}), \dots\}$  satisfying

$$f_{\mu_t}(\tilde{w}[\alpha]) \leq f_{\mu_t}(\tilde{w}) + \sigma \alpha (\nabla f_{\mu_t}(\tilde{w}))^T d,$$

where  $\bar{\alpha}(\tilde{w}) := \min\{1, \eta \alpha_{\max}(\tilde{w})\}$ .

(d) Set  $\tilde{w} = \tilde{w}[\alpha]$ .

**End** (while)

3. Set  $\tilde{w}^t = \tilde{w}$ ,  $\mu_{t+1} = \gamma \mu_t$ ,  $t = t + 1$ , and go to step 2.

In standard convergence analysis of IP methods, the feasible set of the problem is usually assumed to be closed and the objective function is at least continuous on the feasible set (see, for example, [17, Section 3.3], [26, Theorem 19.1] and [18, Section 16.2.3]). An alternative setting was considered in [25, Section 3.2], where the feasible set is assumed to be closed and the objective function is assumed to be  $\beta$ -compatible with the associated barrier function. While these assumptions are quite general, they do not necessarily hold for our problem (8). In particular, the feasible set of (8) is not necessarily closed and moreover it is hard to verify whether the second assumption imposed in [25] holds for (8). The existing convergence results of IP methods in the literature [25, 17, 26, 18] are thus not directly applicable to our above IP method.

The convergence analysis of our IP method is carried out as follows. We first study the convergence of its outer iterations and then discuss the convergence of its inner iterations. The analysis on the outer iterations substantially uses the specific structure of  $\Phi$ , namely, its equivalent reformulation (7) and Assumption 3.1, which enable us to establish the existence of an optimal solution of (8) and show that any accumulation point of the sequence generated from our IP method is an optimal solution of (8).

For notational convenience, in the remainder of this section, we associate with each  $\tilde{w} \in \mathfrak{R}^{n-1}$  a unique  $w \in \mathfrak{R}^n$  by letting  $w := P\tilde{w} + q$ . Analogously, we associate with each  $w \in \mathfrak{R}^n$  a unique  $\tilde{w} \in \mathfrak{R}^{n-1}$  by letting  $\tilde{w}_i = w_i$  for  $i = 1, \dots, n-1$ . Also, we let  $\Phi_{\mathcal{M}}(w) := \Phi(\mathcal{M}(w))$ .

We first observe that if problem (1) has an optimal solution  $w^*$  with  $\mathcal{M}(w^*) \succ 0$ , then there exists a Lagrange multiplier  $u^* \geq 0$  such that  $(w^*, u^*)$  satisfies the following KKT system:

$$\begin{aligned} P^T(\nabla \Phi_{\mathcal{M}}(w) - u) &= 0, \\ e^T w &= 1, \\ u \circ w &= 0, \\ (w, u) &\geq 0. \end{aligned} \quad (12)$$

Given a strictly feasible point  $\tilde{w} \in \mathfrak{R}^{n-1}$  of problem (10), we notice that

$$\nabla f_{\mu}(\tilde{w}) = P^T(\nabla \Phi_{\mathcal{M}}(w) - \mu w^{-1}). \quad (13)$$

Then it is not hard to observe that for each  $\mu > 0$ , the  $w$  associated with the approximate solution  $\tilde{w}$  of (10) obtained by the Newton's method detailed in step 2 above together with  $u := \mu w^{-1}$  satisfies the following perturbed KKT system,

$$\begin{aligned} P^T(\nabla\Phi_{\mathcal{M}}(w) - u) &= v, \\ e^T w &= 1, \\ u \circ w &= \mu e, \\ (w, u) &> 0 \end{aligned} \tag{14}$$

with  $v = \nabla f_{\mu}(\tilde{w}) \in \mathfrak{R}^{n-1}$  whose norm  $\|v\| \leq \epsilon(\mu)$ . Thus, the convergence regarding the outer iterations of our IP method is related to the limiting behavior of the solutions of system (14) as  $(\mu, v) \rightarrow (0_+, 0)$ , that is,  $(\mu, v) \rightarrow (0, 0)$  with  $\mu > 0$ .

We first claim that system (14) has a unique solution for any  $(\mu, v) \in \mathfrak{R}_{++} \times \mathfrak{R}^{n-1}$ . Indeed, it is easy to observe that  $(w, u)$  is a solution of (14) if and only if  $\tilde{w} \in \mathfrak{R}^{n-1}$  is an optimal solution of

$$\min_{\tilde{w}} f_{\mu}(\tilde{w}) - v^T \tilde{w}. \tag{15}$$

Since the objective function of (15) is strictly convex and it has compact level sets, problem (15) has a unique optimal solution, which immediately implies that system (14) has a unique solution. From now on, we denote by  $(w(\mu, v), u(\mu, v))$  the unique solution of (14). Our main theorem below discusses the limiting behavior of  $(w(\mu, v), u(\mu, v))$  as  $(\mu, v) \rightarrow (0_+, 0)$ . The proof of this theorem can be found in the appendix.

**Theorem 3.1.** *Let  $(w(\mu, v), u(\mu, v))$  be defined above for  $(\mu, v) \in \mathfrak{R}_{++} \times \mathfrak{R}^{n-1}$ . Then the following statements hold:*

- (a)  $\lim_{(\mu, v) \rightarrow (0_+, 0)} \Phi(\mathcal{M}(w(\mu, v))) = f^*$  and any accumulation point of  $w(\mu, v)$  as  $(\mu, v) \rightarrow (0_+, 0)$  is an optimal solution of (1).
- (b) Suppose in addition that problem (1) has an optimal solution  $w^*$  with  $\mathcal{M}(w^*) \succ 0$ . Then any accumulation point of  $w(\mu, v)$  as  $(\mu, v) \xrightarrow{\Xi_C} (0, 0)$ , i.e.,  $(\mu, v) \rightarrow (0, 0)$  with  $(\mu, v) \in \Xi_C := \{(\mu, v) : \|v\|_{\infty} < C\mu\}$  for some given  $C > 0$ , is an optimal solution of (1) with maximum cardinality.

As an immediate consequence of Theorem 3.1, we have the following global convergence result regarding the outer iterations of our IP method, whose simple proof is omitted.

**Corollary 3.1.** *Let  $\{\mu_t\}$  and  $\{\tilde{w}^t\}$  be the sequences generated in the IP method. Let  $w^t = P\tilde{w}^t + q$  for all  $t$ . Then the following statements hold:*

- (a)  $\lim_{t \rightarrow \infty} \Phi(\mathcal{M}(w^t)) = f^*$  and any accumulation point of  $\{w^t\}$  is an optimal solution of (1).
- (b) Suppose in addition that problem (1) has an optimal solution  $w^*$  with  $\mathcal{M}(w^*) \succ 0$  and  $\epsilon(\mu_t) = O(\mu_t)$ . Then any accumulation point of  $\{w^t\}$  is an optimal solution of (1) with maximum cardinality.

We emphasize that in Corollary 3.1 (a), we do *not* require the existence of an optimal solution  $w^*$  with  $\mathcal{M}(w^*) \succ 0$ . On the other hand, if such an optimal solution does exist, for example, when  $K = I$ , then Corollary 3.1 (b) states that the accumulation point (with  $\epsilon(\mu_t) = O(\mu_t)$ ) must be an optimal solution of (1) that has the largest number of non-zero entries among all the optimal solutions of (1).

Before ending this section, we establish a convergence result regarding the inner iterations of our IP method.

**Proposition 3.2.** *Let  $\mu_t > 0$  and  $\epsilon(\mu_t) > 0$  be given. Then the Newton's method detailed in step 2 of the IP method starting from any strictly feasible point  $\tilde{w}^{\text{init}}$  of (8) generates a point  $\tilde{w}^t$  satisfying  $\|\nabla f_{\mu_t}(\tilde{w}^t)\| \leq \epsilon(\mu_t)$  within a finite number of iterations.*

*Proof.* First, observe that all iterates generated by the Newton's method lie in the compact level set  $\Upsilon := \{\tilde{w} : f_{\mu_t}(\tilde{w}) \leq f_{\mu_t}(\tilde{w}^{\text{init}})\}$ . Furthermore, it holds that  $\tilde{w} > 0$  and  $1 - e^T \tilde{w} > 0$  for all  $\tilde{w} \in \Upsilon$ . This together with the assumption that  $\mathcal{M}(\Omega) \cap \mathcal{S}_{++}^m \neq \emptyset$  implies that  $\mathcal{M}(\Upsilon) \subset \mathcal{S}_{++}^m$ . Thus  $\nabla f_{\mu_t}$  and  $\nabla^2 f_{\mu_t}$  are continuous in  $\Upsilon$ . Using this observation and the strong convexity of  $f_{\mu_t}$  in  $\Upsilon$ , there exist  $\underline{\lambda}, \bar{\lambda} > 0$  such that  $\underline{\lambda}I \preceq \nabla^2 f_{\mu_t}(\tilde{w}) \preceq \bar{\lambda}I$  for all  $\tilde{w} \in \Upsilon$ . This relation along with the continuity of  $\nabla f_{\mu_t}$  and  $\nabla^2 f_{\mu_t}$  implies that  $d = -(\nabla^2 f_{\mu_t}(\tilde{w}))^{-1} \nabla f_{\mu_t}(\tilde{w})$  is continuous in  $\Upsilon$ . In view of this result and the definition of  $\bar{\alpha}(\tilde{w})$ , it is not hard to show that  $\bar{\alpha}(\tilde{w})$  is positive and continuous in  $\Upsilon$ . This fact together with the compactness of  $\Upsilon$  yields  $\underline{\alpha} := \inf\{\bar{\alpha}(\tilde{w}) : \tilde{w} \in \Upsilon\} > 0$ . Thus, all iterates  $\tilde{w}$  generated by the Newton's method satisfy  $\underline{\lambda}I \preceq \nabla^2 f_{\mu_t}(\tilde{w}) \preceq \bar{\lambda}I$  and  $\bar{\alpha}(\tilde{w}) \in [\underline{\alpha}, 1]$ . The remaining proof follows the same arguments as in the proof of [24, Theorem 3.13].  $\blacksquare$

## 4 IP method for classical optimality criteria

In this section, we discuss how to apply our IP method to solve problem (1) with A-, D- and  $p$ th mean criterion. In particular, we will demonstrate how the Newton direction (11) can be efficiently computed for each criterion.

First, for each optimality criterion  $\Phi$  and the corresponding function  $\Psi$ , we define the associated functions  $\phi$  and  $\psi$  as follows:

$$\phi(x) = \Phi(\mathbf{smat}(x)), \quad \psi(y) = \Psi(\mathbf{smat}(y)) \quad (16)$$

for any  $x \in \mathfrak{R}^{m(m+1)/2}$  and  $y \in \mathfrak{R}^{k(k+1)/2}$ , provided that  $\Phi(\mathbf{smat}(x))$  and  $\Psi(\mathbf{smat}(y))$  are well-defined. It is clear that  $\phi$  and  $\psi$  are convex due to the convexity of  $\Phi$  and  $\Psi$ , respectively. Define

$$M := [\mathbf{svec}(A_1) \dots \mathbf{svec}(A_n)].$$

Clearly,  $M \in \mathfrak{R}^{m(m+1)/2 \times n}$ .

With the notations above, the function  $f_{\mu}$  defined in (10) can be rewritten as

$$f_{\mu}(\tilde{w}) = \phi(M(P\tilde{w} + q)) - \mu \sum_{i=1}^{n-1} \log(\tilde{w}_i) - \mu \log(1 - e^T \tilde{w}).$$

By the chain rule, the gradient and Hessian of  $f_{\mu}$  are given by

$$\begin{aligned} \nabla f_{\mu}(\tilde{w}) &= P^T M^T \nabla \phi(Mw) - \mu P^T w^{-1}, \\ \nabla^2 f_{\mu}(\tilde{w}) &= P^T M^T \nabla^2 \phi(Mw) M P + \frac{\mu}{(1 - e^T \tilde{w})^2} e e^T + \mu \mathcal{D}(\tilde{w}^{-2}), \end{aligned} \quad (17)$$

where  $w = P\tilde{w} + q$ .

The main computational effort of our IP method lies in computing the Newton direction  $d$  by solving the system  $\nabla^2 f_{\mu}(\tilde{w})d = -\nabla f_{\mu}(\tilde{w})$  (see (11)). In applications,  $n$  can be significantly larger than  $m^2$ . Since the rank of  $\nabla^2 \phi(Mw)$  is at most  $m(m+1)/2$ , the first matrix in (17) has ‘‘low’’ rank compared to  $\nabla^2 f_{\mu}(\tilde{w})$ . It is generally more efficient to compute the Newton direction via the Sherman-Morrison-Woodbury formula, *without* explicitly forming the Hessian matrix. To this end, suppose that  $\nabla^2 \phi(Mw)$  has rank  $r$ . Let  $V D V^T$  be the *partial* eigenvalue decomposition of  $\nabla^2 \phi(Mw)$ , where  $D$  is the  $r \times r$  diagonal matrix whose diagonal consists of the  $r$  largest eigenvalues of  $\nabla^2 \phi(Mw)$ , and the columns of  $V$  are the corresponding eigenvectors.<sup>1</sup> Due to the convexity of  $\phi$ , one can observe that  $\nabla^2 \phi(Mw) = V D V^T$ . It then follows from (17) that

$$\begin{aligned} \nabla^2 f_{\mu}(\tilde{w}) &= (P^T M^T V) D (V^T M P) + \frac{\mu}{(1 - e^T \tilde{w})^2} e e^T + \mu \mathcal{D}(\tilde{w}^{-2}) \\ &= \begin{pmatrix} P^T M^T V & e \end{pmatrix} \begin{pmatrix} D & 0 \\ 0 & \frac{\mu}{(1 - e^T \tilde{w})^2} \end{pmatrix} \begin{pmatrix} V^T M P \\ e^T \end{pmatrix} + \mu \mathcal{D}(\tilde{w}^{-2}), \end{aligned}$$

<sup>1</sup>The partial eigenvalue decomposition can be efficiently computed by the package PROPACK [23].

which together with the Sherman-Morrison-Woodbury formula yields the Newton direction

$$d = -(\nabla^2 f_\mu(\tilde{w}))^{-1} \nabla f_\mu(\tilde{w}) = -\left[ \frac{1}{\mu} \mathcal{D}(\tilde{w}^2) - \frac{1}{\mu^2} \mathcal{D}(\tilde{w}^2) (P^T M^T V \quad e) W \begin{pmatrix} V^T M P \\ e^T \end{pmatrix} \mathcal{D}(\tilde{w}^2) \right] \nabla f_\mu(\tilde{w}),$$

where

$$W = \left( \begin{pmatrix} D^{-1} & 0 \\ 0 & \frac{(1-e^T \tilde{w})^2}{\mu} \end{pmatrix} + \frac{1}{\mu} \begin{pmatrix} V^T M P \\ e^T \end{pmatrix} \mathcal{D}(\tilde{w}^2) (P^T M^T V \quad e) \right)^{-1}.$$

When  $n \gg m^2$ , the above approach is much more efficient than solving the Newton system directly by performing Cholesky factorization of  $\nabla^2 f_\mu(\tilde{w})$ . We remark that the ideas of using Sherman-Morrison-Woodbury formula to solve specially structured Newton systems have been explored in the literature (see, for example, [2, 16]).

As seen from above, both  $\nabla \phi(Mw)$  and  $\nabla^2 \phi(Mw)$  are needed for computing the Newton direction. In addition, the rank of  $\nabla^2 \phi(Mw)$  is used in the partial eigenvalue decomposition of  $\nabla^2 \phi(Mw)$ . Since this Hessian could become ill-conditioned as  $\mu \downarrow 0$ , it is more desirable to explicitly determine the rank of  $\nabla^2 \phi(Mw)$  a priori than to compute the rank numerically in each iteration. We will address these two issues in the next theorem. In particular, we derive generic formulas for evaluating  $\nabla \phi$  and  $\nabla^2 \phi$  that will be specialized for the A, D and  $p$ th mean criterion in subsequent subsections. We also determine the rank of  $\nabla^2 \phi$  under some suitable assumptions on  $\Psi$ .

**Theorem 4.1.** *Consider the function  $\Phi(X) = \Psi(\mathcal{C}_K(X))$  with  $\mathcal{C}_K(X) = (K^T X^\dagger K)^{-1}$  and  $\Psi$  satisfying Assumption 3.1. Let  $\phi$  and  $\psi$  be the associated functions of  $\Phi$  and  $\Psi$  as defined in (16). Let  $Q_1 \in \mathfrak{R}^{m(m+1)/2 \times m^2}$  and  $Q_2 \in \mathfrak{R}^{k(k+1)/2 \times k^2}$  be defined in (4), in place of  $Q$ , for the spaces  $\mathcal{S}^m$  and  $\mathcal{S}^k$ , respectively. Then the gradient and Hessian of  $\phi$  at any  $x \in \mathbf{svec}(\mathcal{S}_{++}^m)$  are given by*

$$\nabla \phi(x) = \mathbf{svec}(X^{-1} K \mathcal{C}_K(X) \nabla \Psi(\mathcal{C}_K(X)) \mathcal{C}_K(X) K^T X^{-1}), \quad (18)$$

$$\nabla^2 \phi(x) = Q_1 (G_0 \otimes G_0) Q_2^T \nabla^2 \psi(\mathbf{svec}(\mathcal{C}_K(X))) Q_2 (G_0 \otimes G_0)^T Q_1^T - 2X^{-1} \otimes_s G_2 + 2G_1 \otimes_s G_2, \quad (19)$$

where  $X = \mathbf{smat}(x)$ ,  $G_0 := X^{-1} K \mathcal{C}_K(X)$ ,  $G_1 := X^{-1} K \mathcal{C}_K(X) K^T X^{-1}$ ,  $G_2 := G_0 \nabla \Psi(\mathcal{C}_K(X)) G_0^T$ . In addition, if we suppose that  $-\nabla \Psi(\mathcal{C}_K(X))$  and  $\nabla^2 \psi(\mathbf{svec}(\mathcal{C}_K(X)))$  are positive definite matrices, then the rank of  $\nabla^2 \phi(x)$  is  $m(m+1)/2 - (m-k)(m-k+1)/2$ , which is independent of  $x \in \mathbf{svec}(\mathcal{S}_{++}^m)$ .

*Proof.* To derive the gradient of  $\phi$ , we fix an arbitrary  $x \in \mathbf{svec}(\mathcal{S}_{++}^m)$ . Let  $X = \mathbf{smat}(x)$ . For all sufficiently small  $h \in \mathfrak{R}^{m(m+1)/2}$ , we have  $X + H \succ 0$ , where  $H = \mathbf{smat}(h)$ , and moreover,

$$(X + H)^{-1} = X^{-1} - X^{-1} H X^{-1} + o(H), \quad (20)$$

and hence

$$\begin{aligned} \mathcal{C}_K(X + H) &= (K^T X^{-1} K - K^T X^{-1} H X^{-1} K + o(H))^{-1} \\ &= \mathcal{C}_K(X) + \mathcal{C}_K(X) K^T X^{-1} H X^{-1} K \mathcal{C}_K(X) + o(H). \end{aligned} \quad (21)$$

From (21) and the definition of  $\Phi$ , we obtain further that

$$\begin{aligned} \Phi(X + H) &= \Psi(\mathcal{C}_K(X + H)) = \Psi(\mathcal{C}_K(X) + \mathcal{C}_K(X) K^T X^{-1} H X^{-1} K \mathcal{C}_K(X) + o(H)) \\ &= \Phi(X) + \text{tr}(\nabla \Psi(\mathcal{C}_K(X)) \mathcal{C}_K(X) K^T X^{-1} H X^{-1} K \mathcal{C}_K(X)) + o(H) \\ &= \Phi(X) + \text{tr}(X^{-1} K \mathcal{C}_K(X) \nabla \Psi(\mathcal{C}_K(X)) \mathcal{C}_K(X) K^T X^{-1} H) + o(H). \end{aligned} \quad (22)$$

In view of the definitions of  $\phi$ ,  $\Phi$ ,  $X$  and  $H$ , it follows from (22) and (2) that

$$\phi(x + h) - \phi(x) = \Phi(X + H) - \Phi(X) = h^T \mathbf{svec}(X^{-1} K \mathcal{C}_K(X) \nabla \Psi(\mathcal{C}_K(X)) \mathcal{C}_K(X) K^T X^{-1}) + o(h),$$

We next derive the Hessian of  $\phi$  at any  $x \in \mathbf{svec}(\mathcal{S}_{++}^m)$ . To proceed, we first recall the following well-known results (see, for example, page 243 and Lemma 4.3.1 of [22]):

$$\mathbf{vec}(ABC) = (C^T \otimes A) \mathbf{vec}(B), \quad (A \otimes B)^T = A^T \otimes B^T. \quad (23)$$

Let  $X$ ,  $h$  and  $H$  be defined as above. Using (18), (20) and (21), we have

$$\begin{aligned} \nabla\Phi(X+H) &= (X+H)^{-1}K\mathcal{C}_K(X+H)\nabla\Psi(\mathcal{C}_K(X+H))\mathcal{C}_K(X+H)K^T(X+H)^{-1} \\ &= \nabla\Phi(X) + G_0\nabla^2\Psi(\mathcal{C}_K(X))[G_0^T H G_0]G_0^T - X^{-1}H G_2 - G_2 H X^{-1} \\ &\quad + G_1 H G_2 + G_2 H G_1 + o(H), \end{aligned} \quad (24)$$

where  $G_0$ ,  $G_1$  and  $G_2$  are defined as above, and  $\nabla^2\Psi(\mathcal{C}_K(X))[\cdot]$  denotes the linear operator from  $\mathcal{S}^k$  to  $\mathcal{S}^k$  whose matrix representation is  $\nabla^2\psi(\mathbf{svec}(\mathcal{C}_K(X)))$ . Using (4), (23), and the definition of  $Q_1$  and  $Q_2$ , we obtain that

$$\begin{aligned} \mathbf{svec}(G_0\nabla^2\Psi(\mathcal{C}_K(X))[G_0^T H G_0]G_0^T) &= Q_1 \mathbf{vec}(G_0\nabla^2\Psi(\mathcal{C}_K(X))[G_0^T H G_0]G_0^T) \\ &= Q_1(G_0 \otimes G_0) \mathbf{vec}(\nabla^2\Psi(\mathcal{C}_K(X))[G_0^T H G_0]) = Q_1(G_0 \otimes G_0)Q_2^T \mathbf{svec}(\nabla^2\Psi(\mathcal{C}_K(X))[G_0^T H G_0]) \\ &= Q_1(G_0 \otimes G_0)Q_2^T \nabla^2\psi(\mathbf{svec}(\mathcal{C}_K(X)))\mathbf{svec}(G_0^T H G_0) \\ &= Q_1(G_0 \otimes G_0)Q_2^T \nabla^2\psi(\mathbf{svec}(\mathcal{C}_K(X)))Q_2(G_0 \otimes G_0)^T Q_1^T \mathbf{svec}(H). \end{aligned} \quad (25)$$

Moreover, since  $G_1$  and  $G_2$  are symmetric, we further have from (3) that

$$\mathbf{svec}(G_2 H X^{-1} + X^{-1} H G_2) = 2[X^{-1} \otimes_s G_2] \mathbf{svec}(H), \quad (26)$$

$$\mathbf{svec}(G_2 H G_1 + G_1 H G_2) = 2[G_1 \otimes_s G_2] \mathbf{svec}(H). \quad (27)$$

These together with (16), the definition of  $X$  and  $H$ , and the fact that  $\mathbf{svec}$  is the adjoint operator of  $\mathbf{smat}$  yield

$$\begin{aligned} \nabla\phi(x+h) - \nabla\phi(x) &= \mathbf{svec}(\nabla\Phi(X+H) - \nabla\Phi(X)) \\ &= (Q_1(G_0 \otimes G_0)Q_2^T \nabla^2\psi(\mathbf{svec}(\mathcal{C}_K(X)))Q_2(G_0 \otimes G_0)^T Q_1^T - 2X^{-1} \otimes_s G_2 + 2G_1 \otimes_s G_2)h + o(h), \end{aligned}$$

and hence (19) holds.

Finally, suppose in addition that  $-\nabla\Psi(\mathcal{C}_K(X))$  and  $\nabla^2\psi(\mathbf{svec}(\mathcal{C}_K(X)))$  are positive definite matrices. We show that the rank of  $\nabla^2\phi(x)$  is  $m(m+1)/2 - (m-k)(m-k+1)/2$ . To this end, it suffices to know the dimension of the null space of  $\nabla^2\phi(x)$ , denoted by  $\text{Null}(\nabla^2\phi(x))$ . Notice that  $\phi$  is a twice differentiable convex function in  $\mathbf{svec}(\mathcal{S}_{++}^m)$ . Hence,  $\nabla^2\phi(x) \succeq 0$ . It implies that  $h \in \text{Null}(\nabla^2\phi(x))$  if and only if  $h^T \nabla^2\phi(x)h = 0$ . We will subsequently show that

$$h^T \nabla^2\phi(x)h = 0 \Leftrightarrow K^T X^{-1} H = 0, \quad (28)$$

where  $H = \mathbf{smat}(h)$ . It then follows that

$$h \in \text{Null}(\nabla^2\phi(x)) \Leftrightarrow K^T X^{-1} H = 0.$$

Notice that  $K^T X^{-1}$  has full row rank. Thus, there exist nonsingular matrices  $E_1$  and  $E_2$  such that  $K^T X^{-1} = E_1 \begin{pmatrix} I & 0 \end{pmatrix} E_2$ , where  $I$  is the identity matrix of order  $k$ . It then follows that

$$K^T X^{-1} H = 0 \Leftrightarrow \begin{pmatrix} I & 0 \end{pmatrix} U = 0,$$

where  $U = E_2 H E_2^T \in \mathcal{S}^m$ . It is easy to see that the dimension of  $\{U \in \mathcal{S}^m : \begin{pmatrix} I & 0 \end{pmatrix} U = 0\}$  is  $(m-k)(m-k+1)/2$ . Since  $E_2$  is invertible, we conclude that the dimension of  $\{H \in \mathcal{S}^m : K^T X^{-1} H = 0\}$  and hence of  $\text{Null}(\nabla^2\phi(x))$  is also  $(m-k)(m-k+1)/2$ . Since  $\mathbf{smat}$  is a one-to-one map between  $\mathfrak{R}^{m(m+1)/2}$  and  $\mathcal{S}^m$ , the rank of  $\nabla^2\phi(x)$  is  $m(m+1)/2 - (m-k)(m-k+1)/2$ . To complete the proof, we next show that (28) holds.

To proceed, we first note from (2), (25), (26), (27) and the definition of  $\nabla^2\Psi(\mathcal{C}_K(X))[\cdot]$  that

$$h^T \nabla^2 \phi(x) h = \text{tr}(\nabla^2 \Psi(\mathcal{C}_K(X)) [G_0^T H G_0] G_0^T H G_0) + 2\text{tr}(G_1 H G_2 H) - 2\text{tr}(X^{-1} H G_2 H),$$

where  $H = \text{smat}(h)$ . Hence,  $h^T \nabla^2 \phi(x) h = 0$  if and only if

$$\text{tr}(\nabla^2 \Psi(\mathcal{C}_K(X)) [G_0^T H G_0] G_0^T H G_0) + 2\text{tr}(G_1 H G_2 H) - 2\text{tr}(X^{-1} H G_2 H) = 0. \quad (29)$$

We claim that

$$\text{tr}(G_1 H G_2 H) - \text{tr}(X^{-1} H G_2 H) \geq 0. \quad (30)$$

Indeed, it follows from the definition of  $\mathcal{C}_K(X)$  and Lemma 3.1 that

$$(\mathcal{C}_K(X))^{-1} = K^T X^{-1} K \Rightarrow \begin{pmatrix} X & K \\ K^T & (\mathcal{C}_K(X))^{-1} \end{pmatrix} \succeq 0 \Rightarrow X - K \mathcal{C}_K(X) K^T \succeq 0.$$

Using this relation, the definition of  $G_1$  and  $G_2$  along with the assumption  $-\nabla\Psi(\mathcal{C}_K(X)) \succ 0$ , we have

$$\begin{aligned} & \text{tr}(G_1 H G_2 H) - \text{tr}(X^{-1} H G_2 H) \\ &= \text{tr}(X^{-1} H (-G_2) H) - \text{tr}(G_1 H (-G_2) H) \\ &= \text{tr}([X^{-1} - X^{-1} K \mathcal{C}_K(X) K^T X^{-1}] [H (-G_2) H]) \\ &= \text{tr}([X - K \mathcal{C}_K(X) K^T] [X^{-1} H G_0 (-\nabla\Psi(\mathcal{C}_K(X))) G_0^T H X^{-1}]) \geq 0, \end{aligned} \quad (31)$$

and hence (30) holds. In addition, from the assumption  $\nabla^2\psi(\text{svec}(\mathcal{C}_K(X))) \succ 0$ , we can observe that  $\text{tr}(\nabla^2\Psi(\mathcal{C}_K(X)) [G_0^T H G_0] G_0^T H G_0) \geq 0$ . In view of this relation and (30), we conclude that (29) is equivalent to

$$\text{tr}(\nabla^2\Psi(\mathcal{C}_K(X)) [G_0^T H G_0] G_0^T H G_0) = 0 = \text{tr}(G_1 H G_2 H) - \text{tr}(X^{-1} H G_2 H). \quad (32)$$

Further, using the assumption  $\nabla^2\psi(\text{svec}(\mathcal{C}_K(X))) \succ 0$  and the fact that  $\mathcal{C}_K(X) \succ 0$ , one can see that the first equality in (32) is equivalent to

$$0 = G_0^T H G_0 = \mathcal{C}_K(X) K^T X^{-1} H X^{-1} K \mathcal{C}_K(X) \Leftrightarrow K^T X^{-1} H X^{-1} K = 0. \quad (33)$$

Also, using (31), the definition of  $G_0$ , the assumption  $-\nabla\Psi(\mathcal{C}_K(X)) \succ 0$ , and the fact that  $X - K \mathcal{C}_K(X) K^T \succeq 0$  and  $\mathcal{C}_K(X) \succ 0$ , we can observe that the second equality in (32) is equivalent to

$$0 = G_0^T H X^{-1} [X - K \mathcal{C}_K(X) K^T]^{\frac{1}{2}} \Leftrightarrow K^T X^{-1} H X^{-1} [X - K \mathcal{C}_K(X) K^T]^{\frac{1}{2}} = 0. \quad (34)$$

Furthermore, notice that

$$\begin{aligned} \mathfrak{R}^m &= \text{Range}(X) = \text{Range}(X - K \mathcal{C}_K(X) K^T + K \mathcal{C}_K(X) K^T) \\ &= \text{Range}(X - K \mathcal{C}_K(X) K^T) + \text{Range}(K \mathcal{C}_K(X) K^T) \\ &= \text{Range}([X - K \mathcal{C}_K(X) K^T]^{\frac{1}{2}}) + \text{Range}(K), \end{aligned}$$

where the third equality follows from [29, Section 2.3] and the fact that both  $X - K \mathcal{C}_K(X) K^T$  and  $K \mathcal{C}_K(X) K^T$  are positive semidefinite. Combining this last relation with (33) and (34), we immediately conclude that (32), and hence (29), is equivalent to

$$K^T X^{-1} H X^{-1} = 0 \Leftrightarrow K^T X^{-1} H = 0.$$

Thus, (28) holds. This completes the proof.  $\blacksquare$

**Remark 4.1.** Recall from Assumption 3.1 that  $\Psi$  is convex and decreasing. Using these facts, it is not hard to observe that  $-\nabla\Psi(U)$  and  $\nabla^2\psi(\mathbf{svec}(U))$  are positive semidefinite at any  $U \in \mathcal{S}_{++}^k$ , which are slightly weaker than the conditions required in the second part of the above theorem. Nevertheless, for the commonly used criteria such as classical A-, D- and  $p$ th mean criterion, we can easily verify that they satisfy the latter conditions. ■

In the remainder of this section we tailor the above IP method to problem (1) with classical A-, D- and  $p$ th mean criterion, respectively. Recall that  $\nabla\phi$  and  $\nabla^2\phi$  are needed for computing the Newton direction arising in the IP method, where  $\phi$  is defined in (16). Though generic formulas for  $\nabla\phi$  and  $\nabla^2\phi$  are provided in Theorem 4.1, the structure of the associated function  $\Psi$  for each specific optimality criterion is not exploited. We next derive more computationally efficient formulas for them in the context of classical A-, D- and  $p$ th mean criterion.

#### 4.1 IP method for $p$ th mean criterion

Recall from Section 1 that the  $p$ th mean criterion  $\Phi$  is

$$\Phi(X) = \text{tr}((K^T X^\dagger K)^{-p}) \quad (35)$$

for some  $p < 0$  and  $K \in \mathfrak{R}^{m \times k}$  with full column rank. It is not hard to verify that Assumption 3.1 holds for  $\Phi$ . Hence, problem (1) with this criterion can be suitably solved by our IP method proposed in Section 3. We next derive computationally efficient formulas for the gradient and the Hessian of  $\phi$ . Before proceeding, we state the following classical result (see, for example, [11, Proposition 4.3]) that will be used subsequently.

**Lemma 4.1.** Let  $g : \mathfrak{R} \rightarrow \mathfrak{R}$  be a differentiable function and let  $g^\square : \mathcal{S}^m \rightarrow \mathcal{S}^m$  be defined by

$$g^\square(Y) := V \begin{pmatrix} g(d_1) & & & \\ & g(d_2) & & \\ & & \ddots & \\ & & & g(d_m) \end{pmatrix} V^T,$$

where  $V\mathcal{D}(d)V^T$  is an eigenvalue decomposition of  $Y$  for some  $d \in \mathfrak{R}^m$ . Then the function  $g^\square$  is well-defined, i.e., it is independent of the choice of  $V$  and  $d$ , and is also differentiable. Moreover, let  $S^{g,d} \in \mathcal{S}^m$  be a symmetric matrix whose  $(i, j)$ th entry is given by

$$s_{ij}^{g,d} := \begin{cases} \frac{g(d_i) - g(d_j)}{d_i - d_j} & \text{if } d_i \neq d_j, \\ g'(d_i) & \text{otherwise.} \end{cases}$$

Then the directional derivative of  $g^\square$  at  $Y$  along the direction  $H \in \mathcal{S}^m$  is given by

$$V(S^{g,d} \circ (V^T H V))V^T.$$

**Proposition 4.1.** Let  $\Phi$  be defined in (35) and the associated  $\phi$  be defined in (16). Let  $Q \in \mathfrak{R}^{m(m+1)/2 \times m^2}$  be defined in (4) for  $\mathcal{S}^m$ . Then the gradient and Hessian of  $\phi$  at any  $x \in \mathbf{svec}(\mathcal{S}_{++}^m)$  are given by

$$\nabla\phi(x) = p\mathbf{svec}(X^{-1}K[\mathcal{C}_K(X)]^{p+1}K^T X^{-1}), \quad (36)$$

$$\begin{aligned} \nabla^2\phi(x) = & Q(-p[(X^{-1}KV) \otimes (X^{-1}KV)]\mathcal{D}(\mathbf{vec}(S^{g,d}))[(X^{-1}KV) \otimes (X^{-1}KV)]^T)Q^T \\ & - 2p X^{-1} \otimes_s G, \end{aligned} \quad (37)$$

respectively, where  $X = \mathbf{smat}(x)$ ,  $\mathcal{C}_K(X) = (K^T X^\dagger K)^{-1}$ ,  $V\mathcal{D}(d)V^T$  is an eigenvalue decomposition of  $K^T X^{-1}K$  for some  $d \in \mathfrak{R}^m$ ,  $g(t) = t^{-p-1}$ , and  $G = X^{-1}K[\mathcal{C}_K(X)]^{p+1}K^T X^{-1}$ . In

particular, when  $K = I$ , the above gradient and Hessian reduce to

$$\nabla\phi(x) = p\mathbf{svec}(X^{p-1}), \quad (38)$$

$$\nabla^2\phi(x) = (V \otimes_s V) \mathcal{D}(\mathbf{svec}_0(S^{g,d}))(V \otimes_s V)^T, \quad (39)$$

where  $g(t) = pt^{p-1}$  and  $V \mathcal{D}(d)V^T$  is an eigenvalue decomposition of  $X$  for some  $d \in \mathfrak{R}^m$ .

*Proof.* We first derive (38). To this aim, we fix an arbitrary  $x \in \mathbf{svec}(S_{++}^m)$ . Let  $X = \mathbf{smat}(x)$ . For all sufficiently small  $h \in \mathfrak{R}^{m(m+1)/2}$ , we have  $X + H \succ 0$ , where  $H = \mathbf{smat}(h)$ . Applying Lemma 4.1 with  $g(t) = t^p$  and  $Y = X$ , we obtain that

$$\Phi(X + H) = \text{tr}((X + H)^p) = \Phi(X) + \text{tr}(V(S^{g,d} \circ (V^T H V))V^T) + o(H), \quad (40)$$

where  $V \mathcal{D}(d)V^T$  is an eigenvalue decomposition of  $X$ . Making use of the definition of  $S^{g,d}$  and the fact that  $V^T V = I$ , we further have

$$\begin{aligned} \text{tr}(V(S^{g,d} \circ (V^T H V))V^T) &= \text{tr}(S^{g,d} \circ (V^T H V)) = \sum_{i=1}^m s_{ii}^{g,d} \sum_{j,k} v_{ji} h_{jk} v_{ki} \\ &= p \sum_{j,k} \left( \sum_{i=1}^m v_{ji} d_i^{p-1} v_{ki} \right) h_{jk} = \text{tr}(pX^{p-1}H). \end{aligned} \quad (41)$$

In view of the definition of  $\phi$ ,  $\Phi$ ,  $X$  and  $H$ , it follows from (40), (41) and (2) that

$$\phi(x + h) - \phi(x) = \Phi(X + H) - \Phi(X) = h^T (p\mathbf{svec}(X^{p-1})) + o(h),$$

which yields (38). Formula (36) now immediately follows from (38) and (18) with  $\Psi(U) = \text{tr}(U^p)$  for any  $U \in S_{++}^k$ .

We next derive (37). Let  $X$ ,  $h$  and  $H$  be defined as above. Using (20) and Lemma 4.1 with  $g(t) = t^{-p-1}$  and  $Y = K^T X^{-1} K$ , we have

$$\begin{aligned} \nabla\Phi(X + H) &= \nabla\Phi(X) - p(X^{-1}K)V(S^{g,d} \circ (V^T K^T X^{-1} H X^{-1} K V))V^T (K^T X^{-1}) \\ &\quad - pGHX^{-1} - pX^{-1}HG + o(H), \end{aligned}$$

where  $G = X^{-1}K[\mathcal{C}_K(X)]^{p+1}K^T X^{-1}$ . Notice that by using the definition of  $Q$ , (4) and (23), we have

$$\begin{aligned} &\mathbf{svec}((X^{-1}KV)(S^{g,d} \circ (V^T K^T X^{-1} H X^{-1} K V))(V^T K^T X^{-1})) \\ &= Q[(X^{-1}KV) \otimes (X^{-1}KV)] \mathbf{vec}(S^{g,d} \circ (V^T K^T X^{-1} H X^{-1} K V)) \\ &= Q[(X^{-1}KV) \otimes (X^{-1}KV)] \mathcal{D}(\mathbf{vec}(S^{g,d})) \mathbf{vec}(V^T K^T X^{-1} H X^{-1} K V) \\ &= Q[(X^{-1}KV) \otimes (X^{-1}KV)] \mathcal{D}(\mathbf{vec}(S^{g,d})) [(X^{-1}KV) \otimes (X^{-1}KV)]^T Q^T \mathbf{svec}(H), \end{aligned}$$

and also from (3), we have

$$\mathbf{svec}(GHX^{-1} + X^{-1}HG) = 2[X^{-1} \otimes_s G] \mathbf{svec}(H).$$

Using these relations and proceeding as in Theorem 4.1, we can see that (37) holds.

For the case when  $K = I$ ,  $\nabla^2\phi$  can be directly derived as follows. Since  $\nabla\Phi(X) = pX^{p-1}$ , letting  $g(t) = pt^{p-1}$  and  $V \mathcal{D}(d)V^T$  be an eigenvalue decomposition of  $X$ , we have from Lemma 4.1 that

$$\nabla\Phi(X + H) = \nabla\Phi(X) + V(S^{g,d} \circ (V^T H V))V^T + o(H).$$

In view of (3), one can see that

$$\begin{aligned} \mathbf{svec}(V(S^{g,d} \circ (V^T H V))V^T) &= (V \otimes_s V) \mathbf{svec}(S^{g,d} \circ (V^T H V)) \\ &= (V \otimes_s V) [\mathbf{svec}_0(S^{g,d}) \circ \mathbf{svec}(V^T H V)] \\ &= (V \otimes_s V) (\mathbf{svec}_0(S^{g,d}) \circ [(V \otimes_s V)^T \mathbf{svec}(H)]) \\ &= (V \otimes_s V) \mathcal{D}(\mathbf{svec}_0(S^{g,d})) (V \otimes_s V)^T \mathbf{svec}(H). \end{aligned}$$

Using these relations and proceeding as in Theorem 4.1, we conclude that (39) holds.  $\blacksquare$

For  $p$ th mean criterion, the associated function  $\Psi$  is given by  $\Psi(U) = \text{tr}(U^p)$  for any  $U \in \mathcal{S}_{++}^k$ . It is clear that  $-\nabla\Psi(U)$  and  $\nabla^2\psi(\text{svec}(U))$  are positive definite for any  $U \in \mathcal{S}_{++}^k$ . The following result is an immediate consequence of Theorem 4.1 on the rank of  $\nabla^2\phi(x)$  at any  $x \in \text{svec}(\mathcal{S}_{++}^m)$ .

**Corollary 4.1.** *Let  $\Phi$  be defined in (35) and the associated  $\phi$  be defined in (16). Then the rank of  $\nabla^2\phi(x)$  is  $m(m+1)/2 - (m-k)(m-k+1)/2$  for any  $x \in \text{svec}(\mathcal{S}_{++}^m)$ .*

## 4.2 IP method for A-criterion

Recall from Section 1 that the A-criterion  $\Phi$  is

$$\Phi(X) = \text{tr}(K^T X^\dagger K) \quad (42)$$

for some  $K \in \mathfrak{R}^{m \times k}$  with full column rank. Since A-criterion is a special case of  $p$ th mean criterion with  $p = -1$ , the IP method discussed in Sections 3 and 4.1 can be suitably applied to solve problem (1) with A-criterion. By exploiting the special structure of such  $\Phi$ , we next present a more compact representation of the associated Hessian matrix that is used to compute Newton direction for our IP method. The proof is routine and is thus omitted.

**Proposition 4.2.** *Let  $\Phi$  be defined in (42) and the associated  $\phi$  be defined in (16). Then the gradient and Hessian of  $\phi$  at any  $x \in \text{svec}(\mathcal{S}_{++}^m)$  are given by*

$$\begin{aligned} \nabla\phi(x) &= -\text{svec}(X^{-1} K K^T X^{-1}), \\ \nabla^2\phi(x) &= 2X^{-1} \otimes_s (X^{-1} K K^T X^{-1}), \end{aligned}$$

where  $X = \text{smat}(x)$ .

Since A-criterion is a special case of  $p$ th mean criterion, it follows from Corollary 4.1 that the rank of  $\nabla^2\phi(x)$  is also  $m(m+1)/2 - (m-k)(m-k+1)/2$  for every  $x \in \text{svec}(\mathcal{S}_{++}^m)$ .

## 4.3 IP method for D-criterion

Recall from Section 1 that the D-criterion  $\Phi$  is

$$\Phi(X) = \log \det(K^T X^\dagger K) \quad (43)$$

for some  $K \in \mathfrak{R}^{m \times k}$  with full column rank. It is not hard to verify that Assumption 3.1 is satisfied. Hence, problem (1) with this criterion can be suitably solved by the IP method studied in Section 3. We next present computationally efficient formulas for evaluating gradient and Hessian of the associated function  $\phi$  that are used in the IP method. The proof is routine and is thus omitted.

**Proposition 4.3.** *Let  $\Phi$  be defined in (43) and the associated  $\phi$  be defined in (16). Then the gradient and Hessian of  $\phi$  at any  $x \in \text{svec}(\mathcal{S}_{++}^m)$  are given by*

$$\begin{aligned} \nabla\phi(x) &= -\text{svec}(X^{-1} K \mathcal{C}_K(X) K^T X^{-1}), \\ \nabla^2\phi(x) &= 2X^{-1} \otimes_s (X^{-1} K \mathcal{C}_K(X) K^T X^{-1}) - (X^{-1} K \mathcal{C}_K(X) K^T X^{-1}) \otimes_s (X^{-1} K \mathcal{C}_K(X) K^T X^{-1}), \end{aligned}$$

where  $X = \text{smat}(x)$ .

For D-criterion, the associated function  $\Psi$  is given by  $\Psi(U) = -\log \det(U)$  for any  $U \in \mathcal{S}_{++}^k$ . It is clear that  $-\nabla\Psi(U)$  and  $\nabla^2\psi(\text{svec}(U))$  are positive definite for any  $U \in \mathcal{S}_{++}^k$ . The following result is again an immediate consequence of Theorem 4.1 on the rank of  $\nabla^2\phi(x)$  at any  $x \in \text{svec}(\mathcal{S}_{++}^m)$ .

**Corollary 4.2.** *Let  $\Phi$  be defined in (43) and the associated  $\phi$  be defined in (16). Then the rank of  $\nabla^2\phi(x)$  is  $m(m+1)/2 - (m-k)(m-k+1)/2$  for any  $x \in \text{svec}(\mathcal{S}_{++}^m)$ .*

## 5 Computational results

In this section, we conduct numerical experiments to test the performance of the IP method discussed in this paper for solving problem (1) with A-, D- and  $p$ th mean criterion and also compare its performance with the multiplicative algorithm and the standard IP solver SDPT3.

We develop Matlab codes for our IP method to solve (1) with A-, D- and  $p$ th mean criterion. We also implement the multiplicative algorithm in Matlab for solving (1) with A-, D- and  $p$ th mean criterion. To benchmark the performance of our IP method, we also report the computational results using a general SDP solver, namely, SDPT3 [36, 40] (Version 4.0) on solving a linear SDP reformulation of (1) with A-criterion (see [14, Page 532]) and a log-determinant SDP reformulation of (1) with D-criterion (see [27, Equation (10)]). We shall mention that it is not clear whether problem (1) with  $p$ th mean criterion can be reformulated into a problem that can be efficiently solved by SDPT3. As SDPT3 implements an infeasible path-following algorithm, we project the approximate solution  $w$  found by SDPT3 onto the unit simplex to obtain an approximate optimal feasible solution for problem (1) and the final objective value reported in our tests is based on the latter solution.<sup>2</sup> All computations in this section are performed in Matlab 7.14.0 (2012a) on a workstation with an Intel Xeon E5410 CPU (2.33 GHz) and 8GB RAM running Red Hat Enterprise Linux (kernel 2.6.18).

For our IP method, we set  $\tilde{w}^0 = \frac{1}{n}e \in \mathfrak{R}^{n-1}$ ,  $\mu_1 = 10$ ,  $\beta = \gamma = 0.5$ ,  $\sigma = 0.1$  and  $\eta = 0.95$ . In addition, we set  $\epsilon(\mu) = \max\{\mu, 10^{-10}\}$  and terminate the algorithm once  $\mu_t \leq 10^{-10}$ . On the other hand, for the multiplicative algorithm, similarly as in [45], we set  $\lambda = 1$ ,  $w^0 = \frac{1}{n}e \in \mathfrak{R}^n$ , and terminate the algorithm when it reaches 10000 iterations or

$$\max_{1 \leq i \leq n} d_i(w^t) \leq (1 + \delta) \sum_{i=1}^n w_i^t d_i(w^t)$$

holds with  $\delta = 2 \times 10^{-4}$ , where  $d_i(w)$  is defined in (5).<sup>3</sup> Furthermore, for SDPT3, we use the default tolerance. Finally, we use the mex files skron, smat and svec from the SDPT3 package for efficient operations on symmetric matrices in our implementation of the IP method and the multiplicative algorithm.

In our tests below, we consider the following four design spaces:

$$\begin{aligned} \chi_1(n) &= \{x_i = (e^{-s_i}, s_i e^{-s_i}, e^{-2s_i}, s_i e^{-2s_i})^T, 1 \leq i \leq n\}, \\ \chi_2(n) &= \{x_i = (1, s_i, s_i^2, s_i^3)^T, 1 \leq i \leq n\}, \\ \chi_3(n) &= \{x_{(i-1)\lceil\sqrt{n}\rceil+j} = (1, r_i, r_i^2, t_j, r_i t_j)^T, 1 \leq i, j \leq \lceil\sqrt{n}\rceil\}, \\ \chi_4(n) &= \{x_i = (t_i, t_i^2, \sin(2\pi t_i), \cos(2\pi t_i))^T, 1 \leq i \leq n\}, \end{aligned}$$

where  $s_i = \frac{3i}{n}$ ,  $r_i = \frac{2i}{n} - 1$  and  $t_i = \frac{i}{n}$ . The space  $\chi_1(n)$  represents the linearization of a compartmental model [4]. The space  $\chi_2(n)$  corresponds to polynomial regression. The third space, as described in [46], represents a response surface with a nonlinear effect and an interaction, while the fourth space is the quadratic/trigonometric example proposed in [44]. The test sets  $\chi_1$ ,  $\chi_3$  and a variant of the test set  $\chi_2$  are also used in [46].

In our first test, for each design space, we set  $A_i = x_i x_i^T$  in (1) for  $i = 1, \dots, n$ , with  $n = 10000, 50000, 100000$  for  $\chi_1, \chi_2, \chi_4$ , and  $n = 10000, 40000, 90000$  for  $\chi_3$ . For each  $n$  and each design space, we randomly generate 30 different matrices  $K \in \mathfrak{R}^{m \times 3}$  (i.e., we set  $k = 3$ ), each having i.i.d. Gaussian entries of mean 0 and variance 1. We then apply our IP method and the multiplicative algorithm to solve problem (1) with A-, D- and  $p$ th mean criterion on these instances and also apply SDPT3 to solve (1) with A- and D-criterion. The computational results averaged over the 30 instances are reported in Tables 1–4. In particular, the performance of our IP method, the multiplicative algorithm and SDPT3 are reported under the columns named “IP”, “MUL” and

<sup>2</sup>Such projection makes a difference when SDPT3 terminates early at a solution that is highly infeasible, which could be a consequence of “near infeasibility” of the linear SDP reformulation; see the first three rows of Table 1.

<sup>3</sup>We also tried  $\delta = 10^{-4}$ , but the multiplicative algorithm tends to take a long time for relatively little improvement on some instances.

Table 1: Computational results for A-criterion with random  $K$ 

$\chi_i$	$n$	cpu			obj		
		MUL	IP	SDPT3	MUL	IP	SDPT3
1	10000	13.76	0.69	1.93	193041	191410	211735
1	50000	62.13	3.75	10.64	154584	153219	172514
1	100000	135.26	7.38	19.22	208599	206787	242633
2	10000	17.69	0.77	1.90	215.754	212.356	212.356
2	50000	89.89	4.22	10.94	188.509	185.55	185.551
2	100000	163.28	8.12	21.83	242.414	237.823	237.824
3	10000	33.74	1.02	2.36	54.8551	54.7332	54.7332
3	40000	140.85	5.24	14.17	49.0008	48.9784	48.9791
3	90000	322.52	10.75	35.50	50.8124	50.7906	50.7906
4	10000	13.42	0.90	1.90	572.779	558.088	558.088
4	50000	58.41	4.46	9.89	501.924	487.99	487.991
4	100000	139.62	9.03	20.36	343.827	337.003	337.023

Table 2: Computational results for D-criterion with random  $K$ 

$\chi_i$	$n$	cpu			obj		
		MUL	IP	SDPT3	MUL	IP	SDPT3
1	10000	1.47	0.95	1.47	19.7352	19.7347	19.7356
1	50000	5.40	4.73	6.03	19.9312	19.9307	19.933
1	100000	14.17	9.31	12.34	19.7973	19.7968	19.7987
2	10000	2.10	0.79	1.57	5.95269	5.95229	5.9523
2	50000	20.60	4.07	6.78	5.30436	5.3039	5.3039
2	100000	51.43	8.33	13.32	5.08652	5.08608	5.08609
3	10000	4.31	1.05	1.77	6.58713	6.58694	6.58694
3	40000	18.86	4.28	9.11	6.65124	6.65104	6.65103
3	90000	66.40	9.93	21.48	6.74346	6.74327	6.74382
4	10000	1.67	0.90	1.41	7.40587	7.40535	7.40535
4	50000	13.46	4.35	6.09	7.65401	7.6535	7.6535
4	100000	39.01	7.90	12.03	8.66619	8.66575	8.66574

“SDPT3”, respectively. In addition, the CPU time abbreviated as “cpu” is in seconds and the objective value abbreviated as “obj” is rounded off to six significant digits. We see that our IP method significantly outperforms the multiplicative algorithm in terms of CPU time, and gives a smaller objective value in all instances. Moreover, our IP method also outperforms SDPT3 in CPU time and gives a smaller objective value in most instances. Furthermore, it is worth pointing out that SDPT3 reports infeasibility and hence *early terminates* when solving some instances for  $\chi_1$  with A-criterion, possibly due to bad scaling of  $\mathcal{M}(w)$ . This accounts for its significantly larger objective values in Table 1 corresponding to  $\chi_1$ . Finally, for  $p$ th mean criterion with  $p < -1$ , our IP method achieves significantly better objective values than the multiplicative algorithm, where the objective value of the latter algorithm is chosen to be the minimum over all iterations (see Table 4). This phenomenon is actually not surprising since the multiplicative algorithm is only known to converge for  $p \in (-1, 0)$ , but it may not converge when  $p < -1$ .

In our second test, we consider the case when  $K = I$ . The instances used in this test are the same as those in the first test except  $K = I$ . We also apply our IP method and the multiplicative algorithm to solve problem (1) with A-, D- and  $p$ th mean criterion on these instances and apply SDPT3 to solve (1) with A- and D-criterion. The computational results are reported in Tables 5–8. We again observe that our IP method outperforms the multiplicative algorithm in terms of objective value in all instances, and is generally much faster on large instances. Furthermore, our IP method is usually faster than SDPT3 and produces comparable or smaller objective values.

Table 3: Computational results for  $p$ th mean criterion with random  $K$  for some  $p \in (-1, 0)$

$\chi_i$	$n$	$p = -0.25$				$p = -0.75$			
		cpu		obj		cpu		obj	
		MUL	IP	MUL	IP	MUL	IP	MUL	IP
1	10000	6.39	0.85	25.4567	25.4558	9.17	0.79	5187.73	5187.23
1	50000	36.84	4.18	25.1902	25.1894	45.72	4.09	7128.51	7126.32
1	100000	82.59	8.29	25.1312	25.1304	118.04	8.11	7207.59	7205.68
2	10000	6.19	0.76	5.68067	5.68046	13.67	0.79	46.4144	46.4008
2	50000	28.10	3.89	6.00911	6.00886	73.65	4.28	58.2028	58.1903
2	100000	71.37	8.00	6.12458	6.12434	152.84	8.30	60.9256	60.9108
3	10000	5.47	0.98	5.58387	5.58379	3.73	1.07	24.8691	24.868
3	40000	18.97	4.18	5.56727	5.56718	17.21	4.66	23.5463	23.5451
3	90000	58.35	9.91	5.45907	5.45899	56.66	10.98	24.7679	24.7664
4	10000	1.84	0.90	7.30484	7.30456	2.88	0.90	118.871	118.859
4	50000	8.84	4.63	7.27622	7.27589	10.98	4.97	108.079	108.066
4	100000	31.14	8.81	7.30129	7.30102	45.27	9.70	128.676	128.662

Table 4: Computational results for  $p$ th mean criterion with random  $K$  for some  $p < -1$

$\chi_i$	$n$	$p = -1.1$				$p = -1.2$			
		cpu		obj		cpu		obj	
		mul	IP	mul	IP	mul	IP	mul	IP
1	10000	5.07	0.70	611960	602294	4.66	0.68	1.46813e+06	1.43891e+06
1	50000	20.64	3.59	541355	532904	19.71	3.59	2.0649e+06	2.02777e+06
1	100000	48.85	7.47	371942	365201	51.33	7.35	1.79042e+06	1.75803e+06
2	10000	6.25	0.80	373.376	359.802	5.34	0.79	650.345	629.288
2	50000	21.38	4.12	492.463	476.123	21.15	4.13	667.047	645
2	100000	61.53	8.44	302.083	288.88	62.32	8.41	539.641	514.087
3	10000	19.16	1.11	74.7421	71.4397	20.40	1.18	95.224	88.142
3	40000	71.26	4.76	69.2354	65.8478	68.95	4.79	127.857	116.933
3	90000	204.48	11.78	69.2994	65.8143	160.07	11.85	109.287	100.475
4	10000	6.75	0.92	961.75	910.571	7.17	0.96	1640.19	1524.98
4	50000	27.46	4.91	903.773	846.954	36.12	5.02	1631.8	1520.71
4	100000	75.39	9.85	824.269	776.036	75.90	9.82	1710.28	1596.1

Table 5: Computational results for A-criterion with  $K = I$

$\chi_i$	$n$	cpu			obj		
		mul	IP	SDPT3	mul	IP	SDPT3
1	10000	13.69	0.74	2.29	54286.3	53848.3	53848.4
1	50000	62.08	4.17	12.23	54245.2	53807.3	54103.8
1	100000	133.65	7.37	27.46	54240.1	53802.1	54103.8
2	10000	16.53	0.81	1.82	73.4521	72.4443	72.4443
2	50000	75.99	4.26	11.03	73.391	72.385	72.3853
2	100000	164.18	8.60	20.23	73.3837	72.3778	72.3777
3	10000	1.58	0.93	2.13	21.6203	21.6191	21.6191
3	40000	12.81	4.38	11.38	21.2826	21.2812	21.2812
3	90000	36.66	9.14	30.21	21.1721	21.1706	21.1706
4	10000	12.84	0.96	1.58	174.279	170.775	170.775
4	50000	59.76	5.19	9.51	174.276	170.775	170.775
4	100000	128.73	9.93	17.13	174.277	170.775	170.776

## 6 Concluding remarks

In this paper we propose an IP method for solving problem (1) with a broad class of convex optimality criteria and establish its global convergence. We demonstrate how the Newton direction

Table 6: Computational results for D-criterion with  $K = I$

$\chi_i$	$n$	cpu			obj		
		mul	IP	SDPT3	mul	IP	SDPT3
1	10000	1.11	1.02	0.87	20.5125	20.5119	20.5125
1	50000	4.86	4.67	3.86	20.5098	20.5091	20.5091
1	100000	14.77	9.13	7.59	20.5094	20.5087	20.5088
2	10000	1.92	0.74	1.01	0.410745	0.410221	0.41022
2	50000	16.74	3.80	4.75	0.409964	0.409267	0.40926
2	100000	55.28	6.96	8.89	0.409795	0.409154	0.409145
3	10000	1.76	0.89	1.16	5.14292	5.14267	5.14267
3	40000	15.90	3.99	6.53	5.08236	5.08212	5.08211
3	90000	47.07	8.70	15.38	5.06226	5.06202	5.06201
4	10000	1.35	1.02	0.94	7.25257	7.25189	7.25189
4	50000	11.16	5.04	4.17	7.25253	7.2519	7.25189
4	100000	35.09	9.86	8.14	7.25246	7.2519	7.25189

Table 7: Computational results for  $p$ th mean criterion with  $K = I$  for some  $p \in (-1, 0)$

$\chi_i$	$n$	$p = -0.25$				$p = -0.75$			
		cpu		obj		cpu		obj	
		mul	IP	mul	IP	mul	IP	mul	IP
1	10000	7.48	0.91	23.3728	23.372	3.53	0.82	3635.71	3635.29
1	50000	42.29	4.21	23.3683	23.3675	24.12	4.37	3633.58	3633.2
1	100000	91.80	8.29	23.3677	23.367	57.01	8.86	3633.31	3632.94
2	10000	3.43	0.74	5.58855	5.58838	2.55	0.80	27.4836	27.4811
2	50000	20.56	4.00	5.58796	5.58771	11.69	4.14	27.4691	27.4653
2	100000	69.43	7.67	5.58785	5.58763	37.49	8.60	27.467	27.4634
3	10000	1.65	0.77	6.70457	6.70448	1.56	0.99	14.1435	14.1429
3	40000	14.46	4.24	6.68235	6.68225	13.22	3.82	13.9841	13.9834
3	90000	42.01	8.24	6.675	6.67491	37.85	9.25	13.9318	13.9311
4	10000	1.75	0.92	7.25984	7.25955	1.43	0.92	52.2922	52.286
4	50000	8.97	4.58	7.25988	7.25956	6.05	4.52	52.2937	52.286
4	100000	30.50	8.85	7.25983	7.25957	20.67	9.20	52.2927	52.2861

Table 8: Computational results for  $p$ th mean criterion with  $K = I$  for some  $p < -1$

$\chi_i$	$n$	$p = -1.1$				$p = -1.2$			
		cpu		obj		cpu		obj	
		mul	IP	mul	IP	mul	IP	mul	IP
1	10000	4.62	0.72	162818	159210	4.64	0.67	485415	471459
1	50000	18.15	3.64	162740	159077	18.28	3.63	482380	471030
1	100000	49.28	7.47	162732	159060	47.57	7.56	485149	470975
2	10000	4.59	0.79	108.922	108.171	4.65	0.80	165.133	162.297
2	50000	18.44	4.19	109.588	108.072	19.23	4.27	164.314	162.134
2	100000	47.95	8.66	109.495	108.06	45.78	8.87	165.458	162.114
3	10000	36.72	1.02	25.9565	25.7793	36.40	1.04	31.8264	30.8276
3	40000	142.62	4.17	25.599	25.3307	139.85	4.50	31.5254	30.2362
3	90000	328.00	9.60	25.5115	25.1841	322.93	9.68	31.46	30.0431
4	10000	6.51	0.95	297.604	277.597	7.99	0.89	497.138	453
4	50000	27.17	4.83	297.686	277.597	33.88	4.89	497.287	453
4	100000	63.40	9.57	297.696	277.597	81.00	10.02	497.306	453

can be efficiently computed when the method is applied to (1) with classical optimality criteria.

Our computational results show that the IP method outperforms the widely used multiplicative algorithm as well as the standard IP solver SDPT3 in both speed and solution quality. The codes for this paper, including our implementation of the multiplicative algorithm and our codes generating inputs for SDPT3, are available online at [www.math.sfu.ca/~zhaosong](http://www.math.sfu.ca/~zhaosong).

Finally, we would like to remark that the performance of our IP method depends on whether the Newton direction can be computed accurately and efficiently. In our implementation, we observe that for  $p$ th mean criterion with large  $|p|$ , as well as for the design space  $\{x_i = (1, s_i, s_i^2, s_i^3, s_i^4)^T, 1 \leq i \leq n\}$  with  $n \geq 50000$  and some random  $K \in \mathfrak{R}^{m \times 3}$ , the Newton direction cannot be computed accurately due to numerical errors and hence our IP method fails to terminate with a good approximate solution, compared with the multiplicative algorithm. Indeed, it is known [41, 42] that the performance of a barrier method deteriorates as  $\mu \rightarrow 0$ . It is conceivable that such issues would not arise if a primal-dual IP method was used instead. However, it is much more involved to develop a primal-dual IP method for solving (1): since the feasible set of (1) is not closed in general, one would have to develop a primal-dual IP method on an equivalent nonlinear semidefinite programming reformulation of (1). We leave this as a future research direction.

## Appendix

We present the proof of Theorem 3.1 in this appendix.

*Proof.* In this proof, we denote by  $\tilde{w}(\mu, v)$  the vector obtained from  $w(\mu, v)$  by dropping the last entry for all  $(\mu, v) \in \mathfrak{R}_{++} \times \mathfrak{R}^{n-1}$ . Notice that  $\tilde{w}(\mu, v)$  is the unique optimal solution of (15).

We now prove part (a). Let

$$\bar{f}^* := \inf_w \{\Phi_{\mathcal{M}}(w) : e^T w = 1, w > 0\}. \quad (44)$$

We first show that  $\lim_{(\mu, v) \rightarrow (0_+, 0)} \Phi_{\mathcal{M}}(w(\mu, v)) = \bar{f}^*$ .

Given an arbitrary  $\epsilon > 0$ , there exists a positive  $\tilde{w}$  satisfying  $e^T \tilde{w} < 1$  such that  $f(\tilde{w}) < \bar{f}^* + \epsilon/2$ . Then we have that for any  $v \in \mathfrak{R}^{n-1}$ ,

$$f_{\mu}(\tilde{w}(\mu, v)) - v^T \tilde{w}(\mu, v) \leq f_{\mu}(\tilde{w}) - v^T \tilde{w}. \quad (45)$$

On the other hand, note that  $\tilde{w}(\mu, v) > 0$  and  $e^T \tilde{w}(\mu, v) < 1$ . Hence,

$$-\sum_{i=1}^{n-1} \log(\tilde{w}_i(\mu, v)) - \log(1 - e^T \tilde{w}(\mu, v)) > 0$$

and  $f(\tilde{w}(\mu, v)) \geq \bar{f}^*$ . In view of these inequalities, (45) and the fact that  $\|\tilde{w}(\mu, v)\|_1 \leq 1$  and  $\|\tilde{w}\|_1 \leq 1$ , one can obtain that for any  $(\mu, v) \in \mathfrak{R}_{++} \times \mathfrak{R}^{n-1}$ ,

$$\begin{aligned} \bar{f}^* \leq f(\tilde{w}(\mu, v)) &= f_{\mu}(\tilde{w}(\mu, v)) + \mu \sum_{i=1}^{n-1} \log(\tilde{w}_i(\mu, v)) + \mu \log(1 - e^T \tilde{w}(\mu, v)) \\ &\leq f_{\mu}(\tilde{w}(\mu, v)) \leq f_{\mu}(\tilde{w}) + v^T \tilde{w}(\mu, v) - v^T \tilde{w} \\ &\leq f(\tilde{w}) - \mu \sum_{i=1}^{n-1} \log(\tilde{w}_i) - \mu \log(1 - e^T \tilde{w}) + 2\|v\|_{\infty} \\ &\leq \bar{f}^* + \frac{\epsilon}{2} - \mu \sum_{i=1}^{n-1} \log(\tilde{w}_i) - \mu \log(1 - e^T \tilde{w}) + 2\|v\|_{\infty}. \end{aligned}$$

Thus, there exists some  $\delta > 0$  such that  $\bar{f}^* \leq f(\tilde{w}(\mu, v)) \leq \bar{f}^* + \epsilon$  whenever  $\|(\mu, v)\| < \delta$ ,  $\mu > 0$ . Hence,  $\Phi_{\mathcal{M}}(w(\mu, v)) = f(\tilde{w}(\mu, v)) \rightarrow \bar{f}^*$  as  $(\mu, v) \rightarrow (0_+, 0)$ .

We next show that  $f^* = \bar{f}^*$ . Clearly,  $f^* \leq \bar{f}^*$ . We now suppose for contradiction that  $f^* < \bar{f}^*$ . By the definition of  $f^*$  and  $\bar{f}^*$ , there exist  $w^1$  and  $w^2$  which are feasible points of (1) and (44), respectively, so that  $\Phi_{\mathcal{M}}(w^1) < (f^* + \bar{f}^*)/2$  and  $\Phi_{\mathcal{M}}(w^2) < \bar{f}^* + (\bar{f}^* - f^*)/2$ . Let  $w = (w^1 + w^2)/2$ . Clearly,  $w > 0$ ,  $e^T w = 1$  and  $\text{Range}(K) \subseteq \text{Range}(\mathcal{M}(w))$  due to  $\mathcal{M}(w) \succ 0$ . By convexity of  $\Phi$  in  $\mathcal{S}_+^n(K)$ , we obtain that  $\Phi_{\mathcal{M}}(w) \leq (\Phi_{\mathcal{M}}(w^1) + \Phi_{\mathcal{M}}(w^2))/2 < \bar{f}^*$ , which is a contradiction to the definition of  $\bar{f}^*$ . Thus,  $\lim_{(\mu, v) \rightarrow (0_+, 0)} \Phi_{\mathcal{M}}(w(\mu, v)) = \bar{f}^* = f^*$ .

Now suppose that  $w^*$  is an accumulation point of  $w(\mu, v)$  as  $(\mu, v) \rightarrow (0_+, 0)$ . We next show that  $w^*$  is an optimal solution of (1). Indeed, it follows from (7) that for any feasible point  $w$  of (1),

$$\Phi_{\mathcal{M}}(w) = \inf_U \{ \Psi(U) : \mathcal{M}(w) \succeq KUK^T, U \succ 0 \}. \quad (46)$$

In view of (46), for each  $(\mu, v) \in \mathfrak{R}_{++} \times \mathfrak{R}^{n-1}$ , there exists  $U(\mu, v) \succ 0$  such that

$$\Phi_{\mathcal{M}}(w(\mu, v)) + \|(\mu, v)\| > \Psi(U(\mu, v)) \quad \text{and} \quad \mathcal{M}(w(\mu, v)) \succeq KU(\mu, v)K^T. \quad (47)$$

From the second relation in (47), we see that  $\text{tr}(\mathcal{M}(w(\mu, v))) \geq \lambda_{\min}(K^T K) \text{tr}(U(\mu, v))$ , from which it follows that  $U(\mu, v)$  is bounded and thus it has an accumulation point as  $(\mu, v) \rightarrow (0_+, 0)$ . Let  $U^*$  be such an accumulation point. In view of the first relation in (47) and the assumption on  $\Psi$ , we see that  $U^* \succ 0$ . Moreover, we obtain by taking limit in (47) upon  $(\mu, v) \rightarrow (0_+, 0)$  that

$$\lim_{(\mu, v) \rightarrow (0_+, 0)} \Phi_{\mathcal{M}}(w(\mu, v)) \geq \Psi(U^*), \quad \mathcal{M}(w^*) \succeq KU^*K^T. \quad (48)$$

The second relation in (48) together with Lemma 3.1 implies that

$$\mathcal{M}(w^*) \succeq KU^*K^T \Rightarrow \begin{pmatrix} U^{*-1} & K^T \\ K & \mathcal{M}(w^*) \end{pmatrix} \succeq 0 \Rightarrow \text{Range}(K) \subseteq \text{Range}(\mathcal{M}(w^*)).$$

Hence,  $w^*$  is a feasible point of (1). In view of (46), the first relation in (48) and the result  $\lim_{(\mu, v) \rightarrow (0_+, 0)} \Phi_{\mathcal{M}}(w(\mu, v)) = f^*$ , we have

$$\Phi_{\mathcal{M}}(w^*) \leq \Psi(U^*) \leq \lim_{(\mu, v) \rightarrow (0_+, 0)} \Phi_{\mathcal{M}}(w(\mu, v)) = f^*.$$

Thus,  $w^*$  is an optimal solution of (1). This proves part (a).

We next show that part (b) holds. Let  $w^*$  be an optimal solution of (1) with maximum cardinality. Then it follows immediately from assumption that  $\mathcal{M}(w^*) \succ 0$ . Thus, there exists a corresponding Lagrange multiplier  $u^*$  so that  $(w^*, u^*)$  satisfies (12). Let  $\tilde{w}^*$  be the vector obtained from  $w^*$  by dropping the last entry. In view of (9) and the first equation of (12) and (14), we observe that for any  $(\mu, v) \in \Xi_C$ ,

$$\begin{aligned} & (w(\mu, v) - w^*)^T (u(\mu, v) - u^*) \\ &= (P\tilde{w}(\mu, v) - P\tilde{w}^*)^T (u(\mu, v) - u^*) \\ &= (\tilde{w}(\mu, v) - \tilde{w}^*)^T P^T (\nabla \Phi_{\mathcal{M}}(w(\mu, v)) - \nabla \Phi_{\mathcal{M}}(w^*)) - (\tilde{w}(\mu, v) - \tilde{w}^*)^T v \\ &= (w(\mu, v) - w^*)^T (\nabla \Phi_{\mathcal{M}}(w(\mu, v)) - \nabla \Phi_{\mathcal{M}}(w^*)) - (\tilde{w}(\mu, v) - \tilde{w}^*)^T v \\ &\geq -2C\mu, \end{aligned}$$

where the last inequality holds since  $\Phi$  is convex in  $\mathcal{S}_{++}^n$ ,  $w(\mu, v), w^* \in \Omega$  and  $\|v\|_{\infty} < C\mu$ . Using this inequality and the third equation in (12) and (14), we see that

$$w^{*T} u(\mu, v) + w(\mu, v)^T u^* \leq w^{*T} u^* + w(\mu, v)^T u(\mu, v) + 2C\mu = (2C + n)\mu. \quad (49)$$

Dividing both sides of the above inequality by  $\mu$  and using the third equation of (14), we obtain that

$$\sum_{i=1}^n \frac{w_i^*}{w_i(\mu, v)} + \sum_{i=1}^n \frac{u_i^*}{u_i(\mu, v)} \leq 2C + n. \quad (50)$$

Since  $(w^*, u^*) \geq 0$  and  $(w(\mu, v), u(\mu, v)) > 0$ , it follows from (50) that for all  $i$ ,

$$w_i(\mu, v) \geq \frac{w_i^*}{2C + n}, \quad u_i(\mu, v) \geq \frac{u_i^*}{2C + n}. \quad (51)$$

It immediately implies that the  $i$ th entry of any accumulation point  $w^\diamond$  of  $w(\mu, v)$  as  $(\mu, v) \xrightarrow{\Xi_C} (0, 0)$  must be positive whenever  $w_i^* > 0$ . Since  $w^\diamond$  is an optimal solution of (1) by part (a), we conclude that part (b) holds. ■

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