# The Derivative free trust-region method for the inverse elastography problem 

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#### Abstract

In this work we investigate a mathematical model to reconstruct the mechanical properties of an elastic medium, for the optical coherence elastography imaging modality. To this end, we start by considering a mathematical model for the mechanical deformation based on time-harmonic equations of linear elasticity. The mathematical model for solving this direct problem is the computational basis to address the inverse problem which consists of determining the set of parameters that characterize the mechanical properties of the medium knowing the displacement field for a given excitation. We formulate the inverse problem as PDE-constrained optimization problem, where the objective function measures the discrepancy between observations and predictions. We propose a derivative free trust-region method to solve this inverse problem and we report several computational results which illustrate its behavior in terms of accuracy and efficiency.


Keywords: derivative free trust-region method, inverse problem, linear elasticity, mechanical properties reconstruction

## 1 Linear elasticity model

Let us consider an isotropic elastic material in the configuration space $\Omega \subseteq \mathbb{R}^{3}$, where $\Omega$ is a polyhedron with boundary $\partial \Omega$. The aim is to characterize the field of induced displacements, $\mathbf{u}(x, t)$ with $x \in \Omega$ and $t \in \mathbb{R}_{0}^{+}$.

Let us consider a sinusoidal excitation. The displacement has a time-harmonic form given by [7],

$$
\begin{equation*}
\mathbf{u}(x, t)=\Re\left(\mathbf{u}(x) e^{i \omega t}\right), \tag{1}
\end{equation*}
$$

where $\Re$ is the real part of a complex and $\omega$ is the angular frequency of the sinusoidal excitation. For time-harmonic elastic propagation, the elastic displacement field $\mathbf{u}$ satisfies the Lamé equation

$$
\begin{equation*}
\mu \nabla^{2} \mathbf{u}+(\lambda+\mu) \nabla(\nabla \cdot \mathbf{u})+\omega^{2} \rho \mathbf{u}+\mathbf{f}=0 \text { in } \Omega \tag{2}
\end{equation*}
$$

where $\rho$ is the material density, $\mathbf{f}$ is a given distribution of body forces and the Lamé constants $\mu$ and $\lambda$ are given, respectively, by

$$
\mu=\frac{E}{2(1+\nu)} \text { and } \lambda=\frac{\nu E}{(1+\nu)(1-2 \nu)},
$$

being $E$ is the Young's Modulus and $\nu$ is the Poisson's ratio.
Let $\Gamma_{1}$ and $\Gamma_{2}$ be two open subsets of $\partial \Omega$ such that $\partial \Omega=\bar{\Gamma}_{1} \cup \bar{\Gamma}_{2}, \Gamma_{1} \cap \Gamma_{2}=\emptyset$ and meas $\left(\Gamma_{2}\right)>0$. We impose the traction boundary condition on $\Gamma_{1}$

$$
\begin{equation*}
\sigma(\mathbf{u}) \eta=\mathbf{g} \text { on } \Gamma_{1}, \tag{3}
\end{equation*}
$$

where $\eta$ is the unit outer normal direction. The stress tensor is given by

$$
\sigma(\mathbf{u})=2 \mu \varepsilon(\mathbf{u})+\lambda \operatorname{tr}(\varepsilon(\mathbf{u})) I
$$

where $\varepsilon$ is the strain tensor

$$
\varepsilon(\mathbf{u})=\frac{1}{2}\left(\nabla \mathbf{u}+(\nabla \mathbf{u})^{\top}\right)
$$

Here I is the $3 \times 3$ identity matrix and $\operatorname{tr}(\varepsilon(\mathbf{u}))$ is the trace of $\varepsilon(\mathbf{u})$. It is assumed that the medium is fixed on some non-empty open set $\Gamma_{2}$, where we impose displacement boundary condition

$$
\begin{equation*}
\mathbf{u}=0 \text { on } \Gamma_{2} . \tag{4}
\end{equation*}
$$

To solve this direct problem, we consider the classic finite element method (FEM) which consists of using degree one piecewise polynomials in the approximation. For the case of a nearly incompressible materials, i.e. with Poisson's ratio $\nu$ close to 0.5 , the performance of a classical FEM scheme can deteriorate due to the locking as $\nu \rightarrow 0.5$ [1]. Here we are assuming that we are dealing with media for which the range of values of the Poisson's ratio leads to locking-free FEM solutions. As an example of application, we can mention the aortic elastography [2]. For general materials, some numerical methods have been proposed in the literature, in particular some variations of mixed finite element methods.

To derive the finite element method, we need to consider the weak form of the mathematical model. Let $V=\left\{\mathbf{v} \in H^{1}(\Omega):\left.\mathbf{v}\right|_{\Gamma_{2}}=0\right\}$. The weak formulation of (2)-(4) reads: find $\mathbf{u} \in V$ such that

$$
\begin{equation*}
a(\mathbf{u}, \mathbf{v})=l(\mathbf{v}), \forall \mathbf{v} \in V \tag{5}
\end{equation*}
$$

where

$$
\begin{equation*}
a(\mathbf{u}, \mathbf{v})=\int_{\Omega} 2 \mu \varepsilon(\mathbf{u}): \varepsilon(\mathbf{v})+\lambda(\nabla \cdot \mathbf{u})(\nabla \cdot \mathbf{v})-\omega^{2} \rho \mathbf{u} \cdot \mathbf{v} d x \tag{6}
\end{equation*}
$$

and

$$
l(\mathbf{v})=\int_{\Gamma_{1}} \mathbf{g} \cdot \mathbf{v} d s+\int_{\Omega} \mathbf{f} \cdot \mathbf{v} d x
$$

Let us consider a partition of $\Omega$ into $M$ tetrahedra $K_{j}, j \in\{1, \ldots, M\}$ so that

$$
\begin{equation*}
\Omega=\bigcup_{j=1}^{M} K_{j} \text { and } \operatorname{int}\left(K_{i}\right) \cap \operatorname{int}\left(K_{j}\right)=\emptyset, \forall i, j \in\{1, \ldots, M\}, i \neq j \tag{7}
\end{equation*}
$$

The resulting subdivision (or mesh) is denoted by $\Omega_{h}$ where $h$ represents the diameter of the partition. To each tetrahedron there are associated four vertices
that can be either in the interior or on the border of $\Omega$. For any pair of open tetrahedra in the partition $K_{i}$ and $K_{j}, i \neq j, \bar{K}_{i} \cap \bar{K}_{j}$ is either empty, or a common vertex, side or face of $K_{i}$ and $K_{j}$.

Let us consider the finite dimensional subspace $V_{h} \subset V$ of continuous functions which are linear on each tetrahedron. Assuming that $N$ is the total number of vertices in $\Omega_{h}$ then $\operatorname{dim} V_{h}=3 N$. The finite element formulation of the problem (5) can be written as: find $\mathbf{u}_{h} \in V_{h}$ such that

$$
\begin{equation*}
a\left(\mathbf{u}_{h}, \mathbf{v}_{h}\right)=l\left(\mathbf{v}_{h}\right), \forall \mathbf{v}_{h} \in V_{h} \tag{8}
\end{equation*}
$$

Let $V_{h}=\operatorname{span}\left\{\phi_{11}, \ldots, \phi_{N 1}, \phi_{12}, \ldots, \phi_{N 2}, \phi_{13}, \ldots, \phi_{N 3}\right\}$, where $\phi_{j i}, i=\{1,2,3\}$, $j=\{1, \ldots, N\}$, are the linearly independent basis functions, $\phi_{j i}\left(x^{j}\right)=1, \phi_{j i}\left(x^{k}\right)=$ $0(k \neq j)$, and the support of $\phi_{j i}$ consists in all tetrahedra that share $x^{j}$ as a vertex. In this way, each component of the approximate solution $\mathbf{u}_{h}=$ $\left(u_{1 h}, u_{2 h}, u_{3 h}\right) \in V_{h}$ can be written as a linear combination of the basis functions $\phi_{j i}$ with

$$
\begin{equation*}
u_{i h}(x)=\sum_{j=1}^{N} U_{j i} \phi_{j i}(x), i=1,2,3, \tag{9}
\end{equation*}
$$

where $U_{j i}, i=\{1,2,3\}, j=\{1, \ldots, N\}$ are the coefficients that we want to calculate. Problem (8) is then equivalent to a linear system in the form

$$
\begin{equation*}
A U=F \tag{10}
\end{equation*}
$$

where $U=\left[U_{11}, \ldots, U_{N 1}, U_{12}, \ldots, U_{N 2}, U_{13}, \ldots, U_{N 3}\right]^{\top}, A$ is a $3 N \times 3 N$ matrix which depends on the parameters $\mu$ and $\lambda$, and $F$ a vector of dimension $3 N$. $A$ is a non-singular matrix so (10) has a unique solution.

## 2 Inverse problem

In this section we will analyze the inverse problem, which can be described by the following minimization program:

$$
\begin{array}{cc}
\min _{\mu, \lambda}\left\|U-U_{o b s}\right\|_{L_{h}^{2}(\Omega)}^{2} \\
\text { s.t. } & A U=F \\
& \mu \in\left[\mu_{1}, \mu_{2}\right] \\
& \lambda \in\left[\lambda_{1}, \lambda_{2}\right]
\end{array}
$$

which can be rewritten as

$$
\begin{array}{ll}
\min _{\mu, \lambda} & \left\|A^{-1} F-U_{o b s}\right\|_{L_{h}^{2}(\Omega)}^{2} \\
\text { s.t. } & \mu \in\left[\mu_{1}, \mu_{2}\right]  \tag{11}\\
& \lambda \in\left[\lambda_{1}, \lambda_{2}\right]
\end{array}
$$

where $A$ and $F$ define the linear system to solve the direct problem (10) and $U_{o b s}$ is the vector that contains the information of the given data on the vertices
of the mesh. Here we use the discrete $L_{h}^{2}$-norm, defined for any $3 N \times 1$ vector $y$, as

$$
\|y\|_{L_{h}^{2}(\Omega)}^{2}=\sum_{K \in \Omega_{h}}\|y\|_{L_{h}^{2}(K)}^{2},
$$

with

$$
\|y\|_{L_{h}^{2}(K)}^{2}=\frac{|K|}{4} \sum_{i=1}^{4} \sum_{j=0}^{2} y_{t\left(r_{i}\right)+j N}^{2}
$$

where $|K|$ denotes the volume of the tetrahedron $K$ with vertices $r_{i}, i \in\{1, \ldots, 4\}$. The function $t$ is defined by

$$
\begin{aligned}
t: \mathbb{R}^{3} & \rightarrow\{1, \ldots, N\} \\
r_{i} & \mapsto t\left(r_{i}\right),
\end{aligned}
$$

where $t\left(r_{i}\right)$ is the index that corresponds to vertex of $r_{i}$ in the global numbering.
We are assuming that the range of values for $\mu$ and $\lambda$ set in (11) are compatible with the biological structures.

For convenience, in what follows, we denote the objective function by $l(\mu, \lambda)$, that is,

$$
\begin{equation*}
l(\mu, \lambda)=\left\|A^{-1} F-U_{o b s}\right\|_{L_{h}^{2}(\Omega)}^{2} \tag{12}
\end{equation*}
$$

## 3 Derivative free trust-region method

The configuration of the function $l$ seems to be well approximate, locally, by a quadratic function. This fact, motivated us to apply a variation version of the trust region method presented in [15] where a quadratic model is used to approximate the objective function.

Usually trust region method is based on Taylor's series expansion of $l$ around $\left(\mu_{k}, \lambda_{k}\right)$ [10]. To avoid the computation of the derivatives of $l$, we consider the least squares method to get the quadratic model which approximates $l(\mu, \lambda)$ using a set

$$
P_{k}=\left\{\left(\mu_{i, k}, \lambda_{i, k}\right): i \in\left\{1, \ldots, n_{k}\right\}\right\}
$$

with $n_{k} \geq 9$ random points where the function $l$ is known. Let us consider that these $n_{k}$ points belong to the set $I_{k}=\left[\mu_{i n f, k}, \mu_{\text {sup }, k}\right] \times\left[\lambda_{i n f, k}, \lambda_{\text {sup }, k}\right]$ and that, each of the nine sub-sets, which correspond to each of the cells in a uniform $3 \times 3$ grid in the set $I_{k}$ :

$$
\begin{equation*}
I_{i, j, k}=\left[\mu_{i n f, k}+(i-1) s_{k}, \mu_{i n f, k}+i s_{k}\right] \times\left[\lambda_{i n f, k}+(j-1) r_{k}, \lambda_{i n f, k}+j r_{k}\right], \tag{13}
\end{equation*}
$$

$i, j=1,2,3$, where $s_{k}=\frac{1}{3}\left(\mu_{\text {sup }, k}-\mu_{\text {inf }, k}\right), r_{k}=\frac{1}{3}\left(\lambda_{\text {sup }, k}-\lambda_{\text {inf }, k}\right)$. We impose that each sub-sets contains at least one point to make sure that the points are not too close to each other and to have representative points around all the trust region. Moreover, we assume that $\left(\mu_{k}, \lambda_{k}\right)$ belongs to the central sub-set, that is, $\left(\mu_{k}, \lambda_{k}\right) \in I_{2,2, k}$.

With these $n_{k}$ points, we want to derive the $k$-th quadratic model $l_{k}$, that can be written in the form

$$
\begin{equation*}
l_{k}(\mu, \lambda)=a_{1}+a_{2} \mu+a_{3} \lambda+a_{4} \mu^{2}+a_{5} \lambda^{2}+a_{6} \mu \lambda, k \in \mathbb{Z}_{0}^{+}, \tag{14}
\end{equation*}
$$

for some $a_{i} \in \mathbb{R}, i \in\{1, \ldots, 6\}$. The coefficients $a_{i}, i \in\{1, \ldots, 6\}$, are determined in way to minimize the function

$$
L_{k}\left(a_{1}, \ldots, a_{6}\right)=\sum_{i=1}^{n_{k}}\left(l\left(\mu_{i, k}, \lambda_{i, k}\right)-l_{k}\left(\mu_{i, k}, \lambda_{i, k}\right)\right)^{2}
$$

The necessary condition to obtain the solution is

$$
\begin{equation*}
\frac{\partial L_{k}}{\partial a_{i}}\left(a_{1}, \ldots, a_{6}\right)=0, i \in\{1, \ldots, 6\} \tag{15}
\end{equation*}
$$

and the solution $a^{*}=\left(a_{1}^{*}, \ldots, a_{6}^{*}\right)$ of this system exists and it is unique [6].
It now remains to characterize the minimizer of $l_{k}$ through the trust region. Let $\left(\mu_{0}, \lambda_{0}\right)$ be the initial point and $\Delta_{0}>0$ the initial trust region radius. So the initial trust-region is defined by the set $I_{0}$ such that $\left(\mu_{0}, \lambda_{0}\right) \in I_{2,2,0}$ and for the other sub-sets a random point is generated. After defining the set $P_{0}$ with these points, we are able to calculate $l_{0}$. For the iteration $k$ let us consider the points of $P_{k}$ to obtain the $k$-th approximation $l_{k}$ by the least squares method. To achieve the minimizer, we will obtain the solution of the next sub-problem

$$
\begin{array}{lc}
\min _{(\mu, \lambda) \in \mathbb{R}^{2}} & l_{k}(\mu, \lambda)  \tag{16}\\
\text { s.t. }
\end{array}\left\|(\mu, \lambda)-\left(\mu_{k}, \lambda_{k}\right)\right\|_{\infty} \leq \Delta_{k} .
$$

The solution $\left(\mu_{k+1}, \lambda_{k+1}\right)$ will be the minimizer of $l_{k}$ in the square centered in $\left(\mu_{k}, \lambda_{k}\right)$ and radius $\Delta_{k}$. The critical point is the solution of $\nabla l_{k}(\mu, \lambda)=0$, that is,

$$
\left[\begin{array}{cc}
2 a_{4}^{*} & a_{6}^{*}  \tag{17}\\
a_{6}^{*} & 2 a_{5}^{*}
\end{array}\right]\left[\begin{array}{c}
\mu \\
\lambda
\end{array}\right]=-\left[\begin{array}{l}
a_{2}^{*} \\
a_{3}^{*}
\end{array}\right]
$$

where $a_{i}^{*}, i \in\{2, \cdots, 6\}$ are the solution of (15). We note that the matrix of the system (17) is the hessian, $\nabla^{2} l_{k}(\mu, \lambda)$, of the quadratic model (14). The correspondent eigenvalues are given by

$$
a_{4}^{*}+a_{5}^{*} \pm \sqrt{\left(a_{4}^{*}-a_{5}^{*}\right)^{2}+\left(a_{6}^{*}\right)^{2}} .
$$

So this values are both positive when $4 a_{4}^{*} a_{5}^{*}-\left(a_{6}^{*}\right)^{2}>0$ and $a_{4}^{*}>0$.
If the hessian is positive definite and the solution obtained by resolution of (17) is in the trust region, the solution of (17) is the minimizer of (16). If one of this conditions is not satisfied the minimizer will be determined on the boundary.

Let $\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)$ denote the minimizer of (16) in the iteration $k$. To accept this minimizer, we analyze the ratio

$$
\begin{equation*}
\rho_{k}=\frac{l\left(\mu_{k}, \lambda_{k}\right)-l\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)}{l_{k}\left(\mu_{k}, \lambda_{k}\right)-l_{k}\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)} \tag{18}
\end{equation*}
$$

and we compare it with the value $\gamma \in] 0,1[$ initially fixed, where $\gamma$ is a parameter of the trust-region method that defines whether the iteration is successful. If $\rho_{k} \geq \gamma$, the objective function is well fitted by the quadratic model and we accept the $\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)$ as the new approximation $\left(\left(\mu_{k+1}, \lambda_{k+1}\right)=\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)\right), \Delta_{k+1}=\Delta_{k}$; otherwise, the trust region is too large and the fit between the function and the model is not satisfatory. In this case, we keep the previous approximation $\left(\left(\mu_{k+1}, \lambda_{k+1}\right)=\left(\mu_{k}, \lambda_{k}\right)\right)$ and reduce the trust region radius $\left(\Delta_{k+1}=\Delta_{k} / 2\right)$.

To update $P_{k+1}$ and $I_{k+1}$, regardless of whether the iteration is successful or not, the new trust region $I_{k+1}$ is the square centered on ( $\mu_{k+1}, \lambda_{k+1}$ ) with radius $\Delta_{k+1}$ so $\left(\mu_{k+1}, \lambda_{k+1}\right) \in I_{2,2, k+1}$. Additionally, we keep the points $P_{k} \cap I_{k+1}$ in the corresponding sub-set. Finally, we generate a random point on each sub-set $I_{i, j, k+1}$ that remains empty. $P_{k+1}$ is the set with the new $n_{k+1}$ points obtained by this process and a new quadratic model is obtained by the least squares method. We would like to emphasize that only the admissible portion of each sub-set is considered in $I_{k}$.

This procedure is repeated until the relative error between two consecutive iterations is less than $10^{-5}$ on both parameters. Algorithm 1 presents the sketch of this procedure which is a variation of the derivative free trust-region method presented in [15].

```
Algorithm 1: Derivative free trust-region method
    Initialization Choose \(\left(\mu_{0}, \lambda_{0}\right)\) such that \(l\left(\mu_{0}, \lambda_{0}\right)<\infty, \Delta_{0}>0\) and the
    constant \(\gamma \in] 0,1\left[\right.\). Obtain \(I_{0}\) and consequently \(P_{0}\). Define \(k=0\).
    repeat
        Construct the model for \(l_{k}(\mu, \lambda)\) by the least squares method applied to
        (15) with the points of \(P_{k}\).
        Obtain the critical point by solving problem (17).
        if \(\nabla^{2} l_{k}(\mu, \lambda)\) is positive definite and satisfies the constrain of the
            sub-problem (16) then
            \(\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)\) is the minimizer;
        else
            Obtain the minimizer \(\left(\mu_{k}^{*}, \lambda_{k}^{*}\right)\) over the boundary of the square
            centered in ( \(\mu_{k}, \lambda_{k}\) ) and radius \(\Delta_{k}\).
        Calculate \(\rho_{k}\) in (18).
        if \(\rho_{k} \geq \gamma\) then
            \(\Delta_{k+1}=\Delta_{k} ;\)
            \(\left(\mu_{k+1}, \lambda_{k+1}\right)=\left(\mu_{k}^{*}, \lambda_{k}^{*}\right) ;\)
        else
            \(\Delta_{k+1}=\Delta_{k} / 2 ;\)
            \(\left(\mu_{k+1}, \lambda_{k+1}\right)=\left(\mu_{k}, \lambda_{k}\right) ;\)
        Update \(I_{k+1}, P_{k+1}\) as described in this section.
        \(k=k+1 ;\)
    until \(\frac{\left|\mu_{k}-\mu_{k-1}\right|}{\mu_{k-1}}<10^{-5} \wedge \frac{\left|\lambda_{k}-\lambda_{k-1}\right|}{\lambda_{k-1}}<10^{-5}\);
```


## 4 Computational results

Let us consider the objective function defined by (12), which corresponds to the following setting: $\Omega=[-2,2]^{3}$ with $\partial \Omega=\Gamma_{1} \cup \Gamma_{2}$ where $\Gamma_{1}$ is the face of the cube contained in the plane $z=-2$; the mesh is a partition of $\Omega$ into 48 tetrahedrons; $\rho=1, w=2 \pi \times 10^{6}$; the functions $\mathbf{g}$ and $\mathbf{f}$ are defined respectively by $g_{i}=5.86 \times 10^{-3}$ and $f_{i}=0, i \in\{1,2,3\}$.

To illustrate the performance of the proposed method we used fabricated data obtained by simulating the direct problem. In particular, we considered $U_{o b s}$ as the solution of (10) with $E=4.66 \times 10^{6}$ and $\nu=0.45$ [2]. In this way, $(\mu, \lambda)=\left(1.6069 \times 10^{6}, 1.4462 \times 10^{7}\right)$ is the optimal solution of the inverse problem. For the optimization problem we choose the set $I=[0.9 \mu, 1.1 \mu] \times[0.9 \lambda, 1.1 \lambda]$ and we define $\Delta_{0}=10^{5}\left(\Delta_{0}<0.1 \max \{\mu, \lambda\}\right)$ and $\gamma=0.1$.

In the context of a real application, experimental data are affected by noise. Here we performed experiments with noise free data as well as noisy data in order to assess the sensitivity of our method to noise.

To check the robustness of the proposed method when considering noisy data we consider gaussian noise $R \sim N(0, \sigma)$ where $R$ is a random vector of dimension $3 N \times 1$ and $\sigma$ is the standard deviation. So instead $U_{\text {obs }}$, we consider as data $\bar{U}_{o b s}=\left(R+1_{3 N \times 1}\right) U_{\text {obs }}$, where $1_{3 N \times 1}$ is a $3 N \times 1$ vector with all components equal to one and the $i$-th component of the vector $\bar{U}_{o b s}$ is given by $(R(i)+1) U_{o b s}(i), i \in\{1, \ldots, 3 N\}$.

We consider variations of $\sigma$ in the set $\left\{0,10^{-8}, 10^{-7}, 10^{-6}, 10^{-5}\right\}$ and, for each value of $\sigma$, we consider simulations with 30 random initial points.

Figure 1 presents the relative error averages $\frac{\left|\mu_{k}-\mu_{k-1}\right|}{\mu_{k-1}}$ and $\frac{\left|\lambda_{k}-\lambda_{k-1}\right|}{\lambda_{k-1}}$ obtained from thirty simulations for each value of $\sigma$. As expected, increasing the number of iterations has the effect of decreasing the relative error between two consecutive iterations, allowing the approximations to converge close to the optimal solution and satisfy the stop condition. As we can be seen in Figure 1, the performance of Algorithm 1 is not significantly affected by the increase in noise.

Now since we have the thirty approximations for each value of $\sigma$, we will show in Figure 2 the evolution of the relative error average between the approximation given by the method and the solution $(\mu, \lambda)=\left(1.6069 \times 10^{6}, 1.4462 \times 10^{7}\right)$. As expected, both parameters show an increase in the average error with $\sigma$, which can lead to large perturbations in the optimal solution. Although the noise variation is not high, it generates higher relative errors comparing to the case of no noise ( $\sigma=0$ ). In Figure 2 we can see that for $\sigma \in\left\{10^{-6}, 10^{-5}\right\}$ the relative error in $\mu$ is very closed to 0.1 . This situation happens because the level of noise considered change the behaviour of the objective function where the minimum in attained in the boundary of $I$. Note that in the boundary of the domain $I$ the approximations achieve relative error of the order 0.1.


Fig. 1: Relative error average obtained from thirty simulations for the parameters $\mu$ (left) and $\lambda$ (right) considering different levels of noise.


Fig. 2: Relative error average obtained from thirty simulations for each value of $\sigma$.

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