AN EFFICIENT OPERATOR–SPLITTING METHOD FOR THE EIGENVALUE PROBLEM OF THE MONGE-AMPERE EQUATION

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Dedicated to the memory of Professor Roland Glowinski

Abstract. We develop an efficient operator–splitting method for the eigenvalue problem of the Monge–Ampère operator in the Aleksandrov sense. The backbone of our method relies on a convergent Rayleigh inverse iterative formulation proposed by Abedin and Kitagawa (Inverse iteration for the Monge–Ampère eigenvalue problem, Proceedings of the American Mathematical Society, 148 (2020) 4975–4886). Modifying the theoretical formulation, we develop an efficient algorithm for computing the eigenvalue and eigenfunction of the Monge–Ampère operator by solving a constrained Monge–Ampère equation during each iteration. Our method consists of four essential steps: (i) Formulate the Monge–Ampère eigenvalue problem as an optimization problem with a constraint; (ii) Adopt an indicator function to treat the constraint; (iii) Introduce an auxiliary variable to decouple the original constrained optimization problem into simpler optimization subproblems and associate the resulting new optimization problem with an initial value problem; and (iv) Discretize the resulting initial-value problem by an operator–splitting method in time and a mixed finite element method in space. The performance of our method is demonstrated by several experiments. Compared to existing methods, the new method is more efficient in terms of computational cost and has a comparable rate of convergence in terms of accuracy.

Keywords. Eigenvalue problem; Monge–Ampère equation; Inverse iteration; Operator–splitting method.

1. INTRODUCTION

The Monge–Ampère equation is a second-order fully nonlinear PDE in the form of
\[ \det D^2 u = f, \] (1.1)

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where $\mathbf{D}^2 u$ denotes the Hessian of $u$. The Monge-Ampère equation originates from differential geometry in which it describes a surface with prescribed Gaussian curvature [3, 35]. The existence, uniqueness and regularity of the solution has been extensively studied [3, 26, 44], and related applications can be found in optimal transport [4, 23], seismology [17], image processing [33], finance [45], and geostrophic flows [21].

Due to its broad applications, in the past decade, a lot of efforts have been devoted to developing numerical methods for the Monge-Ampère equation. One line of research is to develop wide-stencil based finite-difference schemes [24, 25] for equation (1.1) with Dirichlet boundary conditions. Such a class of methods utilizes the fact that $\det \mathbf{D}^2 u$ equals the product of the eigenvalues of $\mathbf{D}^2 u$, so that these methods use wide-stencils to estimate the eigenvalues. Later on, such methods were extended to accommodate transport boundary conditions in [23]. Another line of research is to design finite-element based methods. In [20, 22], the authors proposed the vanishing moment method, which approximates a fully nonlinear second-order PDE by a fourth-order PDE. In [9, 10, 14, 15], the authors formulate equation (1.1) as an optimization problem. Fast augmented Lagrangian algorithms are then designed to solve the new problems. Recently, operator–splitting methods have been proposed in [30, 39]. Taking advantage of the divergence form of $\det \mathbf{D}^2 u$, the authors of [30, 39] decouple the nonlinearity of equation (1.1) by introducing an auxiliary variable so that solving equation (1.1) is reduced to finding the steady-state solution of an initial value problem, which is time-discretized by an operator–splitting method and space-discretized by a mixed finite-element method. Other numerical methods for equation (1.1) include [2, 5, 6, 11, 12, 19]; see the survey [18] for more related works.

Existing works discussed above target equation (1.1) with various boundary conditions. Another interesting problem of the Monge–Ampère type is the eigenvalue problem, reading as

$$\begin{cases} \det (\mathbf{D}^2 u) = \lambda |u|^d & \text{in } \Omega, \\ u = 0 & \text{on } \partial \Omega, \end{cases} \quad (1.2)$$

where $\Omega \subset \mathbb{R}^d$ ($d \geq 2$) is an open bounded convex domain, and $\lambda = \lambda[\Omega]$ is the unknown eigenvalue of the Monge–Ampère operator on $\Omega$. Problem (1.2) was first studied by Lions in [38] and later by Tso in [46]. They proved the existence, uniqueness and regularity of the solution on an open, bounded, smooth, uniformly convex domain. The result was then extended by Le in [36] to general bounded convex domains. Theoretically, to find the solution of equation (1.2), a variational formulation was proposed in [46], and a convergent Rayleigh quotient inverse iterative formulation was proposed in [1] which was further improved in [37]. Since, during each Rayleigh quotient iteration, the algorithm in [1] requires solving a Monge–Ampère type equation, how to efficiently implement this formulation numerically has not been studied. The only work on the numerical solution of equation (1.2) we are aware of is [28], in which the authors proposed operator–splitting methods for a class of Monge–Ampère eigenvalue problem. In [28], taking advantage of the divergence form, the authors takes equation (1.2) as the optimality condition of a constrained optimization problem, in which $\lambda$ is considered as the Lagrange multiplier, and an operator–splitting method was proposed to solve the new problem.

Similar to equation (1.1), the eigenvalue problem (1.2) is a fully nonlinear second-order PDE. One effective way to solve such PDEs is the operator–splitting method, which decomposes
complicated problems into several easy-to-solve subproblems by introducing auxiliary variables. Then the new problem will be formulated as solving an initial value problem, which is then time discretized using operator--splittings. All variables will be updated in an alternative fashion, where each subproblem either has an explicit solution or can be solved efficiently. The operator--splitting method has been applied to numerically solving PDEs [30, 39], image processing [16, 40, 41, 42], surface reconstruction [34], inverse problems [29], obstacle problems [43], and computational fluid dynamics [7, 8]. We refer readers to monographs [31, 32] for detailed discussions on operator--splitting methods.

In this work, we propose an efficient numerical implementation of the formulation proposed in [1] to compute the eigenvalue and eigenfunction of the Monge–Ampère operator on an open, bounded, convex domain \( \Omega \). Since each Rayleigh quotient inverse iteration of the formulation in [1] requires solving a Monge–Ampère equation, we first use the divergence form of the Monge–Ampère operator to rewrite the problem as an optimization problem. To stabilize our formulation, we consider a constrained version of the optimization problem by forcing the eigenfunction \( u \) to have unit \( L_2 \)-norm: \( \| u \|_2 = 1 \). The constrained problem is converted to an unconstrained problem by utilizing an indicator function of the constraint set. Then we decouple the nonlinearity of the functional by introducing an auxiliary variable, and we associate it with an initial value problem in the flavor of gradient flow. The initial value problem is time discretized by an operator-splitting method and space discretized by a mixed finite-element method in the space of piecewise-linear continuous functions. The efficiency of the proposed method is demonstrated by several numerical experiments.

We organize the rest of this article as follows: We introduce the background and summarize the convergent formulation of [1] for equation (1.2) in Section 2. Our new operator-splitting approach for implementing this convergent formulation is presented in Section 3. Our operator-splitting scheme is time discretized in Section 4 and space discretized in Section 5. We demonstrate the efficiency of the proposed method by several numerical experiments in Section 6 and conclude this article in Section 7.

2. A CONVERGENT INVERSE ITERATION FOR THE EIGENVALUE PROBLEM

Let \( \Omega \subset \mathbb{R}^d \) be an open bounded convex domain. In equation (1.2), if \( u \) is a convex function, one has \( u \leq 0 \) and \( |u| = -u \). The existence and uniqueness of the eigen-pair was studied in [38]:

**Theorem 2.1.** Assume that \( \Omega \subset \mathbb{R}^d \) is a smooth, bounded, uniformly convex domain. There exist a unique positive constant \( \lambda_{\text{MA}} \) and a unique (up to positive multiplicative constants) nonzero convex function \( u \in C^{1,1}(\bar{\Omega}) \cap C^\infty(\Omega) \) solving the eigenvalue problem (1.2). The constant \( \lambda_{\text{MA}} \) is called the Monge–Ampère eigenvalue of \( \Omega \) and \( u \) is called a Monge–Ampère eigenfunction of \( \Omega \).

Define the Rayleigh quotient of a function \( u \) for the Monge–Ampère operator as

\[
R(u) = \frac{\int_{\Omega} -u \det(D^2u) \, dx}{\int_{\Omega} (-u)^{d+1} \, dx},
\]

and the function space \( \mathcal{K} \) as

\[
\mathcal{K} = \left\{ u \in C^{0,1}(\bar{\Omega}) \cap C^\infty(\Omega) : u \text{ is convex and nonzero in } \Omega, \ u = 0 \text{ on } \partial \Omega \right\}.
\]

Tso [46] showed that \( \lambda_{\text{MA}} \) can be written as the infimum of Rayleigh quotients:
**Theorem 2.2.** Assume that $\Omega \subset \mathbb{R}^d$ is a smooth, bounded and uniformly convex domain. Then

$$
\lambda_{\text{MA}} = \inf_{u \in \mathcal{K}} R(u). 
$$

(2.2)

Based on the property (2.2), the following inverse iterative scheme for the eigenvalue problem (1.2) was proposed by Abedin and Kitagawa in [1]:

$$
\begin{aligned}
&u^0 = u_0, \\
&\det(D^2 u^{k+1}) = R(u^k) |u^k|^d \quad \text{in } \Omega, \\
&u^{k+1} = 0 \quad \text{on } \partial \Omega,
\end{aligned}
$$

(2.3)

where $u_0$ is a given initial condition, and they further proved the convergence of the inverse iteration:

**Theorem 2.3.** Assume that $\Omega \subset \mathbb{R}^d$ is an open bounded convex domain. Let $u_0 \in C(\hat{\Omega})$ satisfy the following:

(i) $u_0$ is convex and $u_0 \leq 0$ on $\partial \Omega$;
(ii) $R(u_0) < \infty$;
(iii) $\det(D^2 u_0) \geq c_0$ in $\Omega$, where $c_0$ is some positive constant.

Then, for $k > 0$, $u^k$ in equation (2.3) converges uniformly on $\hat{\Omega}$ to a nonzero Monge-Ampère eigenfunction, and $R(u_k)$ converges to $\lambda_{\text{MA}}$.

Theorem 2.3 was improved in [37] so that conditions (i) and (iii) are removed; consequently, the inverse iteration converges for all convex initial data having finite and nonzero Rayleigh quotient to a nonzero Monge-Ampère eigenfunction of $\Omega$.

### 3. A Modified Formulation of the Inverse Iteration

Given an initial convex function $u_0$ with bounded nonzero Rayleigh quotient, the inverse iteration (2.3) generates the sequence $\{(R(u^k), u^k)\}$ which is guaranteed to converge to the solution of the eigenvalue problem (1.2). When updating $u^{k+1}$ from $u^k$, one needs to solve a Monge-Ampère equation with the Dirichlet boundary condition, which is a nonlinear problem. It has not been studied yet how to implement the inverse iteration efficiently to produce numerical approximations to the eigenvalue problem of the Monge-Ampère operator. Therefore, we are motivated to develop an efficient algorithm to implement this inverse iterative method.

To achieve this purpose, we adopt a recently developed operator-splitting method (see [28, 30, 39]) to solve equation (2.3) numerically. We focus on the case $d = 2$. Our method can be easily extended to higher dimensional problems.

We first reformulate equation (2.3) using the following identity:

$$
\det(D^2 u) = \frac{1}{2} \nabla \cdot (\text{cof}(D^2 u) \nabla u),
$$

(3.1)

where $\text{cof}(D^2 u)$ is the cofactor matrix of $D^2 u$. 
Incorporating equation (3.1) into equations (2.3) and (2.1) gives rise to
\[
\begin{align*}
\nabla \cdot (\text{cof}(\nabla^2 u_k) \nabla u_{k+1}) &= 2 R(u_k)|u_k|^d \quad \text{in } \Omega, \\
\\nabla u_{k+1} &= 0 \quad \text{on } \partial \Omega,
\end{align*}
\]
with
\[
R(u) = \frac{\int_\Omega (\text{cof}(\nabla^2 u) \nabla u) \cdot \nabla u \, d\mathbf{x}}{2 \int_\Omega (-u)^3 \, d\mathbf{x}},
\]
where we used integration by parts when deriving equation (3.3).

From equation (3.2), updating \(u_{k+1}\) from \(u_k\) is equivalent to solving the optimization problem
\[
\begin{align*}
\min_w \left[ \int_\Omega (\cof(\nabla^2 w) \nabla w) \cdot \nabla w \, d\mathbf{x} + 6 \int_\Omega f_k w \, d\mathbf{x} + I_S(w) \right],
\end{align*}
\]
with \(f = R(u_k)|u_k|^2\), which can be derived from the first-order variational principle; see [30, 39]. Note that if \((\lambda_{\text{MA}}, u^*)\) is a solution to equation (1.2), \((\lambda_{\text{MA}}, \alpha u^*)\) is also a solution for any \(\alpha > 0\) (assuming that we are looking for convex eigenfunctions). To make the solution of equation (1.2) unique, we restrict our attention to looking for the eigenfunction \(u^*\) satisfying \(\|u^*\|_2 = 1\). Therefore it is natural to add the constraint \(\|w\|_2 = 1\) to equation (3.4). However, usually a constrained optimization problem is more challenging to solve than an unconstrained one. Therefore, to remove the constraint while enforcing \(\|w\|_2 = 1\), we utilize an indicator function.

Define the set
\[
S = \{w : w \text{ is smooth}, \|w\|_2 = 1\}
\]
and its indicator function
\[
I_S(w) = \begin{cases} 0 & \text{if } w \in S, \\ +\infty & \text{otherwise}. \end{cases}
\]

Equation (3.4) with constraint \(\|w\|_2 = 1\) can be rewritten as
\[
\begin{align*}
\min_w \left[ \int_\Omega (\cof(\nabla^2 w) \nabla w) \cdot \nabla w \, d\mathbf{x} + 6 \int_\Omega f_k w \, d\mathbf{x} + I_S(w) \right],
\end{align*}
\]
with \(w = 0\) on \(\partial \Omega\).

We follow [30] to introduce a matrix-valued auxiliary variable \(p\) to decouple the nonlinearity in equation (3.5). Then solving equation (3.5) is equivalent to solving
\[
\begin{align*}
\min_{w,p} \left[ \int_\Omega (\cof(p) \nabla w) \cdot \nabla w \, d\mathbf{x} + 6 \int_\Omega f_k w \, d\mathbf{x} + I_S(w) \right],
\end{align*}
\]
with \(w = 0\) on \(\partial \Omega\), \(p = \nabla^2 w\) in \(\Omega\).
After computing the Euler-Lagrange equation, if \((v, p)\) is a solution to equation (3.6), we have

\[
\begin{align*}
\nabla \cdot (\text{cof}(p)\nabla v) - 2f^k + \partial I_S(v) &\geq 0 \quad \text{in } \Omega, \\
v & = 0 \quad \text{on } \partial \Omega, \\
p & = D^2v, \quad \text{in } \Omega,
\end{align*}
\]

(3.7)

where \(\partial I_S\) denotes the sub-differential of \(I_S\).

We associate equation (3.7) with the following initial value problem (in the flavor of gradient flow)

\[
\begin{align*}
\frac{\partial v}{\partial t} + \nabla \cdot ((\varepsilon I + \text{cof}(p))\nabla v) - 2f^k + \partial I_S(v) &\equiv 0 \quad \text{in } \Omega \times (0, +\infty), \\
v & = 0 \quad \text{on } \partial \Omega \times (0, +\infty), \\
\frac{\partial p}{\partial t} + \gamma(p - D^2v) & = 0 \quad \text{in } \Omega \times (0, +\infty), \\
v(0) & = v_0, \quad p(0) = p_0,
\end{align*}
\]

(3.8)

where \(I\) is the identity matrix, \(0\) is the zero matrix, and \(\varepsilon > 0\) is a small constant. The term \(\varepsilon I\) is a regularization term in order to handle the case that \(\inf_{x \in \Omega} f^k(x) = 0\). Then \(u^{k+1}\) is the steady state of \(v\).

In equation (3.8), \(\gamma\) controls the evolution speed of \(p\). A natural choice is to let \(p\) evolve with a similar speed as that of \(v\), leading to

\[\gamma = \beta \lambda_0\]

with \(\lambda_0\) being the smallest eigenvalue of \(-\nabla^2\) and \(\beta > 0\) being some constant.

4. AN OPERATOR SPLITTING METHOD TO SOLVE EQUATION (3.8)

4.1. The operator splitting strategy. The structure of equation (3.8) is well–suited to be time-discretized by the operator splitting method. Among many possible discretization schemes, we choose the simplest Lie scheme.

Let \(\tau > 0\) denote the time step and denote \(t^n = n\tau\). We time-discretize equation (3.8) as follows:

Initialization:

\[
v^0 = v_0, \quad p^0 = p_0.
\]

(4.1)

For \(n > 0\), update \((v^n, p^n) \to (v^{n+1/3}, p^{n+1/3}) \to (v^{n+2/3}, p^{n+2/3}) \to (v^{n+1}, p^{n+1})\) as:

**Step 1:** Solve

\[
\begin{align*}
\frac{\partial v}{\partial t} + \nabla \cdot ((\varepsilon I + \text{cof}(p))\nabla v) - 2f^k & = 0 \quad \text{in } \Omega \times (t^n, t^{n+1}), \\
v & = 0 \quad \text{on } \partial \Omega \times (t^n, t^{n+1}), \\
\frac{\partial p}{\partial t} & = 0 \quad \text{in } \Omega \times (t^n, t^{n+1}), \\
v(t^n) & = v^n, \quad p(t^n) = p^n,
\end{align*}
\]

(4.2)
and set \( v^{n+1/3} = v(t^{n+1}), \quad p^{n+1/3} = p(t^{n+1}) \).

**Step 2:** Solve

\[
\begin{aligned}
\frac{\partial v}{\partial t} &= 0 \quad \text{in } \Omega \times (t^n, t^{n+1}), \\
v &= 0 \quad \text{on } \partial \Omega \times (t^n, t^{n+1}), \\
\frac{\partial p}{\partial t} + \gamma (p - D^2 v) &= 0 \quad \text{in } \Omega \times (t^n, t^{n+1}), \\
v(t^n) &= v^{n+1/3}, \quad p(t^n) = p^{n+1/3},
\end{aligned}
\]

and set \( v^{n+2/3} = v(t^{n+1}), \quad p^{n+2/3} = p(t^{n+1}) \).

**Step 3:** Solve

\[
\begin{aligned}
\frac{\partial v}{\partial t} + \partial I_S(v) &\geq 0 \quad \text{in } \Omega \times (t^n, t^{n+1}), \\
v &= 0 \quad \text{on } \partial \Omega \times (t^n, t^{n+1}), \\
\frac{\partial p}{\partial t} &= 0 \quad \text{in } \Omega \times (t^n, t^{n+1}), \\
v(t^n) &= v^{n+2/3}, \quad p(t^n) = p^{n+2/3},
\end{aligned}
\]

and set \( v^{n+1} = v(t^{n+1}), \quad p^{n+1} = p(t^{n+1}) \).

The scheme (4.1)–(4.4) is only semi-constructive since one still needs to solve the subproblems in equations (4.2)–(4.4). For equation (4.3), we have the explicit solution for \( p^{n+2/3} \):

\[
p^{n+2/3} = e^{-\gamma \tau} p^n + (1 - e^{-\gamma \tau}) D^2 v^{n+1/3}.
\]

Since the solution of equation (1.2) is a convex function, the Hessian \( D^2 u \) is a semi-positive definite matrix. Since \( p \) is an auxiliary variable estimating \( D^2 v \), we project it onto the space of semi-positive definite symmetric matrices once \( p^{n+2/3} \) is computed. We denote the projection operator by \( P_+ \); see more details in Section 5.4.

For other subproblems, we adopt the one-step backward Euler scheme (the Markchuk-Yanenko type). Our updating formulas are summarized as follows:

\[
\begin{aligned}
\frac{v^{n+1/3} - v^n}{\tau} + \nabla \cdot \left[ (\epsilon I + \text{cof}(p^n)) \nabla v^{n+1/3} \right] - 2 f^k &= 0 \quad \text{in } \Omega, \\
v^{n+1/3} &= 0 \quad \text{on } \partial \Omega, \\
p^{n+1} &= P_+ \left( e^{-\gamma \tau} p^n + (1 - e^{-\gamma \tau}) D^2 v^{n+1/3} \right), \\
\frac{v^{n+1} - v^{n+1/3}}{\tau} + \partial I_S(v^{n+1}) &\geq 0 \quad \text{in } \Omega, \\
v^{n+1} &= 0 \quad \text{on } \partial \Omega.
\end{aligned}
\]

**Remark 4.1.** Equation (3.8) is very similar to problem (36) in [28], except that in our current scheme the constraint is \( \|u\|_2 = 1 \) and that in [28] it is \( \|u\|_3 = 1 \). Despite similar formulations, the numerical treatments are very different. In equations (4.5)–(4.7), \( f^k \) and the indicator function \( \partial I_S \) are separately distributed into two sub-steps. Equation (4.7) simply results in a projection to the unit sphere; see Section 4.2 for details.

In [28], \( \lambda d \, du | u | \) with \( d \) being the spatial dimension plays the role of \( f^k \) and the constraint plays the role of \( \partial I_S \), and both terms are arranged in the same sub-step (problem (50b) in [28]):

\[
\begin{aligned}
\int_\Omega |v^{n+2/3}|^3 \, dx &= 1 \\
\int_\Omega u^{n+2/3} = 3 \tau \lambda^{n+1} u^{n+2/3} |du^{n+2/3}|
\end{aligned}
\]
The constraint $\|u\|_3 = 1$ cannot be replaced by $\|u\|_2 = 1$ since equation (4.8) was considered as an optimality condition of a Lagrangian functional and $\tau\lambda^{n+1}$ is the Lagrange multiplier. As a result, $u^{n+2/3}$ solves

$$ u^{n+2/3} \in \arg\min_{v: \int_{\Omega} |v|^3 dx = 1} \left[ \frac{1}{2} \int_{\Omega} |v|^2 dx - \int_{\Omega} u^{n+1/3} v dx \right]. $$

Unlike (4.7), the solution to problem (4.9) does not have an explicit expression, so that an iterative method (such as sequential quadratic programming) was used in [28] to solve problem (4.9).

**Remark 4.2.** Compared to the algorithm (2.3) proposed in [1], our scheme has an additional term related to the constraint $\|u\|_2 = 1$, and such a constraint leads to the projection step (4.7) which helps stabilize our numerical algorithm.

**Remark 4.3.** Scheme (4.2)–(4.4) is an approximation of the gradient flow of the functional in (3.6). The convergence of scheme (4.2)–(4.4) is closely related to that of the gradient flow together with an approximation error. It has been shown that when there is only one variable and the operator in each step has sufficient regularity, the approximation error is of $O(\tau)$ (see [13] and [27, Chapter 6]). However, the terms in (3.6) are non-trivial and non-smooth, traditional analysis techniques are not applicable in this scenario. As the splitting error is closely related to the time step $\tau$, we expect the approximation error reduces (and thus the convergence of scheme (4.2)–(4.4) follows that of the gradient flow) as $\tau$ goes to 0.

### 4.2. On the solution to equation (4.7)

In the scheme above, problems (4.5) and (4.6) are easy to solve. In equation (4.7), $v^{n+1}$ solves

$$ v^{n+1} = \frac{v^{n+1/3}}{\|v^{n+1/3}\|_2}. $$

### 4.3. On the initial condition

We next discuss the initial condition $u_0$ in the outer iteration and $(v_0, p_0)$ in the inner iteration. The convergence theorem for the scheme (2.3), Theorem 2.3, requires the initial condition to be convex and smooth. A simple choice is to set $u_0$ as the solution to

$$ \det D^2 u_0 = 1 \quad \text{in } \Omega, $$

$$ u_0 = 0 \quad \text{on } \partial \Omega. $$

However, solving equation (4.12) is not trivial. Since $u_0$ is only the initial condition and the iterates generated by the inverse iteration are eventually smooth as shown in [37], we do not need to solve equation (4.12) exactly. An operator splitting method is proposed in [30] to solve equation (4.12). To make the initialization simpler, we will choose $u_0$ as the initial condition
according to a strategy used in [30]. Specifically, \( u_0 \) is the solution to the Poisson problem

\[
\begin{align*}
\nabla^2 u_0 &= 2\eta & \text{in } \Omega, \\
u_0 &= 0 & \text{on } \partial \Omega,
\end{align*}
\]

(4.13)

where \( \eta > 0 \) is of \( O(1) \).

For the initial condition \((v_0, p_0)\) in the \( k + 1 \)-th outer iteration, we simply set

\[
v_0 = u^k, \quad p_0 = D^2 v_0.
\]

(4.14)

Our algorithm is summarized in Algorithm 1.

**Algorithm 1: An operator-splitting method for solving problem (1.2)**

**Input:** Parameters \( \gamma, \tau, \varepsilon, N \).

**Initialization:** Set \( k = 0 \). Initialize \( u^0 \) according to equation (4.13).

while not converge do

Step 1. Compute \( f^k = R(u^k)|u^k|^2 \) according to equation (3.3).

Step 2. Set \( n = 0 \). Initialize \((v^0, p^0)\) according to equation (4.14).

while not converge do

Step 3.1. Solve equation (4.5) for \( v^{n+1/3} \).

Step 3.2. Solve equation (4.6) for \( p^{n+1} \).

Step 3.3. Solve equation (4.7) for \( v^{n+1} \).

Step 3.4. Set \( n = n + 1 \).

end while

Step 4. Set \( u^{k+1} \) as the converged \( v^* \).

Step 5. Set \( k = k + 1 \).

end while

**Output:** The converged eigenfunction \( u^* \) and eigenvalue \( \lambda_{\text{MA}} \).

5. A Finite Element Implementation of Scheme (4.5)-(4.7)

5.1. **Generalities.** Let \( \Omega \subset \mathbb{R}^2 \) be an open bounded convex polygonal domain (or it has been approximated by such a domain). Let \( \mathcal{T}_h \) be a triangulation of \( \Omega \), where \( h \) denotes the length of the longest edge of triangles in \( \mathcal{T}_h \). Define the following two piecewise linear function spaces

\[
\begin{align*}
V_h &= \{ \phi \in C^0(\bar{\Omega}) : \phi_T \in \mathbf{P}_1 \text{ for } \forall T \in \mathcal{T}_h \}, \\
V_{0h} &= \{ \phi \in V_h : \phi|_{\partial \Omega} = 0 \},
\end{align*}
\]

where \( \mathbf{P}_1 \) is the space of polynomials of two variables with degree no larger than 1. Let \( H^1(\Omega) \) be the Sobolev space of order 1 and \( H^1_0(\Omega) \) be the collection of functions in \( H^1(\Omega) \) with vanishing trace on \( \partial \Omega \). Then \( V_h \) and \( V_{0h} \) are approximations of \( H^1(\Omega) \) and \( H^1_0(\Omega) \), respectively.

Denote the set of vertices of \( \mathcal{T}_h \) by \( \Sigma_h \). We further denote the interior vertices of \( \mathcal{T}_h \) by \( \Sigma_{0h} = \Sigma_h \setminus (\Sigma_h \cap \partial \Omega) \). We use \( N_h \) and \( N_{0h} \) to denote the cardinality of \( \Sigma_h \) and \( \Sigma_{0h} \), respectively. We have

\[
\dim V_h = N_h \quad \text{and} \quad \dim V_{0h} = N_{0h}.
\]

We order the vertices of \( \mathcal{T}_h \) so that \( \Sigma_{0h} = \{Q_i\}_{i=1}^{N_{0h}} \), where \( Q_i \)'s denote the vertices. For any \( 1 \leq l \leq N_h \), we use \( \omega_l \) to denote the union of triangles in \( \mathcal{T}_h \) that have \( Q_l \) as a common vertex.
Denote the area of $\omega_l$ by $|\omega_l|$. For each vertex $Q_l$, we define the hat function $\phi_l$ so that
\[ \phi_l \in V_h, \quad \phi_l(Q_l) = 1 \text{ and } \phi_l(Q_m) = 0 \text{ for } m \neq l. \]
We have that $\phi_l$ is supported on $\omega_l$. For any function $f \in H^1(\Omega)$, its finite element approximation $f_h \in V_h$ can be written as
\[ f_h = \sum_{l=1}^{N_h} f(Q_l) \phi_l. \]
We further equip $V_h$ with the inner product $(f_h, g_h)_h : V_h \times V_h \rightarrow \mathbb{R}$ defined by
\[ (f_h, g_h)_h = \frac{1}{3} \sum_{l=1}^{N_h} |\omega_l| f_h(Q_l) g_h(Q_l), \forall f_h, g_h \in V_h. \]
The induced norm is defined as
\[ \|f_h\|_h = \sqrt{(f_h, f_h)}. \]

Because of the eventual smoothness of solutions to the inverse iteration (2.3) as shown in [37], our mixed finite-element method uses the space $V_h$ to approximate both the solution $u$ and its second-order partial derivatives $\partial^2 u / \partial x_i \partial x_j$ for $i, j = 1, 2$. In the rest of this section, we denote the finite element approximation of $v$ and $p$ by $v_h \in V_0h$ and $p_h \in (V_h)^{2 \times 2}$, respectively.

5.2. Finite element approximation of the three second-order partial derivatives. In equation (4.6), one needs to compute $D^2 v^{n+1/3}$, the Hessian of $v^{n+1/3}$, which will be numerically computed, and we adopt to our current setting the double regularization method introduced in [30].

The double regularization method is a two-step process to get a smooth approximation of $D^2 u$. In the first step, one solves
\[ \begin{cases} -\varepsilon_1 \nabla^2 \pi_{ij} + \pi_{ij} = \frac{\partial^2 u}{\partial x_i \partial x_j} & \text{in } \Omega, \\ \pi_{ij} = 0 & \text{on } \partial \Omega, \end{cases} \tag{5.1} \]
in which $\varepsilon_1 = O(h^2)$ is a constant, $\pi_{ij}$ is a regularized approximation of $\partial^2 u / \partial x_i \partial x_j$ with zero boundary condition. Although $\pi_{ij}$ is a smooth approximation, the zero boundary condition will have a disastrous influence to the solution $u$ of our scheme, as mentioned in [30]. To mitigate the influence, the second step is a correction step which solves
\[ \begin{cases} -\varepsilon_1 \nabla^2 D^2_{ij} u + D^2_{ij} u = \pi_{ij} & \text{in } \Omega, \\ \frac{\partial D^2_{ij} u}{\partial n} = 0 & \text{on } \partial \Omega, \end{cases} \tag{5.2} \]
where $n$ denotes the outward normal direction of $\partial \Omega$. The resulting $D^2_{ij} u$ is the doubly regularized approximation of $\partial^2 u / \partial x_i \partial x_j$.

From the divergence theorem, one has
\[ \begin{cases} \forall i, j = 1, 2, \forall v \in H^2(\Omega), \\ \int_{\Omega} \frac{\partial^2 v}{\partial x_i \partial x_j} w dx = -\frac{1}{2} \int_{\Omega} \left[ \frac{\partial v}{\partial x_i} \frac{\partial w}{\partial x_j} + \frac{\partial v}{\partial x_j} \frac{\partial w}{\partial x_i} \right] dx, \tag{5.3} \\ \forall w \in H^1_0(\Omega). \end{cases} \]
Based on equation (5.3), the discrete analogues of equations (5.1)-(5.2) read as:

\[
\begin{cases}
\pi_{ijh} \in V_{0h}, \\
c \sum_{T \in \omega_h} |T| \int_T \nabla \pi_{ijh} \cdot \nabla \phi_l dx + \frac{1}{3} |\omega_l| \pi_{ijh}(Q_l) = - \frac{1}{2} \int_{\omega_h} \left[ \frac{\partial u_x}{\partial x_i} \frac{\partial \phi_l}{\partial x_j} + \frac{\partial u_y}{\partial x_j} \frac{\partial \phi_l}{\partial x_i} \right] dx,
\end{cases}
\]

and

\[
\begin{cases}
D_{ijh}^2 u_h \in V_h, \\
c \sum_{T \in \omega_h} |T| \int_T \nabla D_{ijh}^2 u_h \cdot \nabla \phi_l dx + \frac{1}{3} |\omega_l| D_{ijh}^2 u_h(Q_l) = \frac{1}{3} |\omega_l| \pi_{ijh}(Q_l),
\end{cases}
\]

where \(c = O(1)\) is a constant.

5.3. **On the finite-element approximation of problem (4.5).** We first rewrite equation (4.5) in the variational form

\[
\begin{cases}
\psi^{n+1/3} \in V_{0h}, \\
\int_\Omega \psi^{n+1/3} \psi dx + \tau \int_\Omega (\epsilon I + \text{cof}(p^n)) \nabla \psi^{n+1/3} \cdot \nabla \psi dx = 2 \int_\Omega f^k \psi dx,
\end{cases}
\]

If \(p^n\) is semi-positive definite, then problem (5.6) admits a unique solution. Denote \(M = \epsilon I + \text{cof}(p^n)\). The discrete analogue of equation (5.6) reads as

\[
\begin{cases}
\psi^{n+1/3} \in V_{0h}, \\
\frac{1}{3} |\omega_l| \psi^{n+1/3}(Q_l) + \tau \sum_{m=1}^{N_0h} \int_{\omega_l \cap \omega_m} \mathbf{M} \nabla \phi_m \cdot \nabla \phi_l dx = \frac{2}{3} |\omega_l| f^k(Q_l),
\end{cases}
\]

Solving problem (5.7) is equivalent to solving a sparse linear system, for which many efficient solvers, such as the Cholesky decomposition, can be used.

5.4. **On the finite element approximation of problem (4.6).** We first define the projection operator \(P_+\) that projects \(2 \times 2\) real symmetric matrices to the set of real symmetric semi-positive definite matrices. Let \(A\) be a \(2 \times 2\) real symmetric matrix. By spectral decomposition, there exists a \(2 \times 2\) orthogonal matrix \(S\) so that \(A = SAS^{-1}\), where

\[
A = \begin{bmatrix}
\lambda_1 & 0 \\
0 & \lambda_2
\end{bmatrix}
\]

with \(\lambda_1, \lambda_2\) being eigenvalues of \(A\). If \(A\) is semi-positive definite, one has \(\lambda_1, \lambda_2 \geq 0\). Therefore we define \(P_+\) as

\[
P_+(A) = S \begin{bmatrix}
\max(\lambda_1, 0) & 0 \\
0 & \max(\lambda_2, 0)
\end{bmatrix} S^{-1}.
\]
In equation (4.6), we compute

\[ p_{h}^{n+1} = P_+ \left( e^{-\gamma \tau} p_{h}^n + (1 - e^{-\gamma \tau}) \begin{bmatrix} D_{11h}^{n+1/3} & D_{12h}^{n+1/3} \\ D_{21h}^{n+1/3} & D_{22h}^{n+1/3} \end{bmatrix} \right), \]

where the entries \( D_{i,jh}^{n+1/3} \) are computed using equations (5.4)-(5.5).

5.5. On the finite element approximation of problem (4.7). According to equation (4.11), we compute

\[ v_{h}^{n+1} = \frac{v_{h}^{n+1/3}}{\left( \sum_{l=1}^{N_0h} \frac{1}{3} |\omega_l| \left( v_{h}^{n+1/3}(Q_l) \right)^2 \right)^{1/2}}. \]

5.6. On the finite element approximation of equation (3.3). For any \( u_h \in H^1_0(\Omega) \), the discrete analogue of equation (3.3) reads as

\[ R(u_h) = -\frac{N_0h}{2} \sum_{m=1}^{N_0h} \frac{2}{3} |\omega_m| (-u_h(Q_m))^3 \]

where \( D_{h}^2 u_h \) is the finite-element approximation of \( D^2 u \) computed using equations (5.4) and (5.5).

Note that if \( u \) is an eigenfunction of the Monge–Ampère equation (1.2), by Theorem 2.2, one can compute the eigenvalue as \( \lambda_{MA} = \inf_{u \in \mathcal{K}} R(u) \). Therefore, for every time step, we can compute the approximate ‘eigenvalue’ corresponding to \( u_k^h \) as

\[ \lambda_k^h = R(u_k^h) \]

and monitor the evolution of \( \lambda_k^h \), which will monotonically converge to \( \lambda_{MA} \) as shown in [37].

5.7. On the finite element approximation of the initial condition. Denote the finite element of \( u_0 \) and \( (v_0, p_0) \) by \( u_{0h} \) and \( (v_{0h}, p_{0h}) \), respectively. The discrete analogue of the initial condition (4.13) reads as

\[
\begin{cases}
    u_{0h} \in V_{0h}, \\
    \sum_{m=1}^{N_0h} u_{0h}(Q_m) \int_{\omega_l \cap \omega_m} (\text{cof}(D_{h}^2 u_{0h}(Q_m)) \cdot \nabla \phi_m) \cdot \nabla \phi_l \, dx = \frac{2}{3} \eta |\omega_l|, \\
    \forall l = 1, \ldots, N_0h.
\end{cases}
\]

For \( (v_{0h}, p_{0h}) \), we set

\[ v_{0h} = u_k^h, \quad p_{0h} = D_{h}^2 v_{0h}, \]

where \( D_{h}^2 \) is the double regularization approximation using equations (5.4)-(5.5).

6. Numerical Experiments

We demonstrate the efficiency of scheme (4.5)-(4.7) by several numerical experiments. We set the stopping criterion as \( \|u_{h}^{k+1} - u_{h}^{k}\|_h < \xi \) for some small \( \xi > 0 \). Without specification, in
We set $\xi = 10^{-6}$, $\epsilon = 2h^2$, and $c = 2$, where $\epsilon$ and $c$ are regularization parameters in equation (3.8) and scheme (5.4)-(5.5), respectively.

When the exact solution, denoted by $u^*_h$, is given, we define the $L^2$ error and $L^\infty$ error of $u_h$ as

$$\|u_h - u^*_h\|_h \quad \text{and} \quad \max_{m} |u_h(Q_m) - u^*_h(Q_m)|,$$

respectively.

Algorithm 1 consists of two iterations: the outer iteration for $u$ and the inner iteration for $v$ and $p$. Since both $u$ and $v$ are estimates of the solution of equation (1.2), it is not necessary to solve every inner iteration until steady state. Instead, one can just solve the inner iteration for a few steps. In our experiments, we observe that just 1 iteration step for the inner iteration is sufficient for our algorithm to converge. Thus in all of our experiments, we solve the inner iteration for only 1 step in each outer iteration.

6.1. Example 1. In the first example, we test our algorithm on the unit disk

$$\Omega = \{(x_1, x_2) : x_1^2 + x_2^2 < 1\}.$$  \hfill (6.1)

The triangulation of the domain with $h = 1/20$ is visualized in Figure 1(a).

In this case, equation (1.2) has a radial solution. Let $r = \sqrt{x_1^2 + x_2^2}$. For a radial function $g(r)$, one has $\det D^2 g = \frac{g''}{r}$. Therefore, we write the solution to equation (1.2) as $u(r)$, which
satisfies

\[
\begin{cases}
  u \leq 0, \lambda > 0, \\
  u' u'' = -\lambda r u^2 \text{ in } (0,1), \\
  u'(0) = 0, u(1) = 0, \\
  2\pi \int_0^1 |u|^2 r dr = 1.
\end{cases}
\]  

Using a shooting method, we can solve the ODE problem (6.2) very accurately. The exact solution verifies \(u(0) \approx -1.0628\) and \(\lambda \approx 7.4897\). On the domain (6.1), we test our algorithm with \(h = 1/20, 1/40, 1/80\) and \(1/160\). In Figure 2(a)–(d), we show results with \(h = 1/80\). Our numerical result is visualized in Figure 2(a). The contour of Figure 2(a) is shown in Figure 2(b). Our result is a smooth radial function, whose contour consists of several circles with the same center. The convergence histories of the error \(\|u_h^{k+1} - u_h^k\|\) and the computed eigenvalue are shown in Figure 2(c) and Figure 2(d), respectively. Linear convergence is observed for the error \(\|u_h^{k+1} - u_h^k\|\), and the convergence rate is approximately 0.47. The computed eigenvalue converges with just 5 iterations. In Figure 2(e), we show the cross sections of the results with various \(h\) along \(x_2 = 0\). As \(h\) goes to 0, our computed solution converges to the exact solution. For better visualization, the zoomed bottom region of Figure 2(e) is shown in Figure 2(f).

To quantify the convergence of the proposed algorithm, we present in Table 1 the number of iterations needed for convergence, \(L^2\)- and \(L^\infty\)-errors, computed eigenvalues and the minimal value of the computed solution with various \(h\). For all resolutions of mesh, 13 iterations are sufficient for the algorithm to converge. As \(h\) goes to zero, the convergence rate of the \(L^2\)- and \(L^\infty\)-error goes to 1, and the computed eigenvalue and the minimal value converge to the exact solutions. The eigenvalue \(\lambda_h\) converges linearly to the exact eigenvalue with an error of \(O(h)\).

We next compare Algorithm 1 with the method proposed in [28]. For the method from [28], we have to use small time steps to make sure that the method does converge. In the numerical experiment, we set the time step as \(h/2\) and stopping criterion as \(10^{-6}\). Note that the method from [28] finds the solution of equation (1.2) with \(\|u_h\|_3 = 1\). When computing the \(L^2\)- and \(L^\infty\)-errors, we first normalize the solution so that \(\|u_h\|_2 = 1\) and then compute the errors. The comparisons are shown in Table 2. For both \(L^2\)- and \(L^\infty\)-errors, both algorithms have errors with similar magnitudes. We compare the computational efficiency between the two algorithms in Table 3. The number of iterations used by Algorithm 1 is independent of the mesh resolution, while the number of iterations used by [28] grows approximately linearly with \(1/h\). For the CPU time, Algorithm (1) is also much faster than the method in [28]. Note that in Algorithm (1), the constraint \(\|u_h\|_2 = 1\) is enforced by the projection step (4.7). In [28], the constraint is \(\|u_h\|_3 = 1\), which was enforced by a sequential quadratic programming algorithm, which in turn uses around 15 iterations in each outer iteration.

6.2. Example 2. In the second example, we consider the convex smoothed square domain

\[
\Omega = \left\{ (x_1, x_2) : |x_1|^{2.5} + |x_2|^{2.5} < 1 \right\}.
\]  

The triangulation of the domain with \(h = 1/20\) is visualized in Figure 1(b), which has a shape between the unit disk and a square. We test our algorithm with \(h\) varying from \(h = 1/20\) to \(h = 1/160\). Similar to our settings in the previous example, we set the stopping criterion \(\xi = 10^{-6}\). The time step is set as \(\tau = 1/2\). Our results with \(h = 1/80\) are visualized in Figure...
(a) (b) (c)

FIGURE 2. The unit disk domain (6.1). (a) The computed result with $h = 1/80$. (b) The contour of (a). (c) The history of the error $\Vert u_{h}^{k+1} - u_{h}^{k} \Vert_{h}$ with $h = 1/80$. (d) The history of the computed eigenvalue $\lambda_{h}^{k}$ with $h = 1/80$. (e) Comparison of the cross sections along $x_{2} = 0$ of the computed solution with various $h$. (f) Zoomed plot of the bottom region of (e).

TABLE 1. The unit disk domain (6.1). Variations with $h$ of the number of iterations necessary to achieve convergence (2nd column), of the $L^2$ and $L^{\infty}$ approximation errors and of the associated convergence rates (columns 4, 5, 6 and 7), of the computed eigenvalue (8th column) and of the minimal value of $u_{h}$ over $\Omega$ (that is $u_{h}(0)$) (9th column). The exact eigenvalue is around 7.4897. The minimal value of the exact solution is around $-1.0628$.

| $h$     | # Iter. | $\Vert u_{h}^{k+1} - u_{h}^{k} \Vert_{h}$ | $L^2$-error | rate     | $L^{\infty}$-error | rate     | $\lambda_{h}$ | min $u_{h}$ |
|---------|---------|------------------------------------------|-------------|----------|---------------------|----------|----------------|-------------|
| 1/20    | 13      | $2.13 \times 10^{-7}$                    | $4.91 \times 10^{-2}$ |          | $4.29 \times 10^{-2}$ |          | 5.9716         | -1.0189     |
| 1/40    | 13      | $2.91 \times 10^{-7}$                    | $3.36 \times 10^{-2}$ | 0.54     | $3.04 \times 10^{-2}$ | 0.50     | 6.6656         | -1.0362     |
| 1/80    | 13      | $3.56 \times 10^{-7}$                    | $1.94 \times 10^{-2}$ | 0.79     | $1.86 \times 10^{-2}$ | 0.71     | 7.0655         | -1.0484     |
| 1/160   | 13      | $4.04 \times 10^{-7}$                    | $1.01 \times 10^{-2}$ | 0.94     | $1.03 \times 10^{-2}$ | 0.85     | 7.2816         | -1.0556     |

3(a)–(d). Our computed solution is shown in Figure 3(a), whose contour is shown in Figure 3(b). Again, our solution is very smooth. The convergence histories of the error $\Vert u_{h}^{k+1} - u_{h}^{k} \Vert_{h}$ and the computed eigenvalues $\lambda_{h}^{k}$ are shown in Figure 3(c) and Figure 3(d), respectively. The error $\Vert u_{h}^{k+1} - u_{h}^{k} \Vert_{h}$ converges linearly with a rate of 0.38. In this numerical experiment, the stopping criterion is satisfied after 16 iterations. The computed eigenvalue achieves its steady state with about 5 iterations. With various $h$, the comparison of cross sections of our results along $x_{2} = 0$ is shown in Figure 3(e)–(f). As $h$ goes to 0, the convergence of the solution along cross sections is observed.
TABLE 2. The unit disk domain (6.1). Variations with $h$ of the number of iterations necessary to achieve convergence (2nd column), of the $L^2$ and $L^\infty$ approximation errors and of the associated convergence rates (columns 4, 5, 6 and 7), of the computed eigenvalue (8th column) and of the minimal value of $u_h$ over $\Omega$ (that is $u_h(0)$) (9th column). The exact eigenvalue is around 7.4897. The minimal value of the exact solution is around $-1.0628$.

| $h$  | $L^2$-error | rate | $L^\infty$-error | rate | $L^2$-error | rate | $L^\infty$-error | rate |
|------|-------------|------|-----------------|------|-------------|------|-----------------|------|
| 1/20 | 4.91 $\times$ 10^{-2} |  | 4.29 $\times$ 10^{-2} |  | 4.01 $\times$ 10^{-2} |  | 8.40 $\times$ 10^{-2} |  |
| 1/40 | 3.36 $\times$ 10^{-2} | 0.54 | 3.04 $\times$ 10^{-2} | 0.50 | 2.33 $\times$ 10^{-2} | 0.78 | 4.00 $\times$ 10^{-2} | 1.07 |
| 1/80 | 1.94 $\times$ 10^{-2} | 0.79 | 1.86 $\times$ 10^{-2} | 0.71 | 1.37 $\times$ 10^{-2} | 0.76 | 2.05 $\times$ 10^{-2} | 0.96 |
| 1/160| 1.01 $\times$ 10^{-2} | 0.94 | 1.03 $\times$ 10^{-2} | 0.85 | 7.55 $\times$ 10^{-3} | 0.86 | 1.08 $\times$ 10^{-2} | 0.92 |

TABLE 3. The unit disk domain (6.1). Comparison of the number of iterations and the CPU time needed by Algorithm 1 and the method in [28] for convergence.

| $h$  | # Iter. | CPU time | # Iter. | CPU time |
|------|--------|----------|--------|----------|
| 1/20 | 13     | 1.44     | 62     | 3.55     |
| 1/40 | 13     | 4.58     | 101    | 22.39    |
| 1/80 | 13     | 18.35    | 151    | 138.47   |
| 1/160| 13     | 83.95    | 263    | 1206.96  |

TABLE 4. The smoothed square domain (6.3). Variations with $h$ of the number of iterations necessary to achieve convergence (2nd column), of the computed eigenvalue (4th column) and of the minimal value of $u_h$ over $\Omega$ (that is $u_h(0)$) (5th column).

| $h$  | # Iter. | $\|u_h^{k+1} - u_h^k\|_h$ | $\lambda_h$ | min $u_h$ |
|------|--------|--------------------------|-------------|----------|
| 1/20 | 14     | 6.05 $\times$ 10^{-7}    | 5.17        | -0.9833  |
| 1/40 | 14     | 8.00 $\times$ 10^{-7}    | 5.72        | -0.9982  |
| 1/80 | 16     | 2.08 $\times$ 10^{-7}    | 6.05        | -1.0094  |
| 1/160| 18     | 7.77 $\times$ 10^{-7}    | 6.22        | -1.0159  |

We then report the computational cost and convergence behavior of the computed eigenvalue and minimal value with various $h$ in Table 4. The convergence of the eigenvalue is similar to that in [28]: the eigenvalue $\lambda_h$ converges to $\lambda$ uniformly in the rate $\lambda_h \approx \lambda - ch$ with $\lambda \approx 6.4, c \approx 26$. In terms of the computational cost, Algorithm 1 is very efficient since all experiments used less than 20 iterations to satisfy the stopping criterion.

6.3. Example 3. In the third example, we consider an ellipse domain defined by

$$\Omega = \{(x_1, x_2) : x_1^2 + 2x_2^2 < 1\}.$$ (6.4)

A triangulation of the domain with $h = 1/20$ is visualized in Figure 1(c). In this set of experiments, we set stopping criterion $\xi = 10^{-6}$ and time step $\tau = 1/2$. The results with $h = 1/80$ are
Figure 3. The smoothed square domain (6.3). (a) The computed result with $h = 1/80$. (b) The contour of (a). (c) The convergence history of the errors $\|u_{k+1}^h - u_k^h\|_h$. (d) The history of the computed eigenvalue $\lambda_h^k$. (e) Comparison of the cross sections along $x_2 = 0$ of the computed solution with various $h$. (f) Zoomed plot of the bottom region of (e).

Table 5. The ellipse domain (6.4). Variations with $h$ of the number of iterations necessary to achieve convergence (2nd column), of the computed eigenvalue (4th column) and of the minimal value of $u_h$ over $\Omega$ (that is $u_h(0)$) (5th column).

| $h$        | # Iter. | $\|u_{k+1}^h - u_k^h\|_h$ | $\lambda_h^k$ | $\min u_h$ |
|------------|---------|---------------------------|----------------|-----------|
| 1/20       | 16      | $6.80 \times 10^{-7}$     | 21.55          | -1.4277   |
| 1/40       | 16      | $9.68 \times 10^{-7}$     | 25.18          | -1.4525   |
| 1/80       | 17      | $6.44 \times 10^{-7}$     | 27.41          | -1.4734   |
| 1/160      | 17      | $7.00 \times 10^{-7}$     | 28.67          | -1.4875   |

shown in Figure 4(a)–(d). Similar to the results in the previous examples, the computed solution is smooth, and its contour consists of several ellipses with the same center, as shown in Figure 4(a) and Figure 4(b), respectively. In Figure 4(c), linear convergence is observed for the error $\|u_{k+1}^h - u_k^h\|_h$, and the convergence rate is about 0.34. The computed eigenvalue $\lambda_h^k$ attains its steady state with 6 iterations. With various $h$, we compare in Figure 4(e)–(f) the cross sections of the computed results along $x_2 = 0$. Convergence is observed as $h$ goes to 0.

With various $h$, the computational cost, the computed eigenvalue and minimal value of the computed solution are presented in Table 5. The eigenvalue $\lambda_h$ converges to $\lambda$ uniformly in the rate $\lambda_h \approx \lambda - ch$ with $\lambda \approx 29.5, c \approx 161$. In terms of the computational cost, all experiments used less than 20 iterations to satisfy the stopping criterion.
6.4. **Example 4.** We conclude this section by considering an open convex domain with a non-smooth boundary:

\[
\Omega = \{ (x_1,x_2) : -x_1(1-x_1) < x_2 < x_1(1-x_1), \ 0 < x_1 < 1 \}.
\]  

(6.5)

The domain described in the set (6.5) has an eye shape, and its triangulation with \( h = 1/40 \) is visualized in Figure 1(d). Since the domain is not smooth, in our experiments we use a smaller time step \( \tau = 1/8 \) and larger regularization parameters \( \varepsilon = 4h^2 \) and \( c = 4 \). We set stopping criterion \( \xi = 10^{-6} \). The results with \( h = 1/160 \) are shown in Figure 5(a)–(d). The computed solution is smooth, and its level curves have the same center, as shown in Figure 5(a) and Figure 5(b), respectively. In Figure 5(c), linear convergence is observed for the error \( \| u_h^{k+1} - u_h^k \|_h \). The computed eigenvalue \( \lambda_h^k \) attains its steady state with 7 iterations. With various \( h \), we compare in Figure 5(e)–(f) the cross sections of the computed results along \( x_2 = 0 \). Convergence is observed as \( h \) goes to 0.

With various \( h \), the computational cost, the computed eigenvalue, and the minimal value of the computed solution are presented in Table 6. The eigenvalue \( \lambda_h \) converges to \( \lambda \) uniformly in the rate \( \lambda_h \approx \lambda - ch \) with \( \lambda \approx 618, c \approx 7792.3 \). In terms of the computational cost, all experiments used no more than 30 iterations to satisfy the stopping criterion.

---

**Figure 4.** The ellipse domain (6.4). (a) The computed result with \( h = 1/80 \). (b) The contour of (a). (c) The history of the error \( \| u_h^{k+1} - u_h^k \|_h \). (d) The history of the computed eigenvalue \( \lambda_h^k \). (e) Comparison of the cross sections along \( x_2 = 0 \) of the computed solution with various \( h \). (f) Zoomed plot of the bottom region of (e).
FIGURE 5. The eye–shape domain (6.5). (a) The computed result with $h = 1/160$. (b) The contour of (a). (c) The history of the error $\|u_{h}^{k+1} - u_{h}^{k}\|_{h}$. (d) The history of the computed eigenvalue $\hat{\lambda}_{h}^{k}$. (e) Comparison of the cross sections along $x_2 = 0$ of the computed solution with various $h$. (f) Zoomed plot of the bottom region of (e).

TABLE 6. The eye–shape domain (6.5). Variations with $h$ of the number of iterations necessary to achieve convergence (2nd column), of the computed eigenvalue (4th column) and of the minimal value of $u_{h}$ over $\Omega$ (that is $u_{h}(0)$) (5th column).

| $h$   | # Iter. | $\|u_{h}^{k+1} - u_{h}^{k}\|_{h}$ | $\hat{\lambda}_{h}$ | min $u_{h}$ |
|-------|---------|-----------------------------------|----------------------|-------------|
| 1/40  | 15      | $7.80 \times 10^{-7}$             | 425.51               | -3.1091     |
| 1/80  | 20      | $7.53 \times 10^{-7}$             | 516.57               | -3.1256     |
| 1/160 | 27      | $8.35 \times 10^{-7}$             | 568.47               | -3.1617     |
| 1/320 | 30      | $7.87 \times 10^{-7}$             | 597.39               | -3.1913     |

7. Conclusion

We proposed an efficient operator–splitting method to solve the eigenvalue problem of the Monge–Ampère equation. The backbone of our method relies on a convergent algorithm proposed in [1]. In each iteration, we solve a constrained optimization problem whose optimality condition is of the Monge–Ampère type. We remove the constraint by including an indicator function and decouple the nonlinearity by introducing an auxiliary variable. The resulting problem is then converted to finding the steady state solution of an initial value problem which is time discretized by an operator–splitting method. The efficiency and effectiveness of the proposed method is demonstrated with several numerical experiments. In our experiments, we can
choose a large constant time step. On smooth convex domains, our algorithm converges with a few iterations and is much faster than existing methods.

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The preparation of this manuscript has been overshadowed by Roland’s passing away in January 2022. Roland and the authors had intended to write jointly: most of the main ideas were worked out together and the authors have done their best to complete them. In sorrow, the authors dedicate this work to his memory. Roland’s creativity, generosity, and friendship will be remembered.

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