

Electromagnetic Inverse Problems Involving Distributions of Dielectric Mechanisms and Parameters

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Abstract

We consider electromagnetic interrogation problems for complex materials involving distributions of polarization mechanisms and also distributions for the parameters in these mechanisms. A theoretical and computational framework for such problems is given. Computational results for specific problems with multiple Debye mechanisms are given in the case of discrete, uniform, log-normal, and log-Bi-Gaussian distributions.

Keywords: Electromagnetic interrogation with pulsed antenna source microwaves and inverse problems, complex dielectric materials, distributions of relaxation parameters and mechanisms.

1 Introduction

For at least the past century [50, 51], scientific investigators have sought to understand what happens to electromagnetic fields (and how to mathematically model the associated phenomena) when they are introduced into complex materials such as biotissue and more general dielectrics, conductors and magnetics. More specifically, a fundamental question is how to model dispersion and dissipation of the fields in these complex materials. This has most often led to the use of Maxwell's equations in a non-vacuum environment which entails constitutive relationships for polarization (in dielectrics), magnetization (in magnetic materials) and conductivity. We focus here on modeling polarization in dielectric materials for which we develop a new modeling framework. Even though we treat only polarization as our dispersive mechanism in our formulation (adopting Ohm's law for conductivity and considering non-magnetic materials), the approach is sufficiently general so as to be readily extended to treat magnetization and conductivity in materials (each in some type of convolution representation involving susceptibility kernels, e.g., see [2, 3]). We develop a framework that allows not only uncertainty (through distributions of parameters representing molecular variability) at the molecular level, but also allows for the presence of multiple polarization mechanisms in the material.

We first explain a conceptual framework in the context of the 3-D Maxwell equations in a dielectric material. In particular, we use Maxwell's equations which govern the electric field \mathbf{E} and the magnetic field \mathbf{H} in a domain $\mathcal{D} = \Omega_0 \cup \Omega$ with charge density ρ in the material Ω while the ambient Ω_0 is treated as a

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vacuum. Thus we first consider the system

$$\left\{ \begin{array}{ll} \text{(i)} & \frac{\partial \mathbf{D}}{\partial t} + \mathbf{J} - \nabla \times \mathbf{H} = 0, \text{ in } (0, T) \times \mathcal{D}, \\ \text{(ii)} & \frac{\partial \mathbf{B}}{\partial t} + \nabla \times \mathbf{E} = 0, \text{ in } (0, T) \times \mathcal{D}, \\ \text{(iii)} & \nabla \cdot \mathbf{D} = \rho, \text{ in } (0, T) \times \mathcal{D}, \\ \text{(iv)} & \nabla \cdot \mathbf{B} = 0, \text{ in } (0, T) \times \mathcal{D}, \\ \text{(v)} & \mathbf{E} \times \mathbf{n} = 0, \text{ on } (0, T) \times \partial \mathcal{D}, \\ \text{(vi)} & \mathbf{E}(0, \mathbf{x}) = 0, \quad \mathbf{H}(0, \mathbf{x}) = 0, \text{ in } \mathcal{D}. \end{array} \right. \quad (1)$$

As usual, the current \mathbf{J} is composed of the source current \mathbf{J}_s and the conductive current \mathbf{J}_c . Within the domain we have constitutive relations that relate the flux densities \mathbf{D}, \mathbf{B} and the conductive current \mathbf{J}_c to the electric and magnetic fields. We have

$$\left\{ \begin{array}{ll} \text{(i)} & \mathbf{D} = \epsilon_0 \mathbf{E} + \mathbf{P}_T I_\Omega, \\ \text{(ii)} & \mathbf{B} = \mu_0 \mathbf{H}, \\ \text{(iii)} & \mathbf{J}_c = \sigma \mathbf{E} I_\Omega. \end{array} \right. \quad (2)$$

In (2), I_Ω denotes the indicator function on the dielectric medium Ω . Thus $\mathbf{J}_c = 0$ in the ambient or air. The total electric polarization \mathbf{P}_T is given by

$$\mathbf{P}_T = \mathbf{P}_I + \mathbf{P} = \epsilon_0 \chi \mathbf{E} + \mathbf{P},$$

where \mathbf{P}_I is the instantaneous polarization due to the interface between Ω_0 and Ω and \mathbf{P} is the material or dielectric polarization. Hence the constitutive law (2, i) in Ω becomes

$$\mathbf{D} = \epsilon_0 \epsilon_r \mathbf{E} + \mathbf{P},$$

where $\epsilon_r = (1 + \chi)$ is the relative permittivity of the dielectric medium.

Our main focus in this presentation is the dielectric polarization \mathbf{P} which we assume has the general convolution form

$$\mathbf{P}(t, \mathbf{x}) = g \star \mathbf{E}(t, \mathbf{x}) = \int_0^t g(t-s, \mathbf{x}) \mathbf{E}(s, \mathbf{x}) ds, \quad (3)$$

where g is the general dielectric response function (DRF). In every practical example (Debye, Lorentz, etc.) DRFs are parameter dependent as well as time (and possibly space) dependent; we represent this as $g = g(t, \mathbf{x}; \nu)$, where typically $\nu = (\epsilon_\infty, \epsilon_s, \tau)$ contains parameters such as the high frequency limit dielectric permittivity ϵ_∞ , the static permittivity ϵ_s , and relaxation time τ . Examples of often-used DRFs are the Debye [11, 23, 29] in a material region Ω defined in the time domain by

$$g(t, \mathbf{x}) = \epsilon_0 (\epsilon_s - \epsilon_\infty) / \tau \quad e^{-t/\tau},$$

the Lorentz [11, 23, 40] given by

$$g(t, \mathbf{x}) = \epsilon_0 \omega_p^2 / \nu_0 e^{-t/2\tau} \sin(\nu_0 t),$$

and the Cole-Cole [23, 27, 32, 38, 46] defined by

$$g(t, \mathbf{x}) = \mathcal{L}^{-1} \left\{ \frac{\epsilon_0 (\epsilon_s - \epsilon_\infty)}{1 + (s\tau)^\alpha} \right\} = \frac{1}{2\pi i} \int_{\zeta - i\infty}^{\zeta + i\infty} \frac{\epsilon_0 (\epsilon_s - \epsilon_\infty)}{1 + (s\tau)^\alpha} e^{st} ds,$$

where \mathcal{L} is the Laplace transform.

These DRFs also play a fundamental role in convolution representations such as (3) for nonlinear polarization laws [8, 19, 20, 24, 40] of Kerr type and Raman scattering [40]. While the ideas we describe here on distributions of relaxation times and mechanisms can readily be used to treat such nonlinear polarization

laws, we shall in this paper concentrate on linear constitutive laws involving multiple relaxation parameters and multiple dispersive mechanisms in materials in the presence of multiple interrogating frequencies.

The macroscopic polarization model (3) can be derived from microscopic dipole, electron cloud, etc., formulations by passing to a limit over the molecular population [30]. However, such derivations tacitly assume that one has similar individual (molecular, dipole, etc.) parameters; that is, all dipoles, molecules, “electron clouds”, etc., have the same relaxation parameters, plasma frequencies, etc. Historically, such models based on molecular level homogeneity throughout the material have often not performed well when trying to compare models with experimental data. Indeed, in 1907 Von Schweidler [23, 50] observed the need to assume *multiple relaxation times* when considering experimental data and in 1913 Wagner [23, 51] proposed continuous distributions of relaxation times. This idea was subsequently visited by Williams and Ferry [55] in comparing the behavior of viscoelastic polymers and dielectric materials. In the past half century intensive experimental efforts [31, 32, 33, 34, 35, 38, 45] have been pursued in describing data for complex materials with distributions of dielectric parameters (especially relaxation times in multiple Debye [31] or multiple Lorentz [40] mechanisms) in the frequency domain. A significant amount of this work is reviewed in the survey paper by Foster and Schwan [31]. There are now *incontrovertible* experimentally based arguments for distributions of relaxation parameters in mechanisms for heterogeneous materials. Moreover, there is compelling evidence of the presence of *multiple mechanisms* in complex materials such as tissue and modern polymeric composites. These multiple mechanisms may involve interfacial polarization, dipolar orientation, ionic diffusion (e.g., see p. 40, 49, 57 of [31]) and may often require a selection of several types of distributional representations from examples such as the fractional power laws of the Cole-Cole [25, 26, 31, 35, 46], the log normal, the uniform, as well as the Debye and Lorentz (although the fractional power law of Cole-Cole is more the rule rather than the exception – p. 39, [31]). These multiple mechanisms are likely present in some weighted combination (e.g., see [36] and p. 369, [40]) and often are manifested in a frequency-dependent manner. It is therefore advantageous to consider interrogation or inverse problems with multiple frequencies (e.g., ranging from RF (10^6) to GHz (10^{10})) or broadband excitation signals.

To allow for a distribution F of parameters ν over some admissible set \mathcal{N} , we generalize the polarization law (3) to

$$\mathbf{P}(t, \mathbf{x}; F) = h \star \mathbf{E}(t, \mathbf{x}) = \int_0^t \int_{\mathcal{N}} g(t-s, \mathbf{x}; \nu) \mathbf{E}(s, \mathbf{x}) dF(\nu) ds. \quad (4)$$

We expect to chose F from (or from a subspace of) the space $\mathcal{F} = \mathfrak{P}(\mathcal{N})$ of all probability measures F on \mathcal{N} .

We further generalize (4) to allow for dielectric materials with multiple mechanisms or multiple DRFs (i.e., heterogeneous molecular structures) by considering a family \mathcal{G} of possible DRFs and distributions M over this family. This leads to the polarization constitutive relationship

$$\mathbf{P}(t, \mathbf{x}; M, F) = \int_0^t \int_{\mathcal{G}} \int_{\mathcal{N}} g(t-s, \mathbf{x}; \nu) \mathbf{E}(s, \mathbf{x}) dF(\nu) dM(g) ds = \int_0^t K(t-s, \mathbf{x}; M, F) \mathbf{E}(s, \mathbf{x}) ds, \quad (5)$$

where for $F \in \mathcal{F} = \mathfrak{P}(\mathcal{N})$ and $M \in \mathcal{M} = \mathfrak{P}(\mathcal{G})$, K is defined by

$$K(t-s, \mathbf{x}; M, F) = \int_{\mathcal{G}} \int_{\mathcal{N}} g(t-s, \mathbf{x}; \nu) dF(\nu) dM(g). \quad (6)$$

When we use (5) and (6) in the Maxwell system (1)-(2), we are led to a system of partial differential equations where lower order terms (in time) depend on probability measures. These measures are now the “parameters” that characterize the material dielectric properties which one must estimate or identify in interrogation problems.

With the recently growing interest in incorporating uncertainty into models and systems, the need to employ dynamics with probabilistic structures has received increased emphasis. In particular, systems with probability measures embedded in the dynamics (problems involving *aggregate dynamics* as discussed in [8]) have become important in applications in biology [9, 10, 8], electromagnetics [12] and hysteretic [15, 16, 41, 42] and polymeric [17, 18, 21] materials. These systems (in the case of first order ordinary differential equations) have the form

$$\dot{x}(t) = f(t, x(t), F),$$

where F is a probability distribution or measure. In fact such systems are not new and arise in *relaxed* or *chattering* control problems [43, 44, 47, 52, 53, 54] wherein the controls are probability measures. Indeed, such systems date back to the seminal work of L.C. Young on *generalized curves* in the calculus of variations [56, 57].

In the next section we formulate an inverse problem for the Maxwell system with the general polarization law (5) and discuss theoretical and computational aspects of this problem. We then report on our initial computational efforts on a 1-dimensional version of the problem of finding the underlying polarization in a slab of material using signals reflected from a front air/material interface and a metal backing. We present numerical results for inverse problems involving Maxwell's equations with absorbing left boundary condition, a supraconducting right boundary condition and a general polarization term which includes uncertainty in the dielectric parameters. We attempt to determine an unknown distribution of parameters which describes the dielectric properties of the material. We explore both discrete and continuous distributions; for the continuous case appropriate parameterizations and discretizations are used.

2 A General Inverse Problem

We consider the Maxwell system (1)-(2) with polarization $\mathbf{P} = \mathbf{P}(t, \mathbf{x}; M, F)$ given by (5). Let

$$\mathbf{z}(t, \mathbf{x}; M, F) = \begin{pmatrix} \mathbf{E}(t, \mathbf{x}; M, F) \\ \mathbf{H}(t, \mathbf{x}; M, F) \end{pmatrix}$$

with $(t, \mathbf{x}) \rightarrow \mathbf{z}(t, \mathbf{x}; M, F)$ mapping from $(0, T) \times \Omega$ to R^6 . We assume we are given data $\bar{\mathbf{d}} = \{\mathbf{d}_i\}_{i=1}^n$ corresponding to observations of $\mathcal{C}_A \mathbf{z}(t_i, \cdot; M, F)$. Here \mathcal{C}_A denotes evaluation of one or more components of \mathbf{E} or \mathbf{H} at an antenna $\{\mathbf{x}_A\}$. We use this data to estimate the distributions M and F in an ordinary least squares (OLS) formulation, seeking to minimize

$$J(M, F) = \sum_{i=1}^n |\mathcal{C}_A \mathbf{z}(t_i, \cdot; M, F) - \mathbf{d}_i|^2 \quad (7)$$

over $(M, F) \in \mathcal{M} \times \mathcal{F} = \mathfrak{P}(\mathcal{G}) \times \mathfrak{P}(\mathcal{N})$. We thus seek to find (\hat{M}, \hat{F}) such that

$$(\hat{M}, \hat{F}) = \arg \min \{J(M, F) : M \in \mathcal{M}, F \in \mathcal{F}\}.$$

We note that while for simplicity we use an OLS formulation here, most of the results discussed below could readily be developed in the context of other standard estimation formulations such as maximum likelihood estimators (MLE), weighted least squares (WLS), or generalized least squares (GLS).

For theoretical and computational purposes, one needs a topology on \mathcal{M} and \mathcal{F} and for this we choose the Prohorov metric ρ_* of *weak** convergence in \mathcal{M} and \mathcal{F} when they are considered as subsets of the topological duals $C_B^*(\mathcal{G})$ and $C_B^*(\mathcal{N})$ of the spaces $C_B(\mathcal{G})$ and $C_B(\mathcal{N})$ of bounded continuous functions on \mathcal{G} and \mathcal{N} , respectively [5, 8]. That is, $F_k \rightarrow F$ in the ρ_* metric if and only if

$$\int_{\mathcal{N}} \phi(\nu) dF_k(\nu) \rightarrow \int_{\mathcal{N}} \phi(\nu) dF(\nu)$$

for all $\phi \in C_B(\mathcal{N})$, i.e., all bounded continuous ϕ on \mathcal{N} ; similarly for $M_k \rightarrow M$. It is known [5, 8] that if \mathcal{G} and \mathcal{N} are complete metric spaces, then $\mathcal{M} = \mathfrak{P}(\mathcal{G})$ and $\mathcal{F} = \mathfrak{P}(\mathcal{N})$ taken with the Prohorov metric ρ_* are complete metric spaces. Moreover, if \mathcal{G} and/or \mathcal{N} are compact, then so are \mathcal{M} and/or \mathcal{F} . Using these properties and arguments similar to those in [5, 14], the following *problem stability* results can be proven.

Theorem 1: Suppose $(M, F) \rightarrow \mathcal{C}_A \mathbf{z}(t, \cdot; M, F)$ is continuous on $\mathcal{M} \times \mathcal{F}$ and suppose that \mathcal{G} and \mathcal{N} are compact. Then solutions (\hat{M}, \hat{F}) of minimizing (7) exist (generally, non-uniquely) and are continuous in the data $\bar{\mathbf{d}}$ in the following sense. Suppose $(M^*(\bar{\mathbf{d}}), F^*(\bar{\mathbf{d}}))$ and $(M^*(\bar{\mathbf{d}}^k), F^*(\bar{\mathbf{d}}^k))$ are the solution sets of minimizing (7) for data $\bar{\mathbf{d}}$ and $\bar{\mathbf{d}}^k$, respectively, where $\bar{\mathbf{d}}^k \rightarrow \bar{\mathbf{d}}$ as $k \rightarrow \infty$. Then

$$\text{dist} [(M^*(\bar{\mathbf{d}}^k), F^*(\bar{\mathbf{d}}^k)), (M^*(\bar{\mathbf{d}}), F^*(\bar{\mathbf{d}}))] \rightarrow 0$$

as $k \rightarrow \infty$, where $dist[\cdot, \cdot]$ is the Hausdorff distance in the metric space $\mathcal{M} \times \mathcal{F}$.

The problem of minimizing (7) over $\mathcal{M} \times \mathcal{F}$ is in general an infinite dimensional optimization problem that poses formidable computational challenges. But we note that set \mathcal{G} of possible DRFs can in almost all cases be taken as a fixed finite set, i.e., $\mathcal{G} = \{g_1, g_2, \dots, g_K\}$ which means that \mathcal{M} would be given by

$$\mathcal{M} = \left\{ M(g) \in \mathfrak{P}(\mathcal{G}) \mid dM(g) = \sum_{j=1}^K a_j d\Delta_{g_j}(g), a_j \geq 0, \sum_j a_j = 1 \right\},$$

where Δ_{g_j} is the Dirac measure with single atom at g_j , i.e., $d\Delta_{g_j}(g) = \delta_{g_j}(g)dg$. Thus, this part of the optimization reduces to one over a closed bounded convex set in Euclidean space. The minimization over $\mathcal{F} = \mathfrak{P}(\mathcal{N})$ is more interesting since in general one expects \mathcal{N} to involve a continuum of (vector) parameters. We illustrate the possibilities with $\mathcal{N} = \mathcal{T} = \{\tau \mid \tau \in [\tau_a, \tau_b]\}$ for τ the relaxation parameter in, for example, a Debye or Lorentz mechanism. There are a number of ways this problem could be approached:

1. Assume that F is discrete, having the fixed form $dF(\tau) = \sum \alpha_j \delta_{\tau_j}(\tau) d\tau$ and seek to find (α_j, τ_j) minimizing (7) where $\alpha_j \geq 0, \sum \alpha_j = 1, \tau_j \in [\tau_a, \tau_b]$;
2. Assume that F is (absolutely) continuous and given in parametric form $dF(\tau) = f(\tau, \mu, \sigma) d\tau$ where f is known (e. g., normal, log-normal, uniform, etc.) and seek to find (μ, σ) ;
3. Assume F does not have a specific parametric form and seek to find the general form for F through the optimization of (7).

In the first two cases above, one effectively reduces the inverse problem to a computationally tractable (one hopes!—we explore these ideas computationally in subsequent discussions below) optimization problem that is finite dimensional. The third case remains infinite dimensional in nature and one must develop approximation ideas that lead to implementable computational algorithms. In [5], the authors developed approximation ideas based on density results for measures arising in probability theory. We only outline those here, referring the reader to [5] for more details and proofs.

To develop approximation ideas for the nonparametric case, we first consider a family of partition points $\mathcal{T}^N = \{q_j^N\}_{j=1}^N$, $N = 1, 2, \dots$, such that $\cup_{N=1}^\infty \mathcal{T}^N$ is dense in \mathcal{T} . Then define

$$\mathcal{F}^N = \mathfrak{P}^N(\mathcal{T}) = \{F^N \in \mathfrak{P}(\mathcal{T}) \mid dF^N(\tau) = \sum_{j=1}^N p_j^N \delta_{q_j^N}(\tau) d\tau, q_j^N \in \mathcal{T}^N, p_j^N \text{ rational}, \sum p_j^N = 1\}.$$

It can be argued that $\cup_{N=1}^\infty \mathcal{F}^N$ is dense in \mathcal{F} in the Prohorov metric ρ_* . Moreover, if \mathcal{T} is compact, one can prove a *method stability* theorem (see [5, 8]) similar to Theorem 1 above. Specifically, let $F_N^*(\bar{\mathbf{d}})$ be the set of solutions obtained in minimizing (7) with \mathcal{F} replaced by \mathcal{F}^N . Then the method stability theorem guarantees

$$dist[(M^*(\bar{\mathbf{d}}^k), F_N^*(\bar{\mathbf{d}}^k)), (M^*(\bar{\mathbf{d}}), F^*(\bar{\mathbf{d}}))] \rightarrow 0$$

as $N, k \rightarrow \infty$.

More generally, one may wish to consider only classes of distributions \mathcal{F} that arise from densities functions, i.e., absolutely continuous distributions in

$$\mathcal{F}_{AC} = \{F \in \mathcal{F} \mid F' = f, f \in \mathcal{F}_{weak}\}$$

where \mathcal{F}_{weak} is a given weakly compact subset of $L^2(\mathcal{T})$. For example, \mathcal{F}_{weak} could be any given closed, convex, bounded subset of $L^2(\mathcal{T})$. It is proven in [21] that sets \mathcal{F}_{AC} defined in this way are compact (in the ρ_* metric) subsets of \mathcal{F} and the corresponding existence and stability results of Theorem 1 hold for problems using \mathcal{F}_{AC} in place of \mathcal{F} in the optimization for (7). But once again these are infinite dimensional in nature and approximations are needed. We note, that although the computational framework described above utilizing Dirac measures is also valid here, it is often desirable to develop “smoother” approximations to elements of \mathcal{F}_{AC} . In particular, suppose that $f \in \mathcal{F}_{weak}$ and $F \in \mathcal{F} = \mathfrak{P}(\mathcal{T})$ with $F = \int f$. Since \mathcal{F}_{weak}

is a subset of $L^2(\mathcal{T})$, we can use any number of types of splines to formulate an approximation to f , i.e., let $\{S_j^N\}_{j=1}^N$, $N = 1, 2, \dots$ be a standard spline family [14, 48, 49] (e.g., piecewise-linear, cubic, Hermite cubic, etc.) and

$$f^N(\tau) = \sum_{j=1}^N w_j^N S_j^N(\tau),$$

where the rational numbers w_j^N are chosen so that $f^N \rightarrow f$ in $L^2(\mathcal{T})$. This implies that

$$\int_{\mathcal{T}} \phi f^N d\tau \rightarrow \int_{\mathcal{T}} \phi f d\tau$$

for all $\phi \in L^2(\mathcal{T})$, and hence for all $\phi \in C(\mathcal{T})$, which yields

$$\rho(F^N, F) \rightarrow 0,$$

where $F^N = \int f^N$. Thus defining

$$\mathcal{F}_{SPL}^N = \left\{ h \in L^2(\mathcal{T}) \mid h(\tau) = \sum_{j=1}^N w_j^N S_j^N(\tau) \right\},$$

we can conclude that the set

$$\tilde{\mathcal{F}}_{AC} = \{F \in \mathcal{F} \mid F = \int f, f \in \cup_{N=1}^{\infty} \mathcal{F}_{SPL}^N\}$$

is dense in \mathcal{F}_{AC} in the ρ_* -metric. Hence spline families provide an alternative way to approximate elements of \mathcal{F}_{AC} in our computational work. We note that in this case the approximating elements are not probability distributions themselves. Once again one can state and prove a *method stability* theorem [14] using these spline approximations. That is:

Theorem 2: Suppose \mathcal{T} is compact and the solutions $\mathcal{C}_{AZ}(t, \cdot; M, F)$ are continuous. Let \mathcal{F}_{weak} , \mathcal{F}_{AC} , and \mathcal{F}_{SPL}^N be as defined above with $(M^*(\bar{\mathbf{d}}), F_N^*(\bar{\mathbf{d}}))$ the solution sets for minimizing (7) with \mathcal{F} replaced by \mathcal{F}_{SPL}^N . Then for $\bar{\mathbf{d}}^k \rightarrow \bar{\mathbf{d}}$ we have

$$dist[(M^*(\bar{\mathbf{d}}^k), F_N^*(\bar{\mathbf{d}}^k)), (M^*(\bar{\mathbf{d}}), F^*(\bar{\mathbf{d}}))] \rightarrow 0$$

as $N, k \rightarrow \infty$. Thus solutions depend continuously on data and the approximate problems are method stable.

To illustrate computational aspects of the ideas presented in this section, we turn next to a 1-D version of the inverse problem and present results in the next several sections on use of discrete and continuous distributions in the polarization laws.

3 The 1-D Problem Formulation

For our initial numerical efforts, we turned to the 1-D example as explained in detail in [11]. Under the assumptions detailed there, one obtains a domain as depicted in Figure 1.

Restricting to one dimension, and using $D = \epsilon E + P$, we can write Maxwell's equations in second order form as

$$\mu_0 \epsilon \ddot{E} + \mu_0 I_{\Omega} \ddot{P} + \mu_0 \sigma I_{\Omega} \dot{E} - E'' = -\mu_0 \dot{J}_s \text{ in } \Omega \cup \Omega_0, \quad (8)$$

where E is the \hat{k} or z component of the electric field, P is the media's macroscopic electric polarization, J_s is the interrogating signal, Ω is the domain of the material under investigation, Ω_0 is the ambient domain (considered a vacuum), and $\sigma = \sigma(z)$ is the conductivity of the material. Note we are considering a non-magnetic material containing no charge distribution ($\rho = 0$). Also, let $\epsilon = \epsilon_0(1 + I_{\Omega}(\epsilon_r - 1))$ where $\epsilon_r = \epsilon_r(z)$

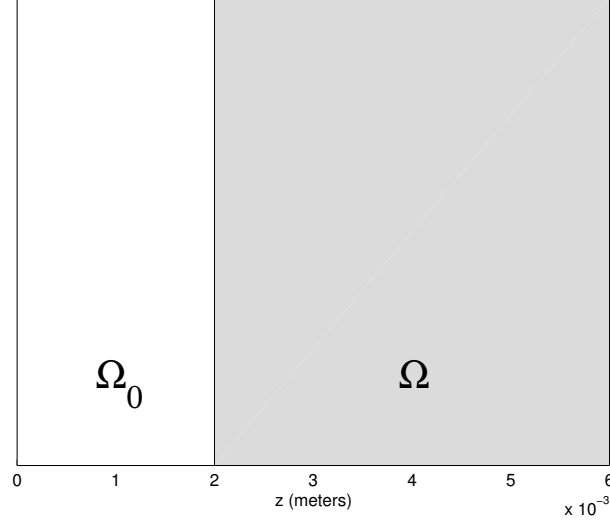


Figure 1: The 1-D domain.

is the dielectric permittivity. The values μ_0 and ϵ_0 are the magnetic permeability and the electric permittivity of free space, respectively. See [11] for more details. The boundary conditions that we are assuming are absorbing at $z = 0$ (to provide a finite computational window) and supraconducting at $z = 1$ (representing a metal backing):

$$\begin{aligned} \left[\dot{E} - cE' \right]_{z=0} &= 0 \\ E(t, 1) &= 0. \end{aligned} \tag{9}$$

We assume homogeneous initial conditions

$$E(0, z) = 0 \tag{10}$$

$$\dot{E}(0, z) = 0. \tag{11}$$

For implementation, we scale time by $\tilde{t} = ct$ and the polarization by $\tilde{P} = P/\epsilon_0$ for convenience. Also note that we have employed the “method of mappings” to obtain a computational domain of $z = [0, 1]$. The actual dimensions of the domains considered in this report depend on the interrogating frequency ω . In particular, for $\omega = 2\pi \times 10^{11}$ we consider a material slab of thickness $.004m$ preceded by a vacuum of depth $.002m$. We scale these dimensions as the wavelength scales when we change interrogating frequencies.

Converting the system (8) to weak form (and dropping the \sim notation) we obtain

$$\langle \epsilon_r \ddot{E}(t, \cdot), \phi \rangle + \langle I_\Omega \ddot{P}(t, \cdot), \phi \rangle + \langle \eta_0 \sigma I_\Omega \dot{E}(t, \cdot), \phi \rangle - \langle E''(t, \cdot), \phi \rangle = -\langle \eta_0 \dot{J}_s(t, \cdot), \phi \rangle,$$

where $\eta_0 = \sqrt{\mu_0/\epsilon_0}$ and $\phi \in V = H_R^1(0, 1) = \{\phi \in H^1(0, 1) : \phi(1) = 0\}$. Finally, we integrate by parts, and apply the boundary conditions (9) to get

$$\langle \epsilon_r \ddot{E}(t, \cdot), \phi \rangle + \langle I_\Omega \ddot{P}(t, \cdot), \phi \rangle + \langle \eta_0 \sigma I_\Omega \dot{E}(t, \cdot), \phi \rangle + \langle E'(t, \cdot), \phi' \rangle + \dot{E}(t, 0)\phi(0) = -\langle \eta_0 \dot{J}_s(t, \cdot), \phi \rangle. \tag{12}$$

To describe the behavior of the electric polarization P , we may employ a general polarization kernel, or dielectric response function, g as follows:

$$P(t, z) = \int_0^t g(t - s, z; \tau) E(s, z) ds \tag{13}$$

where, for instance using a Debye polarization model,

$$g(t; \tau) = \epsilon_0(\epsilon_s - \epsilon_\infty)/\tau e^{-t/\tau}.$$

However, this presupposes that the material may be sufficiently defined by a single relaxation parameter τ , which is generally not the case. In order to account for uncertainty in the polarization mechanisms, we allow for a distribution of relaxation parameters. Thus, we define our polarization model in terms of a distribution dependent dielectric response function h

$$P(t, z) = \int_0^t h(t-s, z; F) E(s, z) ds, \quad (14)$$

where h is determined by a family of polarization laws each described by a different parameter τ , and therefore is given by

$$h(t, z; F) = \int_{\mathcal{T}} g(t, z; \tau) dF(\tau),$$

where $\mathcal{T} = [\tau_1, \tau_2] \subset (0, \infty)$. In particular, if the distribution F were discrete, consisting of a single relaxation parameter, then we would again have (13).

The macroscopic electric polarization becomes

$$P(t, z) = \int_0^t \left[\int_{\mathcal{T}} g(t-s, z; \tau) dF(\tau) \right] E(s, z) ds,$$

or, interchanging integrals, we have

$$P(t, z) = \int_{\mathcal{T}} \mathcal{P}(t, z; \tau) dF(\tau),$$

where

$$\mathcal{P}(t, z; \tau) = \int_0^t g(t-s, z; \tau) E(s, z) ds$$

is the polarization due to the relaxation parameter τ . Thus assuming we have a computational method to compute (13), i.e., \mathcal{P} , we use this as a basis for approximating P , either directly in the discrete case, or using a quadrature rule in the continuous case.

The theory developed in Section 2 for general inverse problems can be directly applied to this 1-D problem once one has established continuity of the solution E with respect to the measures F in the space $\mathcal{F} = \mathfrak{P}(\mathcal{T})$ taken with the Prohorov metric. But these desired continuity results are given in [12].

3.1 Discrete Distribution

Consider the discrete distribution given by

$$dF(\tau) = \sum_{i=0}^{\ell} \frac{\delta(\tau_i)}{\ell+1} d\tau,$$

where δ represents the Dirac distribution, and $\tau_i = \tau_a + i\tau_h$ with $\tau_h = \frac{\tau_b - \tau_a}{\ell}$. Then we have

$$P(t, z) = \sum_{i=0}^{\ell} \alpha_i \mathcal{P}(t, z; \tau_i),$$

where $\alpha_i = \frac{1}{\ell+1}$. Note that for this example, the nodes (τ_i) are linearly spaced and the weights or masses (α_i) are uniform. More generally, one can treat a discrete distribution without linearly spaced nodes and/or uniform masses.

3.2 Uniform Distribution

The simplest continuous distribution is a uniform distribution. Consider

$$dF(\tau) = \frac{1}{\tau_b - \tau_a} d\tau,$$

for $\tau_a \leq \tau \leq \tau_b$, and zero otherwise. Since the distribution function is constant, we may pull it outside the integral giving

$$P(t, z) = \frac{1}{\tau_b - \tau_a} \int_{\tau_a}^{\tau_b} \mathcal{P}(t, z; \tau) d\tau. \quad (15)$$

While this reduces down to essentially integrating τ out of the polarization convolution term, for most polarization models this does not have an analytical solution. Therefore we must resort to numerical quadrature. Recall the Composite Simpson's rule approximation to

$$I_{[a,b]}(L) := \int_a^b L(x) dx$$

is

$$S_{[a,b]}^\ell(L) := \sum_{i=0}^{\ell} c_i^\ell L(x_i) h,$$

where ℓ is even, $x_i = a + ih$, $h = \frac{b-a}{\ell}$, and the weights are given by

$$c_i^\ell = \begin{cases} \frac{1}{3} & \text{if } i = 0 \text{ or } i = \ell \\ \frac{4}{3} & \text{else if } i \text{ odd} \\ \frac{2}{3} & \text{else if } i \text{ even} \end{cases}.$$

Thus, our Composite Simpson's approximation to (15) can be written

$$\begin{aligned} P(t, z) &= \frac{1}{\tau_b - \tau_a} \sum_{i=0}^{\ell} c_i^\ell \mathcal{P}(t, z; \tau_i) \tau_h \\ &= \sum_{i=0}^{\ell} \alpha_i \mathcal{P}(t, z; \tau_i) \end{aligned}$$

with

$$\alpha_i = \frac{c_i^\ell}{\ell},$$

since the numerator in $\tau_h = \frac{\tau_b - \tau_a}{\ell}$ cancels the constant in front of the integral. Note here that while our nodes are still linearly spaced, now our weights are non-uniform. We could have used a uniform discrete distribution to approximate the continuous one, which would have resulted in uniform weights, but the corresponding quadrature rule for $g(t, z; \tau)$ would have only been $O(\tau_h)$. Composite Simpson's rule provides better accuracy for the same number of function evaluations.

3.3 Log-Normal

A more realistic model for the distribution of Debye relaxation times is given by the log-normal distribution (see [23]). Therefore we consider the probability distribution defined by

$$dF(\tau) = \frac{1}{\sqrt{2\pi\sigma^2}} \frac{1}{\ln 10} \frac{1}{\tau} \exp\left(-\frac{(\log \tau - \mu)^2}{2\sigma^2}\right) d\tau,$$

which means that $\log \tau$ is normally distributed with mean μ and variance σ^2 . The $\ln 10$ term appears because we are using base 10 logarithms. Thus

$$P(t, z) = \int_0^\infty \mathcal{P}(t, z; \tau) dF(\tau) \approx \int_{\tau_a}^{\tau_b} \mathcal{P}(t, z; \tau) dF(\tau), \quad (16)$$

for τ_a sufficiently close to zero and τ_b sufficiently large. Since $\log \tau$ is normally distributed, we choose our nodes based on a uniform discretization of $\log \tau$, denoted $\{\xi_i\}_{i=0}^\ell$ from $\log \tau_a := \mu - 6\sigma$ to $\log \tau_b := \mu + 6\sigma$, where τ_a and τ_b are thusly defined by their logarithms. In this way we reduce our support to a compact set which should be large enough to effectively approximate our integral. In order to use Composite Simpson's rule we prefer to have a uniform discretization of nodes. Therefore we first change variables so that the integral is in terms of $\xi = \log \tau$. First, let

$$\hat{f}(\xi) = \frac{1}{\sqrt{2\pi}\sigma^2} \exp\left(-\frac{(\log \tau - \mu)^2}{2\sigma^2}\right)$$

and note that

$$d\xi = d\left(\frac{\ln \tau}{\ln 10}\right) = \frac{1}{\ln 10} \frac{1}{\tau} d\tau.$$

Then we have that

$$P(t, z) \approx \int_{\tau_a}^{\tau_b} \mathcal{P}(t, z; \tau) dF(\tau) = \int_{\mu-6\sigma}^{\mu+6\sigma} \mathcal{P}(t, z; 10^\xi) \hat{f}(\xi) d\xi$$

where implicitly $dF = f(\tau)d\tau$ and $d\hat{F} = \hat{f}(\xi)d\xi$. Applying Composite Simpson's rule gives finally

$$P(t, z) \approx \sum_{i=0}^{\ell} \alpha_i \mathcal{P}(t, z; \tau_i)$$

where we define $\tau_i = 10^{\xi_i}$ and

$$\alpha_i = c_i^\ell \hat{f}(\xi_i) \xi_h = \frac{c_i^\ell 12\sigma}{\ell} \frac{1}{\sqrt{2\pi}\sigma^2} \exp\left(-\frac{(\log \tau_i - \mu)^2}{2\sigma^2}\right).$$

For this case our nodes are non-linearly spaced in τ and our weights are non-uniform.

The method used in the bi-gaussian case is derived in a similar fashion. Consider that the polarization is driven by two distinct mechanisms, one with a dielectric distribution determined by mean μ_1 and standard deviation σ_1 , and the other determined by mean μ_2 and standard deviation σ_2 . Then we define $F_1(\tau; \mu_1, \sigma_1)$ and $F_2(\tau; \mu_2, \sigma_2)$ as log normal distributions as above, and let our macroscopic electric polarization be a function of some combination of these distributions (e.g., determined by the relative volume percentage of each of two substances in a material). Thus if $dF(\tau) = \beta_1 dF_1(\tau) + \beta_2 dF_2(\tau)$ then we again have the representation (16). For the discretization, we prefer to apply Composite Simpson's to each distribution separately and then combine at the end. In other words

$$P(t, z) \approx \beta_1 \sum_{i=0}^{\ell} \alpha_i^1 \mathcal{P}(t, z; \tau_i^1) + \beta_2 \sum_{i=0}^{\ell} \alpha_i^2 \mathcal{P}(t, z; \tau_i^2),$$

where $\{\alpha_i^1\}$ and $\{\tau_i^1\}$ are determined by μ_1 and σ_1 , and $\{\alpha_i^2\}$ and $\{\tau_i^2\}$ are determined by μ_2 and σ_2 , as explained above.

4 Inverse Problem Formulation

Our goal is to estimate the probability distribution function (PDF) $F \in \mathcal{P}(\mathcal{T})$ of relaxation parameters in a given model of the polarization, where $\mathcal{P}(\mathcal{T})$ is the set of all PDFs on the admissible region $\mathcal{T} = [\tau_1, \tau_2] \subset (0, \infty)$, by using reflections of electromagnetic interrogating signals off a metallic backing of a

dielectric material. To this end we attempt to minimize the difference between simulations and observations of time-domain data. In our formulation, the observations, \hat{E}_j , are of the electric field E at discrete times t_j taken at $z = 0$. Each simulation is a solution of Maxwell's equations given in (12) using the polarization model (14) with candidate values for the distribution of relaxation parameters. We propose a non-standard least-squares measurement (see [13] for a detailed justification) for the objective function, which is given by

$$J(F) = \sum_j \left| |E(t_j, 0; F)| - |\hat{E}_j| \right|^2, \quad (17)$$

where $E(\cdot, \cdot; F)$ is the solution to (12) and (14) corresponding to the distribution F . Thus the inverse problem is to solve

$$\min_{F \in \mathcal{P}(\mathcal{T})} J(F).$$

5 Inverse Problem Results Using a Discrete Distribution

Consider a discrete distribution with two atoms at τ_1 and τ_2 . Essentially, we are decomposing the electric polarization into two components, each dependent on distinct relaxation times as follows:

$$P = \alpha_1 \mathcal{P}(t, z; \tau_1) + \alpha_2 \mathcal{P}(t, z; \tau_2), \quad (18)$$

where each $\mathcal{P}(\cdot, \cdot; \tau_i)$ satisfies a Debye polarization equation with parameter τ_i . For now we assume the proportions α_1 and $\alpha_2 = 1 - \alpha_1$ are known. Thus we are attempting to solve the following least squares optimization problem:

$$\min_{(\tau_1, \tau_2)} \sum_j \left| |E(t_j, 0; (\tau_1, \tau_2))| - |\hat{E}_j| \right|^2, \quad (19)$$

where \hat{E}_j is synthetic data generated using the true solution (τ_1^*, τ_2^*) in our simulator (with a highly refined mesh), and $E(t_j, 0; (\tau_1, \tau_2))$ depends on each τ_i through its dependence on P , see for example (12). Figures 2 and 3 depict an example of the objective function and the log of the objective function respectively, both plotted versus the logs of τ_1 and τ_2 (using a frequency of $10^{11} Hz$, $\alpha_1 = \alpha_2 = .5$, $\tau_1^* = 10^{-7.5}$ and $\tau_2^* = 10^{-7.8}$).

5.1 Analysis of Objective Function

We can see clearly from the log surface plot in Figure 3 that there exists a relation for the relaxation times for which the corresponding simulations best match the data. We will refer to this relation as the “curve of best fit”.

Note that the appearance of many local minima is due to the steep decent near the “curve of best fit” since the lattice points of the mesh used do not always lie near the curve. If we trace along this curve, as displayed in Figure 4, we see that there are actually two global minima, the exact solution of $\log(\tau_1) = -7.5$ and $\log(\tau_2) = -7.8$, and since the proportion used in this case was $\alpha_1 = \alpha_2 = .5$, we also have the symmetric solution where τ_1 and τ_2 are swapped. If we superimpose the data from the contour plot and the data from the surface plot on a lattice, as demonstrated in Figure 5, we can more easily see the structure of the “curve of best fit”, and both global minima.

Unfortunately, the scale of these plots shows that the difference between the objective function at the minimizers and any other point on this curve in our parameter space is less than 4×10^{-10} . Therefore the exact minimizing parameters are not likely to be identifiable in a practical, experimental setting.

The equation for this “best fit curve” can be derived by combining equations (12) and (14). Note that

$$\ddot{P}(t, z) = \int_0^t \ddot{\mathcal{G}}(t - s, z) E(s, z) ds + \mathcal{G}(0, z) \dot{E}(t, z) + \dot{\mathcal{G}}(0, z) E(t, z), \quad (20)$$

where

$$\ddot{\mathcal{G}}(t, z; F) = \int_{\mathcal{T}} \ddot{g}(t, z; \tau) dF(\tau).$$

So that now substituting (20) into (12) we obtain

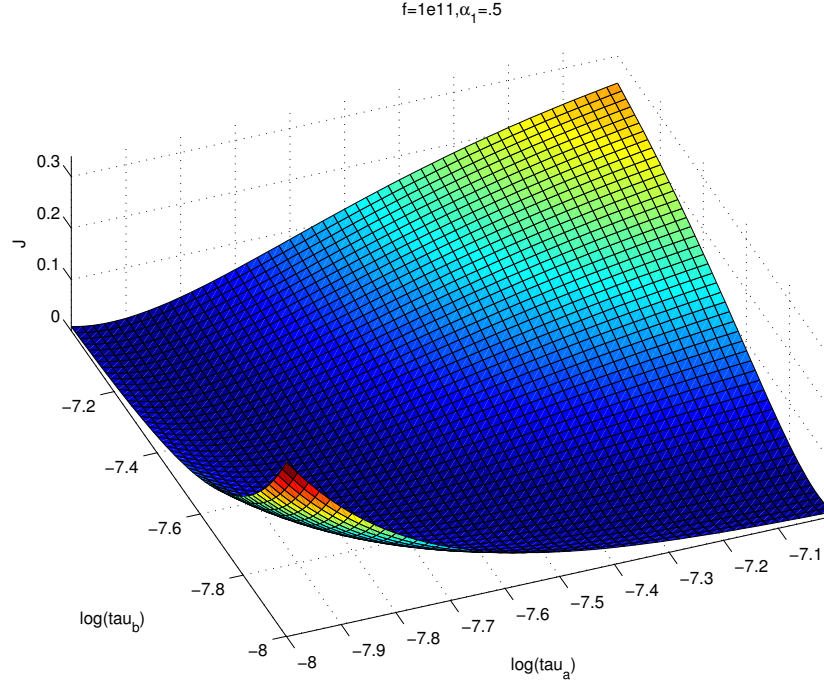


Figure 2: The objective function for the relaxation time inverse problem versus the log of τ_1 and the log of τ_2 using a frequency of $10^{11} Hz$.

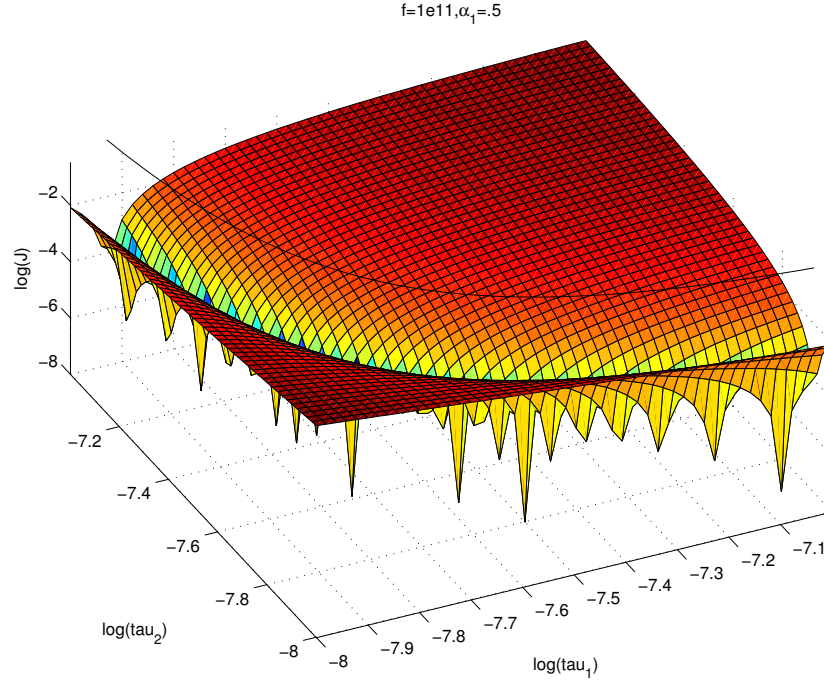


Figure 3: The log of the objective function for the relaxation time inverse problem versus the log of τ_1 and the log of τ_2 using a frequency of $10^{11} Hz$. The solid line above the surface represents the curve of constant $\tilde{\lambda}$.

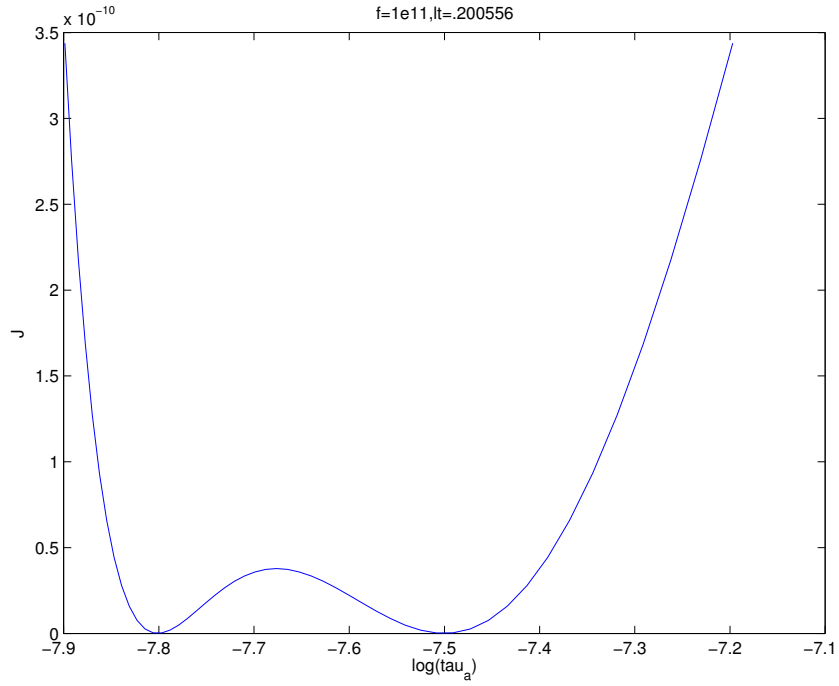


Figure 4: The objective function for the relaxation time inverse problem plotted along the curve of constant $\tilde{\lambda}$ using a frequency of $10^{11} Hz$.

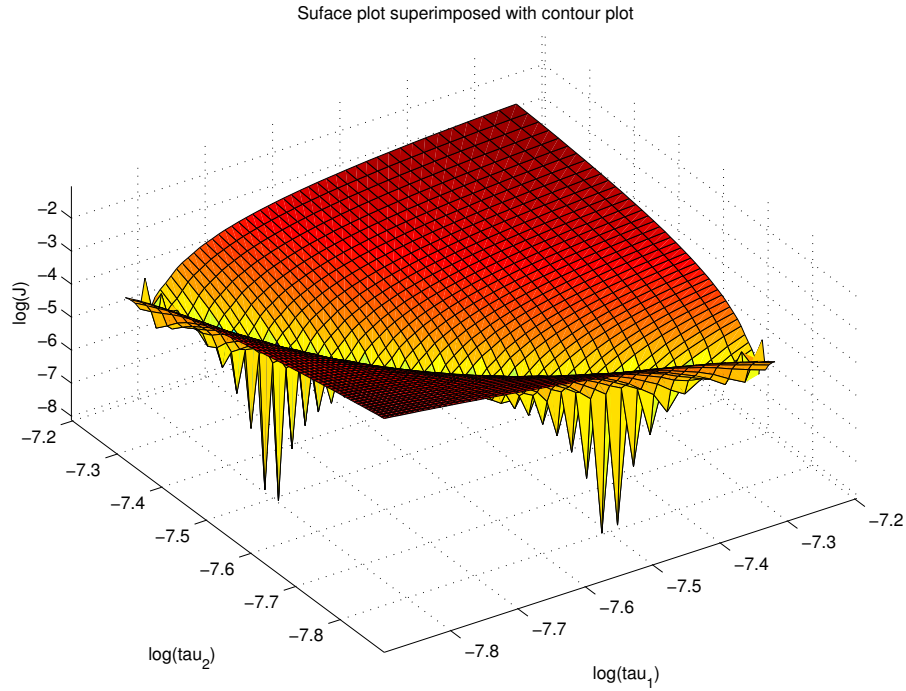


Figure 5: The surface plot of the objective function for the relaxation time inverse problem superimposed with the contour plot along the curve of constant $\tilde{\lambda}$, using a frequency of $10^{11} Hz$. The fact that there are only two minima is clearly more visible in this plot.

$$\begin{aligned} \langle \epsilon_r \ddot{E}(t, \cdot), \phi \rangle + \langle I_\Omega [\eta_0 \sigma + \mathcal{G}(0, z)] \dot{E}(t, \cdot), \phi \rangle + \langle I_\Omega \dot{\mathcal{G}}(0, \cdot) E(t, \cdot), \phi \rangle \\ + \langle \int_0^t I_\Omega \ddot{\mathcal{G}}(t-s, \cdot) E(s, \cdot) ds, \phi \rangle + \langle E'(t, \cdot), \phi' \rangle + \dot{E}(t, 0) \phi(0) = -\langle \eta_0 \dot{J}_s(t, \cdot), \phi \rangle. \end{aligned} \quad (21)$$

For large frequencies, the \dot{E} term in (21) dominates over all other terms which depend on \mathcal{G} . We can see this by considering the frequency domain where the time derivative causes an increase by a factor of ω in the norm. Thus the equation for the “curve of best fit” is simply that of constant $\mathcal{G}(0, z)$. For the discrete Debye example we have

$$\mathcal{G}(0, z) = \int_{\mathcal{T}} g(0, z; \tau) dF(\tau) = \alpha_1 \epsilon_0 (\epsilon_s - \epsilon_\infty) / \tau_1 + \alpha_2 \epsilon_0 (\epsilon_s - \epsilon_\infty) / \tau_2,$$

therefore the equation is

$$\frac{\alpha_1}{\tau_1} + \frac{\alpha_2}{\tau_2} = \frac{\alpha_1}{\tau_1^*} + \frac{\alpha_2}{\tau_2^*}. \quad (22)$$

This is precisely the curve that is plotted above the surface in Figure 3.

For scaling purposes, and to be consistent with the presentation in [11], we define

$$\tilde{\lambda} := \frac{1}{c\epsilon_0(\epsilon_s - \epsilon_\infty)} \mathcal{G}(0, z),$$

which in this example means

$$\tilde{\lambda} = \frac{\alpha_1}{c\tau_1} + \frac{\alpha_2}{c\tau_2}$$

and thus we may say that the “curve of best fit” is the line of constant $\tilde{\lambda}$.

The frequency dependence of the \dot{E} term in (21) suggests that for smaller frequencies it may not be the dominant contributor, and therefore, there may be a fundamentally different structure to the surface plot. This is in fact what we observe in our simulations. Figures 6 through 9 display the surface plots and the log surface plots for frequencies of $10^9 Hz$ and $10^6 Hz$. Note that in the latter case the concavity of the “curve of best fit” has swapped orientation! This demonstrates that the surface plots are very much dependent on ω even though the relaxation mechanisms are the same.

Through our numerical calculations we have determined that for the case using a frequency of $10^6 Hz$ the “curve of best fit” is actually that of constant $\tilde{\tau} := \alpha_1 \tau_1 + \alpha_2 \tau_2$ which is what one might expect as this is the weighted average of the relaxation times. For the example given here, $\tilde{\tau} \approx 2.37000 \times 10^{-8}$. In our simulations, angular frequencies less than one divided by this number were characterized by a constant $\tilde{\tau}$ while larger frequencies resulted in dominance by the $\tilde{\lambda}$ term given above. The fact that the regime characterized by $\omega\tau < 1$ is fundamentally different in many respects from that of the $\omega\tau > 1$ regime is well documented (see, for example, [23]). Still, the behavior of the objective function along its corresponding “curve of best fit” is similar for each frequency, despite the curves themselves being fundamentally different, as demonstrated by comparing Figure 4 to Figure 10 which uses a frequency of $10^6 Hz$. Note, however, that the scale of Figure 10 is several orders of magnitudes larger, suggesting that the global minimizers may be easier to find for smaller frequencies.

Remark 1 *We should mention here that in order to analyze only the relationship between the interrogating frequency and the relaxation time, we scaled the dimensions of the interrogated object according to the change in scale of frequency. For example, when the frequency is divided by 100 to go from 10^{11} to 10^9 , we accordingly multiply the slab thickness by 100 to get .4m in order that the ratio of the dimensions to the wavelength of the signal remains the same.*

5.2 Optimization Procedure and Results

We attempt to apply a two parameter Levenberg-Marquardt optimization routine to the modified least squares error between the given data and our simulations, as defined in (17), to try to identify the two distinct

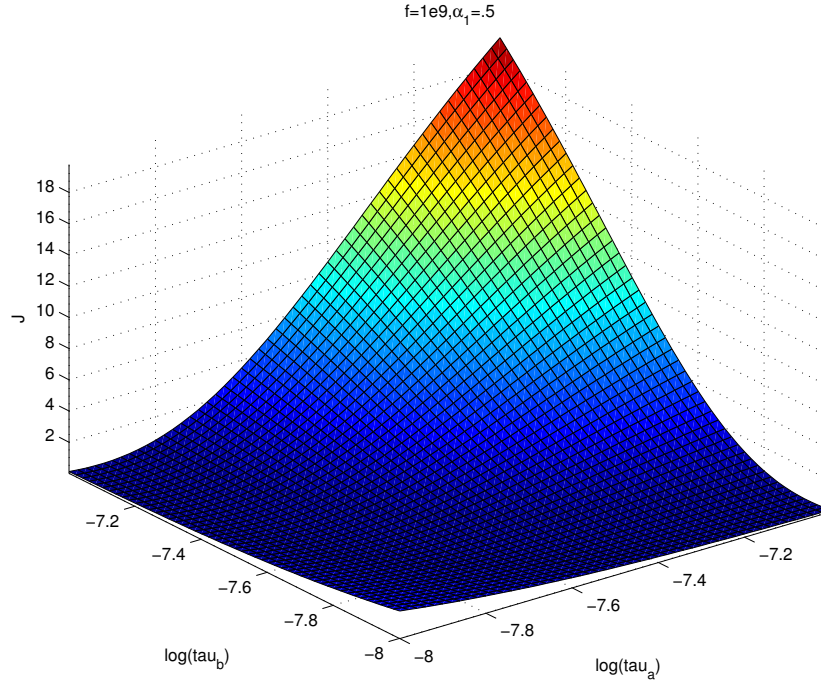


Figure 6: The objective function for the relaxation time inverse problem versus the log of τ_1 and the log of τ_2 using a frequency of $10^9 Hz$.

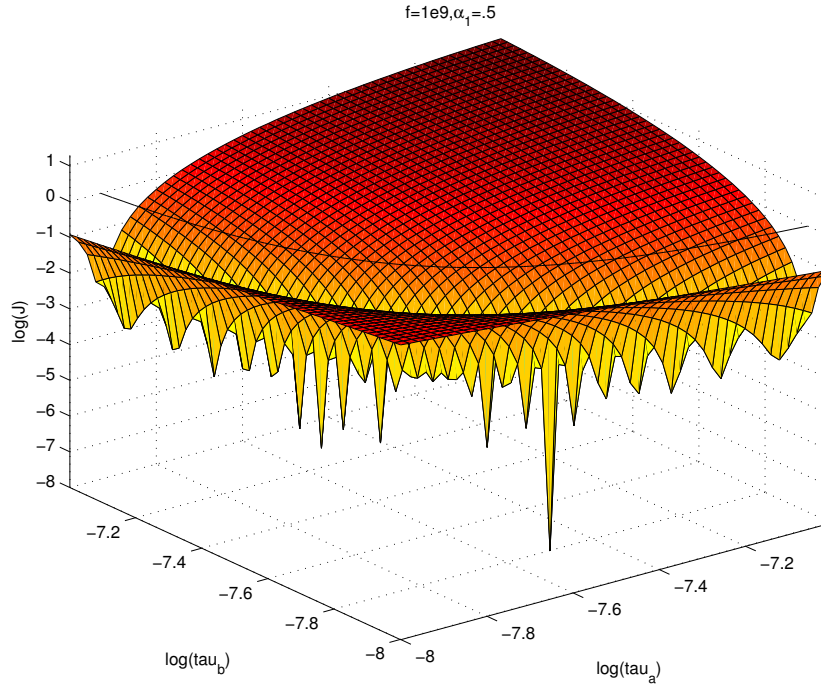


Figure 7: The log of the objective function for the relaxation time inverse problem versus the log of τ_1 and the log of τ_2 using a frequency of $10^9 Hz$. The solid line above the surface represents the curve of constant $\tilde{\lambda}$.

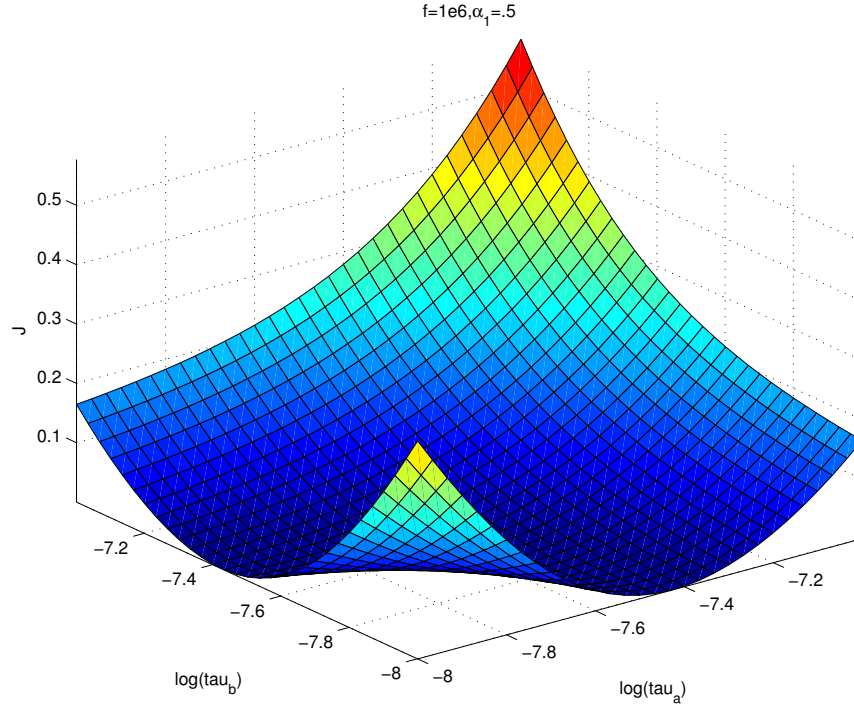


Figure 8: The objective function for the relaxation time inverse problem versus the log of τ_1 and the log of τ_2 using a frequency of $10^6 Hz$.

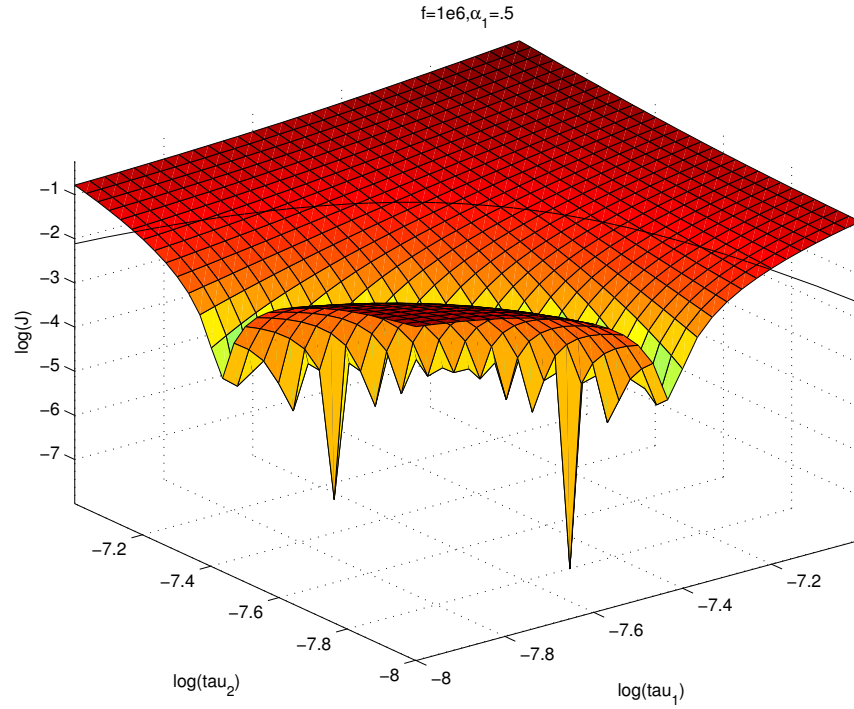


Figure 9: The log of the objective function for the relaxation time inverse problem versus the log of τ_1 and the log of τ_2 using a frequency of $10^6 Hz$. The solid line above the surface represents the curve of constant $\tilde{\tau}$.

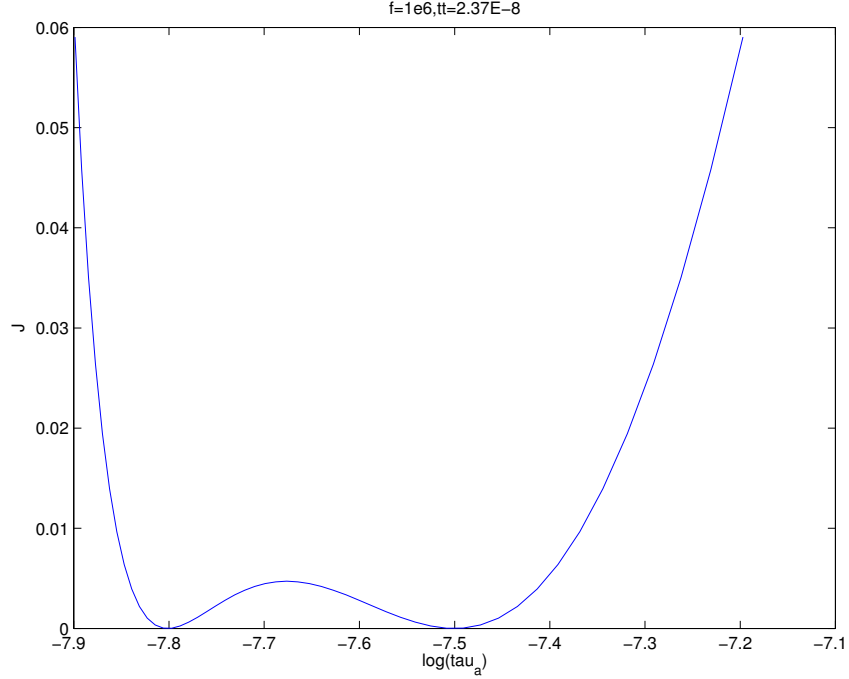


Figure 10: The objective function for the relaxation time inverse problem plotted along the curve of constant $\tilde{\tau}$ using a frequency of $10^6 Hz$.

relaxation times that generated the data. We are assuming that the corresponding volume proportions of the two materials (α_1 and $\alpha_2 := 1 - \alpha_1$) are known. We consider three different scenarios with respect to the volume proportions: $\alpha_1 \in \{.1, .5, .9\}$. We also perform our inverse problem using the frequencies $10^{11} Hz$, $10^9 Hz$ and $10^6 Hz$. Lastly we test the optimization procedure with three various initial conditions given in Table 1. The actual values are $\tau_1 = 10^{-7.50031} \approx 3.16 \times 10^{-8}$ and $\tau_2 = 10^{-7.80134} \approx 1.58 \times 10^{-8}$. Note that the first set of initial conditions is the farthest from the exact solution, while the third is the closest.

Table 1: Three sets of initial conditions for the relaxation time inverse problem representing $(\tau_1^0, \tau_2^0) = (C\tau_1^*, \tau_2^*/C)$ for $C \in \{5, 2, 1.25\}$ respectively (case 0 represents exact solution), also given are the \log_{10} of each relaxation time, as well as the % relative error from the exact value.

case	τ_1	τ_2	$\log(\tau_1)$	$\log(\tau_2)$	% τ_1	% τ_2
0	3.1600e-8	1.5800e-8	-7.50031	-7.80134	0	0
1	1.5800e-7	3.1600e-9	-6.80134	-8.50031	400	80
2	6.3200e-8	7.9000e-9	-7.19928	-8.10237	100	50
3	3.9500e-8	1.2640e-8	-7.40340	-7.89825	25	20

The results of the optimization are given in Tables 2 through 7. Most of the cases appear not to have converged. Only a few of the cases corresponding to the closest initial conditions converged close to the original values of the relaxation times. But if we recall the shape of the objective function that we are trying to minimize (for example see Figures 2 and 3) then it is understandable how a gradient based method may converge directly to the “curve of best fit”. However, once it reaches this curve, the optimization routine may jump back and forth across the “trench” and may not be able to traverse the curve to find the true global minimum. This is in fact what is occurring.

Table 2: Resulting values of τ_1 from the Levenberg-Marquardt routine using a frequency of $10^{11}Hz$ (recall the exact solution $\tau_1^* = 3.1600e-8$).

case	α_1		
	.1	.5	.9
1	1.57218e-07	1.7488e-07	2.26538e-07
2	6.34973e-08	9.73445e-08	4.90169e-08
3	3.98062e-08	4.13321e-08	3.43569e-08

Table 3: Resulting values of τ_1 from the Levenberg-Marquardt routine using a frequency of 10^9Hz (recall the exact solution $\tau_1^* = 3.1600e-8$).

case	α_1		
	.1	.5	.9
1	1.60128e-07	1.60397e-07	7.0395e-08
2	6.37815e-08	6.68561e-08	3.99819e-08
3	3.97416e-08	4.03239e-08	3.59503e-08

Table 4: Resulting values of τ_1 from the Levenberg-Marquardt routine using a frequency of 10^6Hz (recall the exact solution $\tau_1^* = 3.1600e-8$).

case	α_1		
	.1	.5	.9
1	3.79957e-08	3.23393e-08	3.16236e-08
2	3.17753e-08	3.32218e-08	3.21036e-08
3	3.17897e-08	3.19001e-08	3.16068e-08

Table 5: Resulting values of τ_2 from the Levenberg-Marquardt routine using a frequency of $10^{11}Hz$ (recall the exact solution $\tau_2^* = 1.5800e-8$).

case	α_1		
	.1	.5	.9
1	1.51271e-08	1.1208e-08	3.2429e-09
2	1.53697e-08	1.18119e-08	6.10283e-09
3	1.56197e-08	1.41322e-08	1.1607e-08

Table 6: Resulting values of τ_2 from the Levenberg-Marquardt routine using a frequency of 10^9Hz (recall the exact solution $\tau_2^* = 1.5800e-8$).

case	α_1		
	.1	.5	.9
1	1.51255e-08	1.12739e-08	4.54532e-09
2	1.53692e-08	1.25029e-08	8.12599e-09
3	1.56222e-08	1.4257e-08	1.02271e-08

Table 7: Resulting values of τ_2 from the Levenberg-Marquardt routine using a frequency of 10^6Hz (recall the exact solution $\tau_2^* = 1.5800e-8$).

case	α_1		
	.1	.5	.9
1	1.51064e-08	1.50634e-08	1.63437e-08
2	1.57826e-08	1.41845e-08	1.12902e-08
3	1.57792e-08	1.55032e-08	1.57796e-08

Table 8 displays the values of $\tilde{\lambda}$ corresponding to each set of initial conditions, where case 0 denotes the exact solution. Tables 9 and 10 display the resulting values of $\tilde{\lambda}$ from the Levenberg-Marquardt routine using frequencies $10^{11}Hz$ and 10^9Hz respectively. Since smaller frequencies are expected to converge to the curve of constant $\tilde{\tau}$, we give the initial values of $\tilde{\tau}$ as well, in Table 11, and Table 12 displays the resulting values of $\tilde{\tau}$ after running the Levenberg-Marquardt routine using a frequency of 10^6Hz . Note that while the initial values in case 1 are farthest from the exact solution, some of the corresponding $\tilde{\tau}$ are actually closest to the actual value (for example, $\alpha_1 = .1$). While this suggests that this should more easily converge to the “curve of best fit”, it will most likely be farther away from the exact solution on this curve, and therefore should still be considered the hardest of the cases to solve.

Recognizing that we have converged to the “curve of best fit” in each of the above cases, we may now restart an optimization routine that traverses this curve. Some modifications to the optimization routine’s parameters are required to address the vast difference in scales of this subproblem as compared to the two parameter optimization problem. The final results of this two step optimization approach are given in Tables 13 through 18. As expected, the results from case 3 are generally better than the other two sets of initial conditions. The highest frequency attempted, $10^{11}Hz$ seemed to perform the most poorly. This suggests that the higher the frequency, the more difficult to accurately resolve the polarization mechanisms. Although the 10^6Hz case used the curve of constant $\tilde{\tau}$ while the 10^9Hz case used the curve of constant $\tilde{\lambda}$, there was no evidence to suggest that one case performed better than the other. Lastly, it appears that when material 1 (corresponding to relaxation parameter τ_1) is of the highest proportion, the optimization routine is best able to resolve τ_1 . Likewise, if material 1 is of a lower proportion, the routine instead does a better job of resolving τ_2 . Note that there are several instances where the optimizer “switched” τ_1 and τ_2 , for example in case 1 of frequency 10^9Hz with $\alpha_1 = .5$. In this scenario each material comprises 50% of the whole, so the problem is symmetric and swapping τ_1 and τ_2 has no effect. However, in case 1 of frequency 10^9Hz with $\alpha_1 = .1$, it appears that τ_1 is converging toward the exact τ_2^* value, but since this problem is not symmetric, τ_2 converges to a meaningless solution.

Table 8: The initial values of $\tilde{\lambda} := \frac{\alpha_1}{c\tau_1} + \frac{\alpha_2}{c\tau_2}$ for each set of initial conditions (case 0 represents the exact solution).

case	.1	.5	.9
0	0.200556	0.158333	0.116111
1	0.952112	0.538334	0.124556
2	0.385278	0.237500	0.0897223
3	0.245945	0.174167	0.102389

Table 9: Resulting values of $\tilde{\lambda}$ from the Levenberg-Marquardt routine using a frequency of $10^{11}Hz$ for each set of initial conditions (case 0 represents the exact solution).

case	α_1		
	.1	.5	.9
0	0.200556	0.158333	0.116111
1	0.200573	0.15834	0.116109
2	0.200573	0.158327	0.1159
3	0.200573	0.158363	0.116114

Table 10: Resulting values of $\tilde{\lambda}$ from the Levenberg-Marquardt routine using a frequency of 10^9Hz for each set of initial conditions (case 0 represents the exact solution).

case	α_1		
	.1	.5	.9
0	0.200556	0.158333	0.116111
1	0.200555	0.158331	0.116029
2	0.200556	0.158337	0.116132
3	0.200556	0.158339	0.116119

Table 11: The initial values of $\tilde{\tau} := \alpha_1\tau_1 + \alpha_2\tau_2$ for each set of initial conditions (case 0 represents the exact solution).

case	.1	.5	.9
0	1.7380e-08	2.3700e-08	3.0020e-08
1	1.8644e-08	8.0580e-08	1.42516e-07
2	1.3430e-08	3.5550e-08	5.7670e-08
3	1.5326e-08	2.6070e-08	3.6814e-08

Table 12: Resulting values of $\tilde{\tau}$ from the Levenberg-Marquardt routine using a frequency of 10^6Hz for each set of initial conditions (case 0 represents the exact solution).

case	α_1		
	.1	.5	.9
0	1.7380e-08	2.3700e-08	3.0020e-08
1	1.73954e-08	2.37014e-08	3.00956e-08
2	1.73819e-08	2.37031e-08	3.00222e-08
3	1.73803e-08	2.37016e-08	3.00241e-08

Table 13: Final estimates for τ_1 from two step optimization approach using a frequency of $10^{11}Hz$ for each set of initial conditions (recall the exact solution $\tau_1^* = 3.1600e-8$).

	α_1		
case	.1	.5	.9
1	1.57221e-07	2.81579e-08	2.61108e-08
2	6.34951e-08	4.63981e-08	3.15232e-08
3	3.98043e-08	3.33803e-08	3.2833e-08

Table 14: Final estimates for τ_1 from two step optimization approach using a frequency of 10^9Hz for each set of initial conditions (recall the exact solution $\tau_1^* = 3.1600e-8$).

	α_1		
case	.1	.5	.9
1	1.12709e-08	1.58005e-08	3.15929e-08
2	3.1574e-08	3.16002e-08	3.16017e-08
3	3.16009e-08	3.16009e-08	3.16009e-08

Table 15: Final estimates for τ_1 from two step optimization approach using a frequency of 10^6Hz for each set of initial conditions (recall the exact solution $\tau_1^* = 3.1600e-8$).

	α_1		
case	.1	.5	.9
1	3.17673e-08	3.16031e-08	3.16986e-08
2	3.16206e-08	3.16068e-08	3.16031e-08
3	3.16031e-08	3.16039e-08	3.16053e-08

Table 16: Final estimates for τ_2 from two step optimization approach using a frequency of $10^{11}Hz$ for each set of initial conditions (recall the exact solution $\tau_2^* = 1.5800e-8$).

	α_1		
case	.1	.5	.9
1	1.51273e-08	1.68275e-08	2.93245e-07
2	1.53699e-08	1.36276e-08	1.61391e-08
3	1.56196e-08	1.53854e-08	1.35142e-08

Table 17: Final estimates for τ_2 from two step optimization approach using a frequency of 10^9Hz for each set of initial conditions (recall the exact solution $\tau_2^* = 1.5800e-8$).

	α_1		
case	.1	.5	.9
1	1.75594e-08	3.15995e-08	1.58782e-08
2	1.58008e-08	1.57994e-08	1.57801e-08
3	1.58001e-08	1.5799e-08	1.57925e-08

Table 18: Final estimates for τ_2 from two step optimization approach using a frequency of 10^6Hz for each set of initial conditions (recall the exact solution $\tau_2^* = 1.5800e-8$).

	α_1		
case	.1	.5	.9
1	1.52384e-08	1.52286e-08	1.61666e-08
2	1.5787e-08	1.45007e-08	1.18861e-08
3	1.57845e-08	1.55744e-08	1.5783e-08

5.3 Determination of Volume Proportions

We now attempt to apply a one parameter Levenberg-Marquardt optimization routine to our problem to identify the relative amounts of two materials with known, distinct relaxation times. Thus we are trying to find the corresponding volume proportions of the two materials (α_1 and $\alpha_2 := 1 - \alpha_1$). We again consider the three scenarios with respect to the exact volume proportions: $\alpha_1 \in \{.1, .5, .9\}$. We also perform our inverse problem using the frequencies $10^{11}Hz$, 10^9Hz and 10^6Hz . Lastly we test the optimization procedure with three various initial conditions: $\alpha_1^0 \in \{.9999, .0001, .5\}$ (except in the case when $\alpha_1^* = .5$ in which case we used $\alpha_1^0 \in \{.9999, .0001, .4\}$). We refer to these as Cases 1, 2, and 3, respectively. In all of the following we assume that the known relaxation times are $\tau_1 = 10^{-7.50031} \approx 3.16 \times 10^{-8}$ and $\tau_2 = 10^{-7.80134} \approx 1.58 \times 10^{-8}$.

Figures 11 through 13 depict the graphs of the functions that we are attempting to minimize. For each case the curves appear well behaved. The results for this one parameter inverse problem, displayed in Tables 19 through 21, verify that the relative proportions of known materials are generally easily identifiable. Tables 22 through 24 display the final objective function values for each case. Note that typical initial values for J were around 0.1, further, the tolerance was set at 10^{-9} , thus all but a few cases converged before reaching the maximum of 20 iterations.

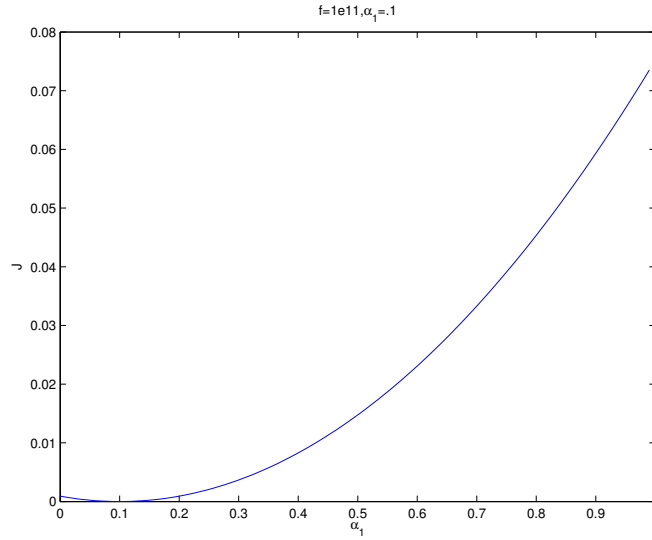


Figure 11: The objective function for the relaxation time inverse problem versus α_1 using a frequency of $10^{11}Hz$ and $\alpha_1^* = .1$.

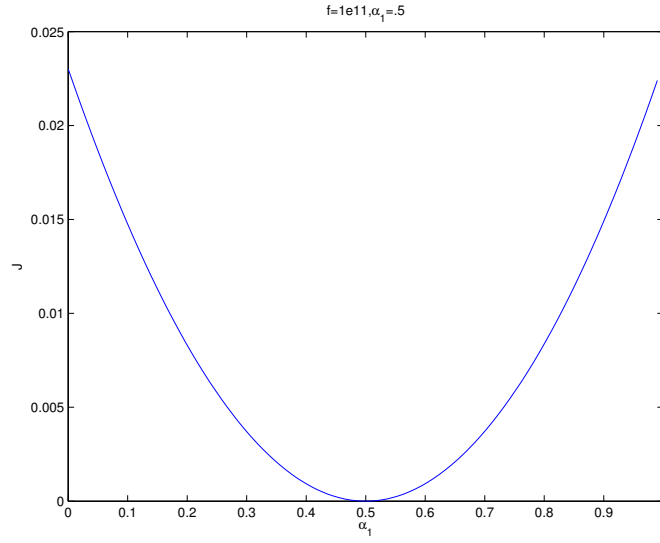


Figure 12: The objective function for the relaxation time inverse problem versus α_1 using a frequency of $10^{11} Hz$ and $\alpha_1^* = .5$.

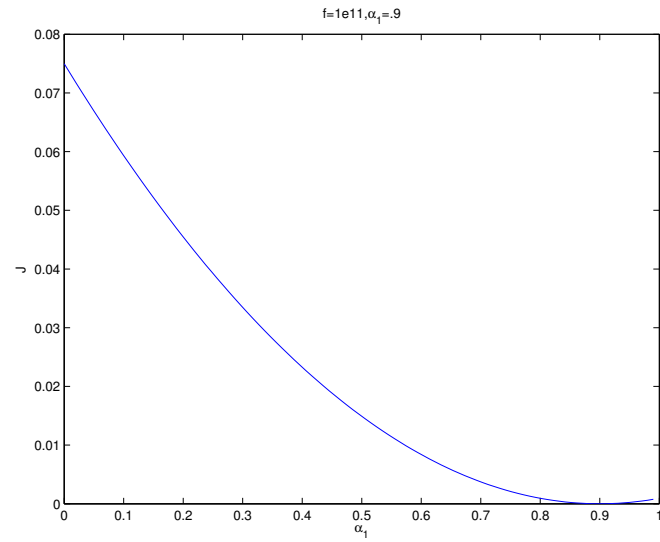


Figure 13: The objective function for the relaxation time inverse problem versus α_1 using a frequency of $10^{11} Hz$ and $\alpha_1^* = .9$.

Table 19: Results for the one parameter inverse problem to determine the relative proportion of two known Debye materials using a frequency of $10^{11}Hz$ (α_1 estimates are shown).

case	α_1		
	.1	.5	.9
1	0.10095	0.501353	0.900013
2	0.0999994	0.5	0.9
3	0.100643	0.499994	0.899994

Table 20: Results for the one parameter inverse problem to determine the relative proportion of two known Debye materials using a frequency of 10^9Hz (α_1 estimates are shown).

case	α_1		
	.1	.5	.9
1	0.10086	0.5	0.9
2	0.0999161	0.5	0.9
3	0.100017	0.5	0.9

Table 21: Results for the one parameter inverse problem to determine the relative proportion of two known Debye materials using a frequency of 10^6Hz (α_1 estimates are shown).

case	α_1		
	.1	.5	.9
1	0.0995386	0.5	0.900008
2	0.0999965	0.499969	0.899996
3	0.100013	0.499999	0.899999

Table 22: Final objective function values for the inverse problem to determine the relative proportion of two known Debye materials using a frequency of $10^{11}Hz$.

case	α_1		
	.1	.5	.9
1	8.27897e-08	1.69802e-07	1.81577e-11
2	3.46217e-14	1.64325e-18	1.8982e-15
3	3.79016e-08	3.76649e-12	2.91698e-12

Table 23: Final objective function values for the inverse problem to determine the relative proportion of two known Debye materials using a frequency of 10^9Hz .

case	α_1		
	.1	.5	.9
1	1.63248e-05	1.25127e-12	6.89036e-14
2	1.55308e-07	5.89446e-13	5.8576e-14
3	6.53272e-09	2.37274e-12	5.77714e-15

Table 24: Final objective function values for the inverse problem to determine the relative proportion of two known Debye materials using a frequency of 10^6Hz .

case	α_1		
	.1	.5	.9
1	7.74602e-08	7.31725e-15	4.39638e-12
2	4.51572e-12	1.40031e-10	1.03552e-12
3	6.11308e-11	2.0608e-13	1.53426e-13

5.4 Determination of Volume Proportions and Relaxation Times Simultaneously

Although we anticipate this problem formulation to be underdetermined, we run our optimization routine for the problem where neither the relaxation times, nor the relative volume proportions, of two distinct Debye materials are known. Thus this is a three parameter inverse problem for $\{\tau_1, \tau_2, \alpha_1\}$ (since $\alpha_2 = 1 - \alpha_1$). Again we consider the following scenarios with respect to the actual volume proportions: $\alpha_1 \in \{.1, .5, .9\}$. We also perform our inverse problem using the frequencies $10^{11}Hz$, 10^9Hz and 10^6Hz . Lastly we test the optimization procedure with the three various initial conditions that were given in Table 1. Note that the actual values are still $\tau_1 = 10^{-7.50031} \approx 3.16 \times 10^{-8}$ and $\tau_2 = 10^{-7.80134} \approx 1.58 \times 10^{-8}$.

Our initial condition for the volume distribution is $\alpha_1^0 = .9999$, which means there is essentially only material 1. We chose this initial condition mainly to see whether the optimization routine would try to compensate for the altered volume distribution by significantly changing the relaxation times and not sufficiently correcting the volume distribution. This is in fact what seems to have occurred. The final α_1 values are given in Tables 25 through 27, while the final τ_1 and τ_2 values are in Tables 28 through 30 and 31 through 33 respectively. Tables 34 through 36 display the final objective function values. Any values of the final objective function under the tolerance of 10^{-9} should be considered as indicative of convergence, while values greater than 1 clearly indicate stagnation, (e.g., 10^9Hz : case 1).

Table 25: Resulting values of α_1 for the underdetermined inverse problem using a frequency of $10^{11}Hz$.

case	α_1^*		
	.1	.5	.9
1	0.606255	0.881722	0.912728
2	0.761665	0.817744	0.877667
3	0.818091	0.879896	0.942175

Table 26: Resulting values of α_1 for the underdetermined inverse problem using a frequency of 10^9Hz .

case	α_1		
	.1	.5	.9
1	0.941041	0.941091	0.880293
2	0.999811	0.993342	0.994368
3	0.990988	0.96514	0.965754

Table 27: Resulting values of α_1 for the underdetermined inverse problem using a frequency of 10^6Hz .

case	α_1		
	.1	.5	.9
1	0.687959	0.746478	0.956578
2	0.83678	0.900072	0.908588
3	0.890569	0.841634	0.945684

Of the cases that converged, none converged to the correct solution. For example, the final values of α_1 were all greater than .6, (although the estimates corresponding to a smaller α_1^* were on average less than the estimates corresponding to $\alpha_1^* = .9$). When looking at the final estimates for the relaxation times, however, there appears to be no rhyme nor reason. This is quite similar to what happened when we first tried to determine unknown relaxation times, but with a fixed volume proportion, in Section 5. Thus we may expect that here again no single parameter, τ_1 , τ_2 , nor α_1 , may converge, but instead we may see convergence of a general relation involving them all, namely $\tilde{\lambda} := \frac{\alpha_1}{c\tau_1} + \frac{\alpha_2}{c\tau_2}$ (or respectively $\tilde{\tau} := \alpha_1\tau_1 + \alpha_2\tau_2$ for angular frequencies less than $\frac{1}{\tilde{\tau}}$).

We redisplay the exact values of $\tilde{\lambda}$ and $\tilde{\tau}$ in Table 37 for reference. Also, the initial values of $\tilde{\lambda}$ and $\tilde{\tau}$ for each case of initial conditions (recall Table 1) is given in Table 38. Finally, Tables 39 through 41 display the error of the final resulting $\tilde{\lambda}$ (or $\tilde{\tau}$ appropriately) estimates from the exact solutions (i.e., the absolute

Table 28: Resulting values of τ_1 for the underdetermined inverse problem using a frequency of $10^{11}Hz$ (recall the exact solution $\tau_1^* = 3.1600e-8$).

case	α_1		
	.1	.5	.9
1	3.9469e-08	5.28765e-07	3.23635e-08
2	1.52948e-08	2.49314e-08	2.60749e-08
3	1.79134e-08	2.40213e-08	3.06153e-08

Table 29: Resulting values of τ_1 for the underdetermined inverse problem using a frequency of 10^9Hz (recall the exact solution $\tau_1^* = 3.1600e-8$).

case	α_1		
	.1	.5	.9
1	1.2046e-07	1.24415e-07	5.50848e-07
2	1.67567e-08	2.16631e-08	2.93323e-08
3	1.65557e-08	2.25001e-08	3.0285e-08

Table 30: Resulting values of τ_1 for the underdetermined inverse problem using a frequency of 10^6Hz (recall the exact solution $\tau_1^* = 3.1600e-8$).

case	α_1		
	.1	.5	.9
1	2.12522e-08	2.81013e-08	3.10195e-08
2	1.91221e-08	2.3216e-08	3.15199e-08
3	2.58248e-08	2.67889e-08	3.11248e-08

values of the differences). We show the errors instead of the actual values because some cases converged so well that the number of decimal places required to see discrepancies is impractical to show! Thus while the inverse problem was unable to accurately resolve the individual values τ_1 , τ_2 , and α_1 , in all but a few cases (particularly in the 10^9Hz scenario) the optimization routine did converge as well as to be expected to the “curve of best fit”. In the previous section, with the fixed volume proportions, we were able to traverse the curve of constant $\tilde{\lambda}$ (or $\tilde{\tau}$ appropriately) to find the global minimum. However, here we truly have an underdetermined problem and therefore we cannot extract any further information without providing more data.

Table 31: Resulting values of τ_2 for the underdetermined inverse problem using a frequency of $10^{11}Hz$ (recall the exact solution $\tau_2^* = 1.5800e-8$).

case	α_1		
	.1	.5	.9
1	8.79556e-09	2.58245e-09	1.32076e-08
2	2.30775e-08	1.24249e-08	1.06318e-07
3	1.25825e-08	1.10813e-08	1.4329e-08

Table 32: Resulting values of τ_2 for the underdetermined inverse problem using a frequency of 10^9Hz (recall the exact solution $\tau_2^* = 1.5800e-8$).

case	α_1		
	.1	.5	.9
1	1.81687e-10	2.32623e-10	1.2353e-09
2	4.08525e-10	4.11677e-09	6.18566e-09
3	3.36629e-08	7.61993e-09	1.17231e-08

Table 33: Resulting values of τ_2 for the underdetermined inverse problem using a frequency of 10^6Hz (recall the exact solution $\tau_2^* = 1.5800e-8$).

case	α_1		
	.1	.5	.9
1	8.82441e-09	1.06962e-08	7.87058e-09
2	8.40313e-09	2.78555e-08	1.5107e-08
3	6.19113e-09	7.15491e-09	1.0713e-08

Table 34: Resulting values of the objective function J for the underdetermined inverse problem using a frequency of $10^{11}Hz$.

case	α_1^*		
	.1	.5	.9
1	1.87304e-09	1.39638e-07	2.96012e-12
2	1.87016e-15	3.3089e-14	1.60601e-17
3	2.79754e-14	2.80156e-16	1.38435e-13

Table 35: Resulting values of the objective function J for the underdetermined inverse problem using a frequency of 10^9Hz .

case	α_1		
	.1	.5	.9
1	9.59505	13.787	25.045
2	0.00021558	2.49585e-06	2.03659e-07
3	3.48546e-06	5.0347e-07	1.18429e-08

Table 36: Resulting values of the objective function J for the underdetermined inverse problem using a frequency of 10^6Hz .

case	α_1		
	.1	.5	.9
1	1.19225e-05	4.76542e-07	1.18516e-07
2	6.34493e-07	3.86474e-05	4.41757e-10
3	4.07294e-06	1.99923e-06	2.75692e-08

Table 37: The exact values of $\tilde{\lambda} := \frac{\alpha_1}{c\tau_1} + \frac{\alpha_2}{c\tau_2}$ (first row) and $\tilde{\tau} := \alpha_1\tau_1 + \alpha_2\tau_2$ (second row) for each set of volume distributions.

	α_1^*		
	.1	.5	.9
$\tilde{\lambda}^*$	0.200556	0.158333	0.116111
$\tilde{\tau}^*$	1.7380e-08	2.3700e-08	3.0020e-08

Table 38: The initial values of $\tilde{\lambda} := \frac{\alpha_1}{c\tau_1} + \frac{\alpha_2}{c\tau_2}$ (first column) and $\tilde{\tau} := \alpha_1\tau_1 + \alpha_2\tau_2$ (second column) for each set of initial conditions.

case	$\tilde{\lambda}^0$	$\tilde{\tau}^0$
1	0.0212146	1.57985e-07
2	0.0528147	6.31945e-08
3	0.0844624	3.94973e-08

Table 39: Error of the resulting values of $\tilde{\lambda}$ from the exact values for the underdetermined inverse problem using a frequency of $10^{11}Hz$ for each set of initial conditions.

case	α_1		
	.1	.5	.9
1	2.99851e-08	7.35704e-08	3.49983e-09
2	1.23759e-10	4.68118e-10	2.30235e-10
3	1.14292e-09	2.16657e-10	1.273e-10

Table 40: Error of the resulting values of $\tilde{\lambda}$ from the exact values for the underdetermined inverse problem using a frequency of 10^9Hz for each set of initial conditions.

case	α_1		
	.1	.5	.9
1	0.907928	0.711589	0.212452
2	4.93369e-06	1.05337e-05	1.47737e-06
3	3.86466e-06	4.68234e-06	6.40273e-08

Table 41: Error of the resulting values of $\tilde{\tau}$ from the exact values for the underdetermined inverse problem using a frequency of 10^6Hz for each set of initial conditions.

case	α_1		
	.1	.5	.9
1	5.78065e-12	1.12783e-11	5.70111e-12
2	7.42983e-12	2.03867e-11	4.54227e-13
3	2.3737e-11	2.04292e-11	3.88617e-12

5.5 Multiple Frequency Interrogation

In certain situations it may be feasible to add to the data used in the inverse problem in the previous section by interrogating the material again with a different frequency. In particular, if for example the first interrogating angular frequency is less than $\omega_c := 1/\tilde{\tau}$ then choosing the second frequency to be greater than this critical value should give different information, i.e., $\omega_L < 1/\tilde{\tau} < \omega_H$. We have already seen that for frequencies lower than the critical angular frequency we are able to accurately resolve $\tilde{\tau}$. Thus, using this information we may choose a second frequency to be higher than the critical level, and therefore expect to be able to resolve $\tilde{\lambda}$, as was done in our previous examples for high frequencies. We applied this approach to the test cases described for the $f = 10^6 Hz$ interrogating signal. Using the results of the previous section, i.e., fixing $\tilde{\tau}$ and using the corresponding values of τ_1 , τ_2 , and α_1 as initial conditions, we attempted an inverse problem to match data obtained by interrogating the same material with a frequency $f = 10^9 Hz$, which is higher than the critical value. In particular, we performed a one parameter search for each of τ_1 and τ_2 successively, then a two parameter search for both simultaneously for fine-tuning. As in previous inverse problems, this allowed determination of $\tilde{\lambda}$. Now with $\tilde{\lambda}$ and $\tilde{\tau}$ both fixed, we traverse along both contours simultaneously with a one parameter inverse problem for τ_1 . Lastly, again, we do a full three parameter inverse problem for fine-tuning. The results for Cases 1, 2, and 3 are given in Tables 42 - 44. Note that in order to ensure better accuracy when the wavelength is decreased, we double the number of finite elements used to solve the high frequency inverse problem.

In each case the lower frequency inverse problem was able to determine the value of $\tilde{\tau}$ to at least three decimal places, while the values of τ_1 , τ_2 , α_1 and $\tilde{\lambda}$ are in general not necessarily even improved over the initial estimates (e.g., τ_1 for Case 3). This is similar to the results of previous sections before traversing the contours of constant $\tilde{\tau}$ or $\tilde{\lambda}$. The higher frequency inverse problem with $\tilde{\tau}$ fixed was able to determine the value of $\tilde{\lambda}$ to at least three decimal places in each case, while again the values of τ_1 , τ_2 and α_1 are in general not necessarily even improved over the previous estimates. The final rows show the results from running a one parameter inverse problem for τ_1 with both $\tilde{\tau}$ and $\tilde{\lambda}$ fixed (and then using a three parameter search for fine tuning). Each of the final estimates for τ_1 , τ_2 , and α_1 have at least two decimal places accuracy.

Recall that, as explained in Remark 1, the slab thickness here is $400m$. These dimensions are merely arbitrary examples, but there is a direct restriction on the width of a slab if it is to be reliably simultaneously interrogated with frequencies just above and just below the critical angular frequency for a given medium, which is described by the material properties through the value of $\tilde{\tau}$. Thus, for example, materials with relaxation times on the order of 10^{-11} , like water, that have dimensions on the order of $.4m$ are quite feasible to interrogate with multiple frequencies above and below its critical frequency.

Table 42: Numerical results from the inverse problem using data from interrogating the same medium with one frequency above, and one frequency below the critical angular frequency for that medium. Given are the τ_1 , τ_2 , and α_1 values along with the relative error for $\tilde{\lambda}$ and $\tilde{\tau}$ for each step of the inverse problem: exact solution, initial conditions for Case 1, and the solutions from the low frequency inverse problem, from the high frequency inverse problem with fixed $\tilde{\tau}$, and finally from fixing both $\tilde{\tau}$ and $\tilde{\lambda}$.

Case 1	τ_1	τ_2	α_1	relerr($\tilde{\lambda}$)	relerr($\tilde{\tau}$)
Exact	3.16e-8	1.58e-8	0.1	0	0
Initial	1.58e-7	3.16e-9	0.9999	0.894221	8.09005
$f = 10^6 Hz$	2.12522e-8	8.82441e-9	0.687957	0.126493	3.33717e-4
$f = 10^9 Hz$	1.82157e-8	8.82198e-9	0.910418	1.24653e-4	3.33717e-4
Contour	3.13819e-8	1.57826e-8	0.102596	4.98614e-5	1.72612e-4

Table 43: Numerical results from the inverse problem using data from interrogating the same medium with one frequency above, and one frequency below the critical angular frequency for that medium. Given are the τ_1 , τ_2 , and α_1 values along with the relative error for $\tilde{\lambda}$ and $\tilde{\tau}$ for each step of the inverse problem: exact solution, initial conditions for Case 2, and the solutions from the low frequency inverse problem, from the high frequency inverse problem with fixed $\tilde{\tau}$, and finally from fixing both $\tilde{\tau}$ and $\tilde{\lambda}$.

Case 2	τ_1	τ_2	α_1	relerr($\tilde{\lambda}$)	relerr($\tilde{\tau}$)
Exact	3.16e-8	1.58e-8	0.1	0	0
Initial	6.32e-8	7.9e-9	0.9999	0.736658	2.63605
$f = 10^6 Hz$	1.91221e-8	8.40313e-9	0.836785	5.08337e-2	4.25777e-4
$f = 10^9 Hz$	1.81458e-8	8.47912e-9	0.920012	1.2964e-4	4.25777e-4
Contour	3.11405e-8	1.57683e-8	0.105044	5.48475e-5	1.78366e-4

Table 44: Numerical results from the inverse problem using data from interrogating the same medium with one frequency above, and one frequency below the critical angular frequency for that medium. Given are the τ_1 , τ_2 , and α_1 values along with the relative error for $\tilde{\lambda}$ and $\tilde{\tau}$ for each step of the inverse problem: exact solution, initial conditions for Case 3, and the solutions from the low frequency inverse problem, from the high frequency inverse problem with fixed $\tilde{\tau}$, and finally from fixing both $\tilde{\tau}$ and $\tilde{\lambda}$.

Case 3	τ_1	τ_2	α_1	relerr($\tilde{\lambda}$)	relerr($\tilde{\tau}$)
Exact	3.16e-8	1.58e-8	0.1	0	0
Initial	3.95e-8	1.264e-9	0.9999	0.578859	1.27257
$f = 10^6 Hz$	1.81785e-8	1.09823e-8	0.887927	1.79102e-2	4.60299e-4
$f = 10^9 Hz$	1.87871e-8	1.09078e-8	0.820402	1.09695e-4	4.60299e-4
Contour	3.12584e-8	1.57761e-8	0.103767	4.48752e-5	1.49597e-4

6 Inverse Problem Results Using a Uniform Distribution

The previous polarization inverse problems have assumed a discrete distribution with two atoms. According to experimental reports [23], most materials demonstrate polarization effects described by a range of relaxation times. Here we consider the simplest of distributions by exploring the possibility of a uniform distribution of relaxation times (τ) between a lower and upper limit (τ_a and τ_b respectively). This presents us again with a two parameter inverse problem, namely, to try to resolve the end points of the distribution used to generate the given data. Computationally we still approximate this distribution with discrete nodes, but instead of just one at each endpoint, we use $\ell = 13$ uniformly distributed within the interval (recall Section 3.2). Note that we do not wish to restrict our optimization routine to search only for $\tau_a < \tau_b$, therefore subsequently if $\tau_a > \tau_b$ it is assumed without loss of generality that τ_b is the lower limit of the distribution and τ_a is the upper limit.

Figures 14 and 15 depict the objective function and the log of the objective function, respectively, for a frequency of 10^{11} Hz . The solid line in Figure 15 is the curve of constant $\tilde{\lambda}$. Given a uniform distribution of relaxation times in a Debye medium, this parameter is given (analytically) by

$$\tilde{\lambda} := \frac{1}{c(\tau_b - \tau_a)} \int_{\tau_a}^{\tau_b} \frac{d\tau}{\tau} = \frac{\ln \tau_b - \ln \tau_a}{c(\tau_b - \tau_a)}. \quad (23)$$

Note that in computations we must use the same quadrature method to evaluate this integral as we do for \mathcal{G} to ensure a correct correlation. Although the curve appears slightly different from the discrete distribution case in Figure 3, the fact that this objective function is also small along this curve suggests that this problem should behave similarly to the discrete distribution case in Section 5. Figures 16 and 17 depict the objective function and the log of the objective function, respectively, for a frequency of 10^6 Hz . We notice that again the orientation of the “curve of best” fit is different from the higher frequency case. The solid line in Figure 17 is the curve of constant $\tilde{\tau} := \int_{\tau_a}^{\tau_b} \tau dF(\tau) = \frac{\tau_b - \tau_a}{2}$. Again, the fact that the $\tilde{\lambda}$ term only dominates the behavior when the interrogating frequency is greater than $\frac{1}{2\pi\tilde{\tau}}$ is consistent with the discrete distribution case.

Based on our previous experience with the discrete distribution, we anticipate that the two parameter inverse problem will simply converge to the “line of best fit”. Thus instead, we use an equivalent method to converge to this curve, namely minimizing over just one parameter, τ_a , leaving τ_b fixed at its initial value (the same values as those given as τ_2^0 in Table 1). It should be noted that one may actually perform two of the one parameter constrained optimizations, one in each of the directions τ_a and τ_b , to allow for the possibility of the first direction not converging, for example, if $\tau_a^0, \tau_b^0 \leq 10^{-9}$. In general, the second direction of the optimization requires only enough iterations to verify convergence, as the first direction has usually already converged to the “line of best fit”. Using this one parameter at a time approach, we still converge to the “line of best fit”, but theoretically use half as many function evaluations as the two parameter inverse problem, since we only compute one gradient at each step. (In practice, only a third as many function evaluations were actually needed to get the same order of accuracy as the two parameter inverse problem.)

We performed the one parameter inverse problem using the three initial condition cases described above in Table 1 and again using frequencies 10^{11} Hz , 10^9 Hz and 10^6 Hz . The τ_a estimates from running Levenberg-Marquardt on the modified least squares objective function are given in Table 45. Again, as in the discrete case, the values of the relaxation times do not appear to be converging to the correct solution. But we expect that the optimization routine is converging to the “curve of best fit”. To test this we must look at the approximations to $\tilde{\lambda}$ and $\tilde{\tau}$.

The initial values of $\tilde{\lambda}$ and $\tilde{\tau}$ are given in Table 46 (note that these values were computed by $\tilde{\lambda} \approx \sum_i \frac{\alpha_i}{c\tau_i}$ and $\tilde{\tau} \approx \sum_i \alpha_i \tau_i$ using appropriately defined $\{\alpha_i\}_{i=0}^{N-1}$ determined by the Composite Simpson’s rule). The exact values of each are $\tilde{\lambda}^* = 0.248369$ and $\tilde{\tau}^* = 5.42467 \times 10^{-8}$. The resulting $\tilde{\lambda}$ and $\tilde{\tau}$ values from running the one parameter Levenberg-Marquardt routine are given in Table 47. Clearly each case has converged to the “line of best fit”; in general the closer initial conditions converged closer to the actual value of $\tilde{\lambda}$ (or $\tilde{\tau}$ for $f = 10^6 \text{ Hz}$).

After our one parameter optimization routine resolved $\tilde{\lambda}$ (or $\tilde{\tau}$ for $f = 10^6 \text{ Hz}$), we minimized for τ_a along the line of constant $\tilde{\lambda}$ (or $\tilde{\tau}$). Again, this is a one parameter inverse problem, and therefore very efficient. The results of these computations are given in Tables 48 and 49, for τ_a and the corresponding τ_b (given constant $\tilde{\lambda}$ or $\tilde{\tau}$), respectively.

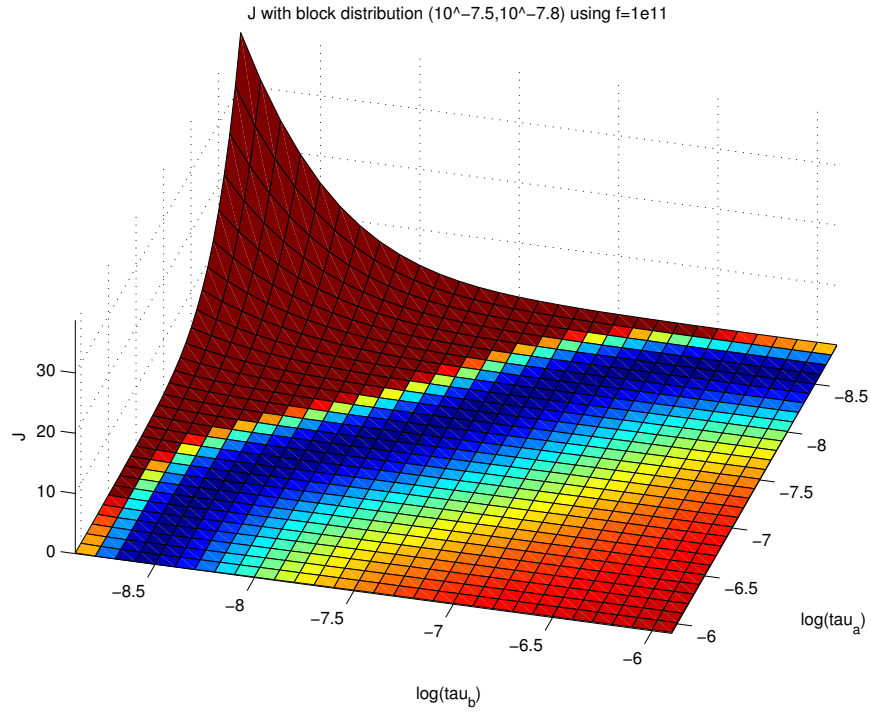


Figure 14: The objective function for the uniform distribution inverse problem versus the log of τ_a and the log of τ_b using a frequency of $10^{11} Hz$.

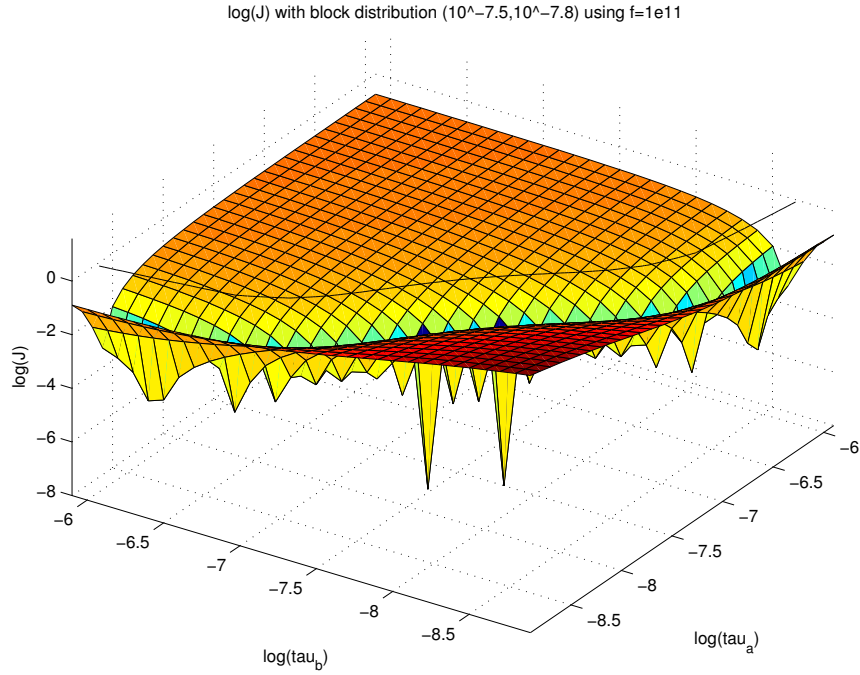


Figure 15: The log of the objective function for the uniform distribution inverse problem versus the log of τ_a and the log of τ_b using a frequency of $10^{11} Hz$. The solid line above the surface represents the curve of constant $\tilde{\lambda}$.

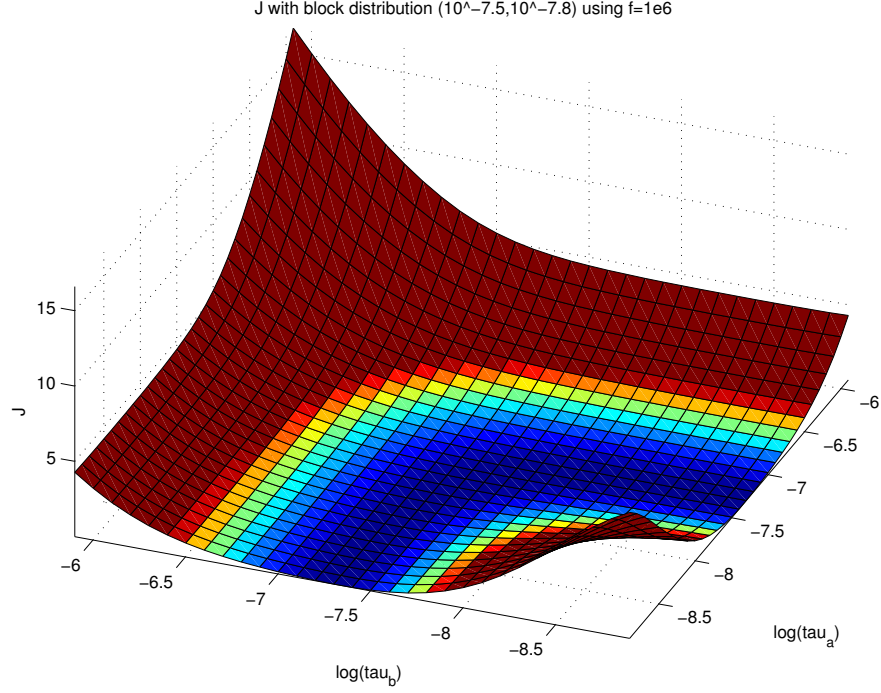


Figure 16: The objective function for the uniform distribution inverse problem versus the log of τ_a and the log of τ_b using a frequency of $10^6 Hz$.

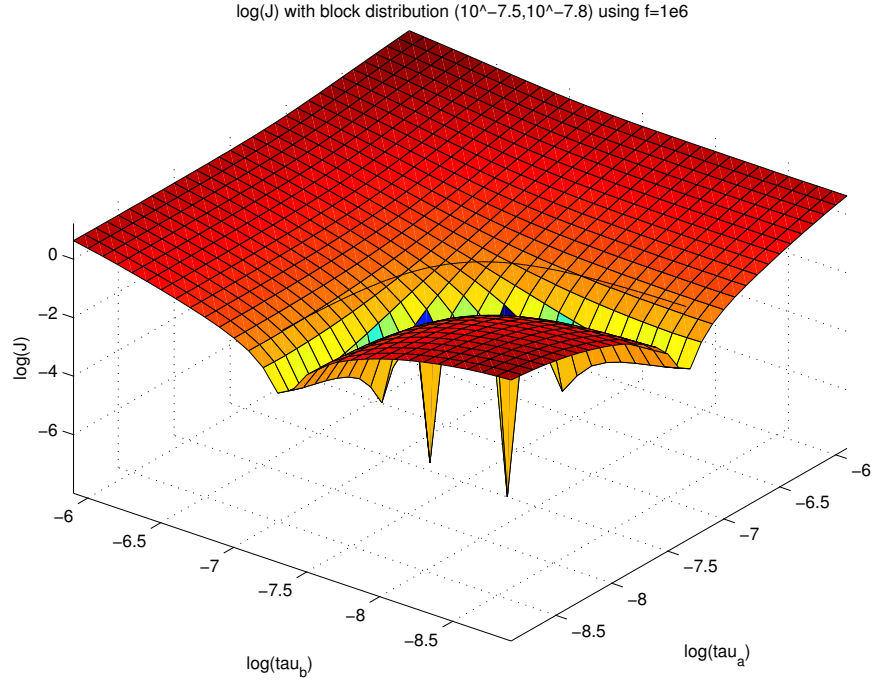


Figure 17: The log of the objective function for the uniform distribution inverse problem versus the log of τ_a and the log of τ_b using a frequency of $10^6 Hz$. The solid line above the surface represents the curve of constant $\tilde{\tau}$.

Table 45: Resulting values of τ_a from the one parameter Levenberg-Marquardt routine for the inverse problem to determine the endpoints of a uniform distribution of relaxation times (recall the exact solution $\tau_a^* = 3.16000\text{e-}8$).

case	Initial	Frequency (Hz)		
		10^{11}	10^9	10^6
1	1.58000e-7	5.41874e-8	5.40910e-8	3.59781e-8
2	6.32000e-8	4.01846e-8	4.01788e-8	3.43510e-8
3	3.95000e-8	3.41828e-8	3.41820e-8	3.27009e-8

Table 46: The initial values of $\tilde{\lambda} := \sum_i \frac{\alpha_i}{c\tau_i}$ (or $\tilde{\tau} := \sum_i \alpha_i \tau_i$ for $f = 10^6$) for each set of initial conditions for the inverse problem to determine the endpoints of a uniform distribution of relaxation times (case 0 represents the exact solution).

case	Frequency (Hz)	
	$10^{11} - 10^9$	10^6
0	0.248369	5.42467e-8
1	0.115739	2.33313e-7
2	0.176705	9.66433e-8
3	0.223136	6.42533e-8

Table 47: Resulting values of $\tilde{\lambda}$ (or $\tilde{\tau}$ for $f = 10^6$) from the Levenberg-Marquardt routine for the inverse problem to determine the endpoints of a uniform distribution of relaxation times for each set of initial conditions (case 0 represents the exact solution). The values in parenthesis denote the absolute value of the difference as the number of digits shown here would not sufficiently distinguish the approximations from the exact solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
0	0.248369	0.248369	5.42467e-8
1	(1.64144e-8)	0.248701	5.43479e-8
2	(1.51187e-9)	0.248396	5.43315e-8
3	(5.12895e-11)	0.248373	5.42813e-8

Table 48: Resulting values of τ_a from minimizing along the line of constant $\tilde{\lambda}$ (or $\tilde{\tau}$ for $f = 10^6 \text{ Hz}$), for the inverse problem to determine the endpoints of a uniform distribution of relaxation times (recall the exact solution $\tau_a^* = 3.16000\text{e-}8$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	3.08694e-08	3.1653e-08	3.17617e-08
2	3.16401e-08	3.16043e-08	3.17357e-08
3	3.15905e-08	3.16007e-08	3.16559e-08

Finally, for fine tuning, we apply the full two parameter Levenberg-Marquardt routine using the estimates from minimizing along constant $\tilde{\lambda}$ (or $\tilde{\tau}$). These results are shown in Tables 50 and 51. We see that the estimates change very little, if at all, which suggests that our approximation method is not only efficient, but quite accurate as well.

Table 49: Resulting values of τ_b from the minimizing along the line of constant $\tilde{\lambda}$ (or $\tilde{\tau}$ for $f = 10^6 Hz$), for the inverse problem to determine the endpoints of a uniform distribution of relaxation times (recall the exact solution $\tau_b^* = 1.5800e-8$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	1.68829e-08	1.56448e-08	1.55282e-08
2	1.57433e-08	1.57874e-08	1.55715e-08
3	1.58134e-08	1.5798e-08	1.57053e-08

Table 50: Resulting values of τ_a from the two parameter Levenberg-Marquardt routine for the inverse problem to determine the endpoints of a uniform distribution of relaxation times (recall the exact solution $\tau_a^* = 3.16000e-8$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	3.08694e-08	3.15996e-08	3.16014e-08
2	3.16401e-08	3.16005e-08	3.16008e-08
3	3.15905e-08	3.16008e-08	3.16002e-08

Table 51: Resulting values of τ_b from the two parameter Levenberg-Marquardt routine for the inverse problem to determine the endpoints of a uniform distribution of relaxation times (recall the exact solution $\tau_b^* = 1.58000e-8$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	1.68829e-08	1.58005e-08	1.57988e-08
2	1.57433e-08	1.57993e-08	1.57991e-08
3	1.58134e-08	1.57988e-08	1.57998e-08

7 Inverse Problem Results Using a Gaussian Distribution

In this section we explore the possibility of determining dielectric parameters that are normally distributed. We apply a Gaussian distribution to the logarithms of the relaxation times (thus actually the relaxation times are log-normally distributed, but all of our computations are in “log”-space so we refer to these parameters as normally distributed). Based on experimental data in [23], this seems to be a likely model for the relaxation times of a material.

Recall from Section 3.3 that our PDF is

$$dF(\tau; \mu, \sigma) = \frac{1}{\sqrt{2\pi\sigma^2}} \frac{1}{\ln 10} \frac{1}{\tau} \exp\left(-\frac{(\log \tau - \mu)^2}{2\sigma^2}\right) d\tau \quad (24)$$

for a log-normal distribution in τ . Here μ is the mean of the relaxation times and σ is the standard deviation. Recall further that the density is truncated to a support interval $[\tau_a, \tau_b]$ where τ_a and τ_b are determined based on μ and σ as follows

$$\begin{aligned} \tau_a &:= 10^{\mu-6\sigma} \\ \tau_b &:= 10^{\mu+6\sigma}. \end{aligned}$$

Thus, this problem again presents us with a two parameter inverse problem, namely, to try to resolve the mean and standard deviation of the distribution used to generate the given data. Learning from the success in the previous section, we expect to apply a three step approach: one parameter Levenberg-Marquardt for

the mean μ (which we expect to converge to the line of constant $\tilde{\lambda}$ or $\tilde{\tau}$), then traverse the line of constant $\tilde{\lambda}$ or $\tilde{\tau}$ (which is again one parameter), and finally, perform two parameter Levenberg-Marquardt for fine-tuning. We can see from the surface plot using $f = 10^{11}$ shown in Figure 18 that we do again have a trench representing a “line of best fit”, which does in fact have a single minimizer, but it appears to be independent of σ . The trench is actually slightly slanted, as can be seen in Figure 19 which zooms in on μ , and in fact, for the lower frequencies the slant is indeed in the opposite direction. Therefore while at first glance of Figure 18 traversing a line of constant $\tilde{\lambda}$ or $\tilde{\tau}$ may seem unnecessary, it does actually give better results than simply minimizing in the σ direction alone.

Remark 2 *Note that the slant does not appear to be a consequence of our quadrature rule. Changing the number of nodes used in our computations had no effect on the slope. Further, increasing the support on which the integral was computed had no effect, e.g., to 12σ on either side of μ . While $\int \mathcal{G}(0, z) = \mu$ is actually independent of σ we still expect σ to affect the objective function in some other way, and this is in fact what is happening. Lastly, the apparent local minima in Figure 19 are again just an artifact of the grid points not falling exactly on the “line of best fit”.*

We performed our inverse problem solution approach to the sample problem of $\mu = 10^{-7.62525}$ and $\sigma = 0.0457575$ (this choice of parameters results in a distribution function comparable to the uniform distribution case of the previous section in that the corresponding densities have the same mean and roughly similar support, see Figure 20). The cases of initial values we considered for μ and σ , and the corresponding initial $\tilde{\lambda}$ and $\tilde{\tau}$ values, are given in Table 52. Note that Case 3 is specifically designed to test whether the standard deviation can be determined when the mean is known.

Table 52: Initial estimates and corresponding $\tilde{\lambda}$ and $\tilde{\tau}$ values for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times (case 0 corresponds to the exact solution).

case	$\log(\mu_0)$	σ_0	$\tilde{\lambda}_0$	$\tilde{\tau}_0$
0	-7.62525	0.0457575	0.141524	2.38319e-08
1	-6.92628	0.036606	0.0282483	1.18922e-07
2	-8.32422	0.054909	0.709351	4.77804e-09
3	-7.62525	0.0411817	0.141375	2.37329e-08

The μ estimates resulting from the one parameter Levenberg-Marquardt are given in Table 53. In each case at least three decimals places of accuracy was achieved. Also, at least seven decimal places of accuracy are obtained on the estimates of the corresponding $\tilde{\lambda}$ and $\tilde{\tau}$ values, as displayed in Table 54. The μ and σ estimates resulting from the one parameter tracing of the “line of best fit” are given in Tables 55 and 56, respectively. In each case at least six decimals places of accuracy was achieved for the log of μ . Also, at least three decimal places of accuracy are obtained on the estimates of σ . Finally, the μ and σ estimates resulting from the two parameter Levenberg-Marquardt fine-tuning are given in Tables 57 and 58, respectively. In all but a couple cases we have as much degree accuracy as the number of digits shown here allows. Figure 21 gives a graphical representation of how well our inverse problem solution method converged given the initial condition Case 1 from Table 52. We conclude that the inverse problem involving a Gaussian distribution of relaxation times for a Debye polarization model is computationally feasible for the sample parameters presented here.

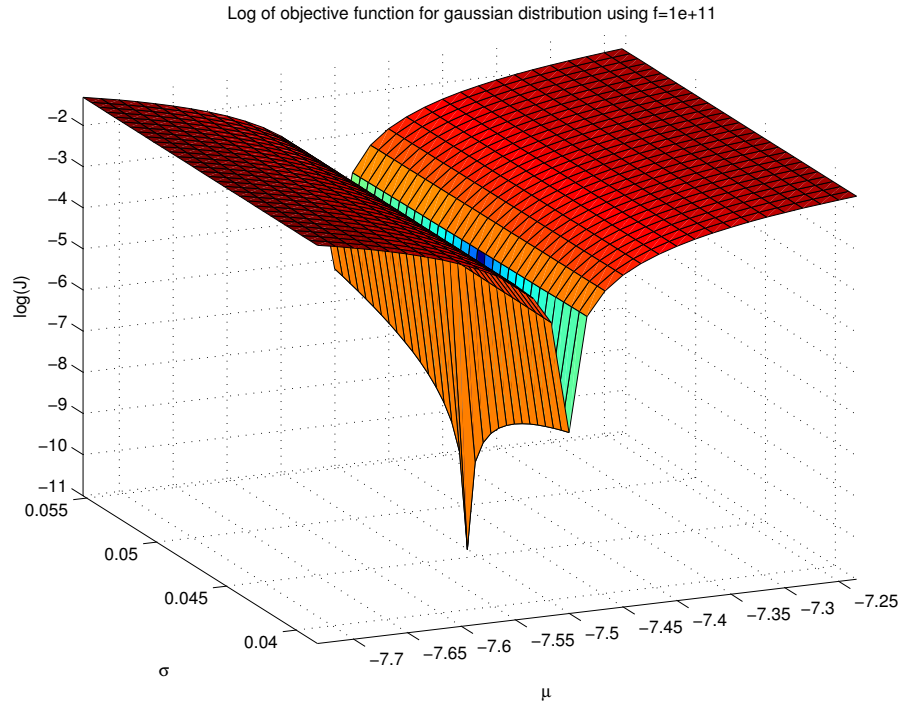


Figure 18: The log of the objective function for the gaussian distribution inverse problem versus the log of μ and σ using a frequency of $10^{11} Hz$.

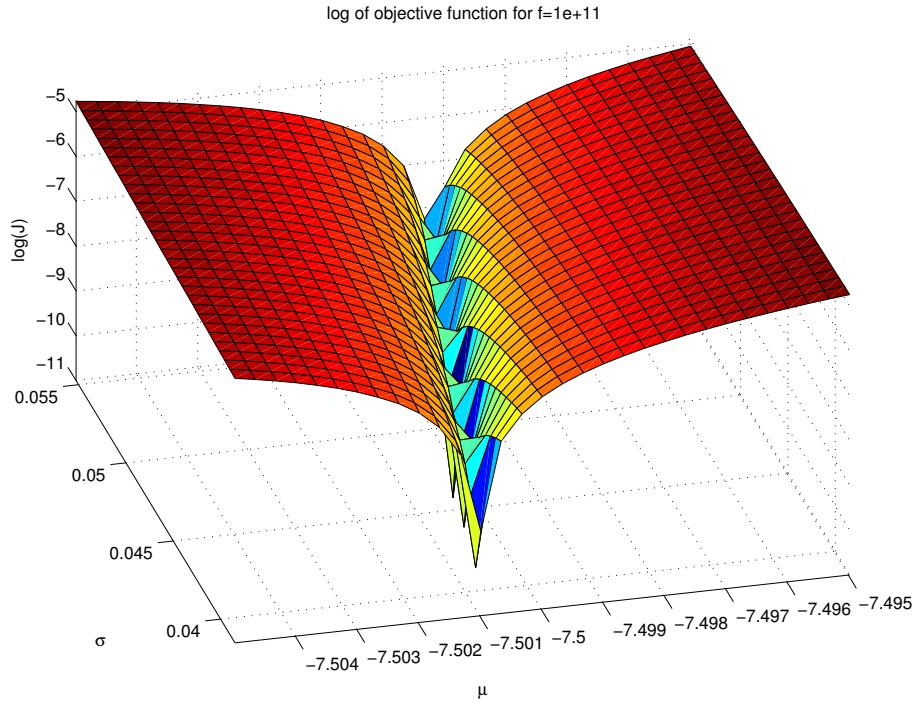


Figure 19: The log of the objective function for the gaussian distribution inverse problem versus the log of μ and σ , zoomed in with respect to μ , and using a frequency of $10^{11} Hz$. The appearance of several local minima is merely an artifact of the lattice points used to plot the surface.

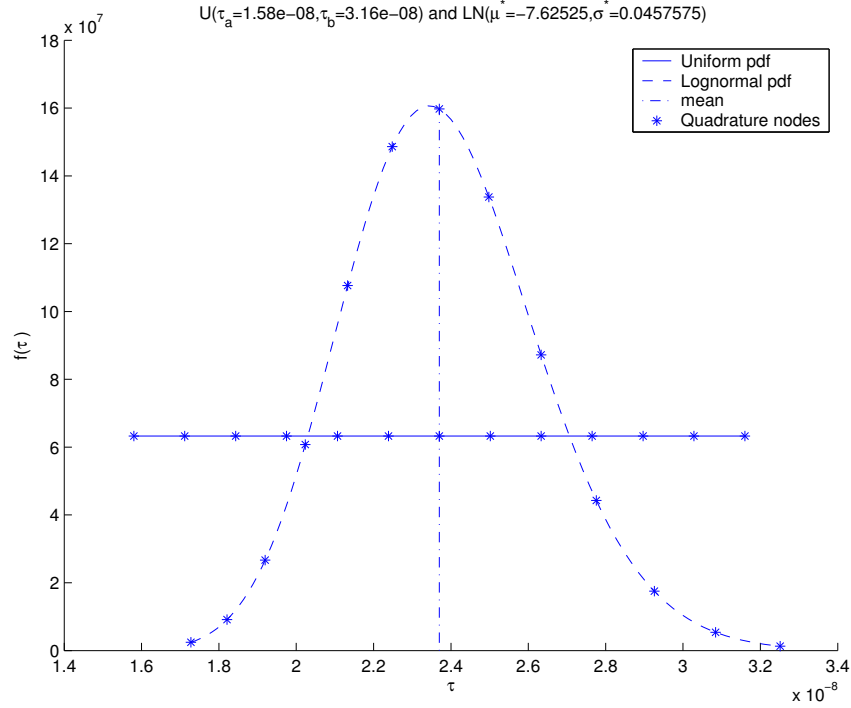


Figure 20: This plot compares the uniform density considered in Section 6 to the lognormal density. Shown are the two density functions and the corresponding quadrature nodes used in integrating them.

Table 53: The resulting $\log(\mu)$ values from performing a one parameter Levenberg-Marquardt optimization for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu^*) = -7.62525$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.62612	-7.6261	-7.6244
2	-7.62419	-7.62422	-7.62629
3	-7.62571	-7.62701	-7.62348

Table 54: The corresponding $\tilde{\lambda}$ and $\tilde{\tau}$ values from performing a one parameter Levenberg-Marquardt optimization for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case. Note that the absolute values of the differences between the estimates and the exact values are shown as the number of digits shown here would not sufficiently distinguish the approximations from the exact solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	2.62469e-11	7.38776e-08	8.98789e-13
2	3.61986e-11	9.10257e-08	1.10726e-12
3	4.25993e-13	1.56168e-07	1.86751e-12

Table 55: The resulting $\log(\mu)$ values from traversing the “line of best fit” for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu^*) = -7.62525$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.62526	-7.62525	-7.62526
2	-7.62524	-7.62525	-7.62524
3	-7.62524	-7.62525	-7.62527

Table 56: The resulting σ values from traversing the “line of best fit” for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\sigma^* = 0.0457575$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	0.0456376	0.0457564	0.0456871
2	0.0458453	0.0457592	0.0458442
3	0.0458207	0.0457551	0.0456102

Table 57: The resulting $\log(\mu)$ values from performing a two parameter Levenberg-Marquardt optimization for fine-tuning of the solution for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu^*) = -7.62525$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.62526	-7.62525	-7.62525
2	-7.62524	-7.62525	-7.62525
3	-7.62524	-7.62525	-7.62525

Table 58: The resulting σ values from performing a two parameter Levenberg-Marquardt optimization for fine-tuning of the solution for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\sigma^* = 0.0457575$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	0.0456376	0.0457576	0.0457575
2	0.0458453	0.0457576	0.0457575
3	0.0458207	0.0457575	0.0457575

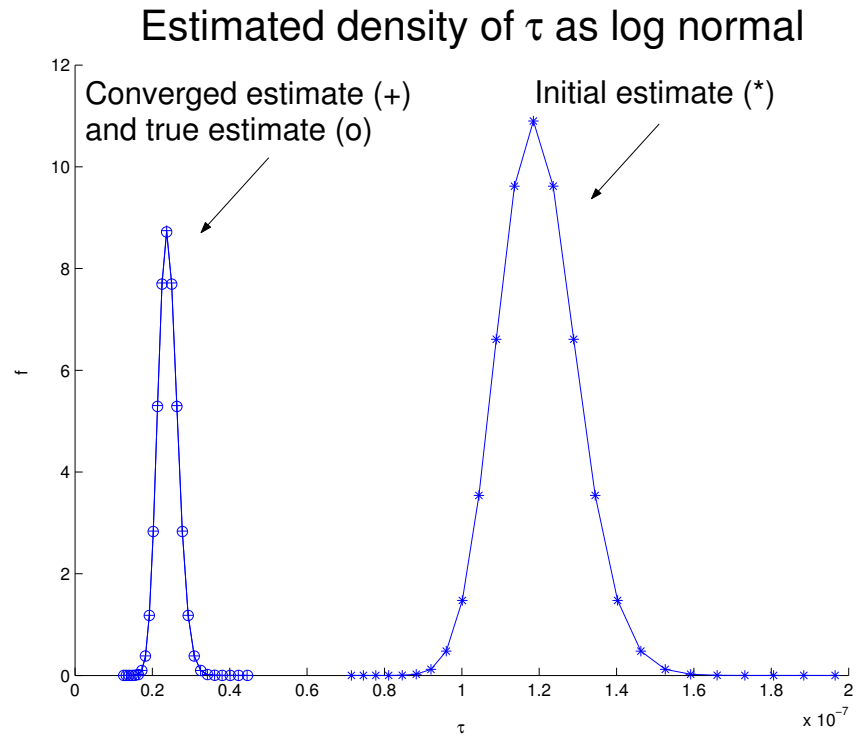


Figure 21: Shown are the initial density function, the minimizing density function and the true density function (the latter two being practically identical) corresponding to Case 1 in Table 52.

8 Inverse Problem Results Using a Bi-Gaussian Distribution

In this section we combine the two ideas from Sections 5 and 7 in that we consider a material comprised of primarily two distinct polarization mechanisms, but rather than assuming two atoms at τ_1 and τ_2 as before we let $\mu_1 = \tau_1$ and $\mu_2 = \tau_2$ in two Gaussian distributions. Essentially, we are decomposing the distribution function into two components, each dependent on distinct means and standard deviations as follows:

$$d\hat{F} = \alpha dF(\tau; \mu_1, \sigma_1) + (1 - \alpha) dF(\tau; \mu_2, \sigma_2), \quad (25)$$

where dF is given by (24) (we take $\alpha = .5$ in our numerical experiments).

We now have a four parameter inverse problem, namely

$$\min_{(\mu_1, \sigma_1, \mu_2, \sigma_2) \in Q} J(\hat{F})$$

where Q is the admissible region for our unknown parameters (e.g., $\sigma_i > 0$). For the following examples, the true values of the means were taken to be $(\mu_1^*, \mu_2^*) = (10^{-7.80134}, 10^{-7.50031})$, with the standard deviations ranging from .02 to .07. As in the previous gaussian case, we expect the objective function to be relatively less dependent on the standard deviations than the means, therefore we address the dependency on the means first. Figure 22 shows a surface plot of the objective function with respect to μ_1 and μ_2 . Here we can see the similarities to the distribution function from Section 5, namely the presense of the “line of best fit” with two (symmetric) global minima. (When comparing to Figure 3, note that here we have rotated the axis for better viewing of the minima under the surface.) Our minimization approach is thus the same as that which worked well in the previous sections. We perform a one parameter Levenberg-Marquardt search in the μ_1 direction (and just to be sure, a one parameter Levenberg-Marquardt search in the μ_2 direction), then optimize along the “line of best fit”, and finally fine-tune with a two parameter Levenberg-Marquardt search for both μ_1 and μ_2 .

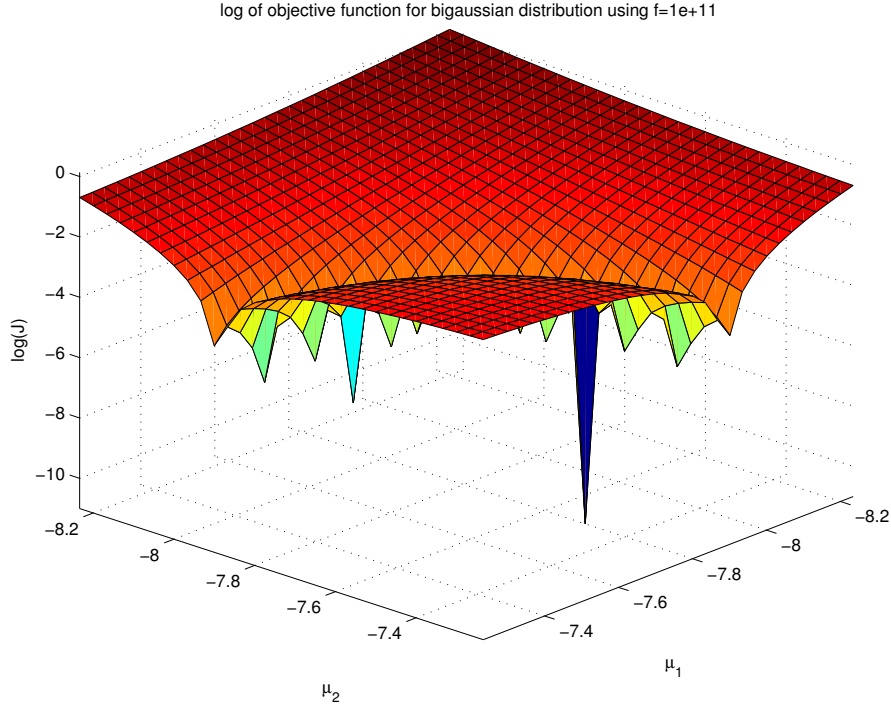


Figure 22: The log of the objective function for the bi-gaussian distribution inverse problem versus the log of μ_1 and μ_2 using a frequency of $10^{11} Hz$.

Thus far we have completely ignored the standard deviations. Figures 23 and 24 show the objective function versus μ_1 and μ_2 for two distinct cases: using the exact values for the standard deviations and using significantly incorrect ones. The fact that the location of the “line of best fit” does not change drastically suggests that we should be able to at least determine the values of $\tilde{\lambda}$ or $\tilde{\tau}$ with just our one parameter searches, as in the previous sections. However, the location of the local minima along this curve has changed, therefore we should not expect our μ_1 and μ_2 estimates so far to be our final answer. We must first attempt to optimize with respect to the standard deviations, and then finally attempt a full four parameter minimization to make sure we have accounted for all interdependencies. These are the last two steps of our now six step optimization process for the four parameter inverse problem involving a Bi-Gaussian distribution.

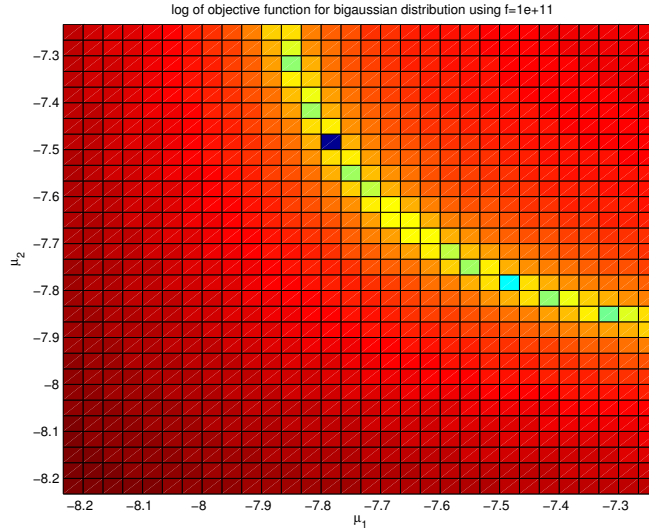


Figure 23: The log of the objective function for the bi-gaussian distribution inverse problem versus the log of μ_1 and μ_2 using correct values for the standard deviations.

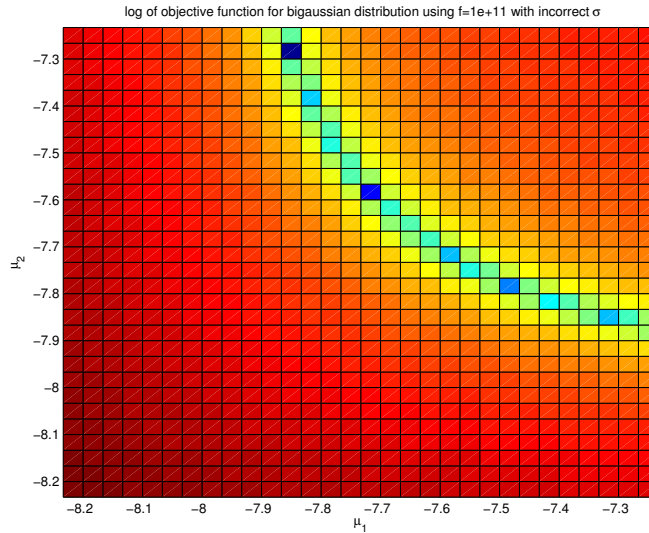


Figure 24: The log of the objective function for the bi-gaussian distribution inverse problem versus the log of μ_1 and μ_2 using incorrect values for the standard deviations.

The true values of the optimization parameters and the initial conditions considered for each case we attempted are given in Table 59. These values correspond to roughly the same relative initial error as the initial value cases in the single gaussian problem of Section 7. The μ_1 and μ_2 estimates resulting from the one parameter Levenberg-Marquardt searches are given in Tables 60 and 61, respectively.

Remark 3 *Note that in Case 2 at $f = 10^{11}$ the value of μ_1 diverged toward zero in our first attempts. Therefore the results shown in the tables for Case 2 at each frequency represent first minimizing in the μ_2 direction followed by minimization in the μ_1 direction. In general better results are achieved when optimizing the smaller of the values (i.e., corresponding to a more negative logarithm) first.*

Several of the μ_2 estimates did not change from their initial values during the one parameter optimization for μ_2 suggesting that the μ_1 search arrived sufficiently close to the $\tilde{\lambda}$ or $\tilde{\tau}$ curve (likewise for Case 2, the two higher frequency estimates for μ_1 did not change suggesting that the μ_2 search converged sufficiently). The values of $\tilde{\lambda}$ or $\tilde{\tau}$ from the two one-parameter searches are given in Table 62. These values, on average, have about two fewer decimal places of accuracy as those in the single gaussian case. Still, as the following results shall show, in general this is sufficiently close to allow for good final estimates.

We now hold $\tilde{\lambda}$ or $\tilde{\tau}$ fixed and search for μ_1 and μ_2 . The μ_1 and μ_2 estimates resulting from the one parameter tracing of the “line of best fit” are given in Tables 63 and 64, respectively. Finally, the μ_1 and μ_2 estimates resulting from the two parameter Levenberg-Marquardt fine-tuning are given in Tables 65 and 66, respectively. The lack of a significant difference after the two parameter search suggests that we have found local minima with respect to μ_1 and μ_2 given our initial standard deviation estimates.

Remark 4 *Note that in some cases (actually, about half!) μ_1 is converging toward μ_2^* (we have highlighted these cases in each subsequent table). This is the same problem that we encountered in the discrete case in Section 5. For our examples here, $\alpha = .5$ so the problem is symmetric provided the standard deviations converge to the corresponding symmetric solution as well (i.e., σ_1^0 converges to σ_2^*). We could restrict our optimization problem to, for instance, only allow $\mu_1 > \mu_2$, but we prefer the idea that there are two equally viable solutions to choose from since it may actually double our chances of finding one of them! If our initial estimates for each σ_i were based on a certain ordering of the μ_i ’s (for example, say we know the distribution of the smaller relaxation time is much more narrow) we could easily test for that ordering and rearrange our σ_i estimates if necessary (though we would need to optimize with respect to each μ_i once more to reflect the changes).*

While the results for determining the means of the unknown Bi-Gaussian distribution are surprisingly decent, the optimization of the standard deviations was expectedly more difficult. Figure 25 demonstrates that the objective function behaves in a similar way with respect to the standard deviations as it did to the means in that there is a “line of best fit”. Unfortunately, when compared to Figure 22 it becomes obvious that the scale involved in traversing along this curve is too small to be accurate enough for optimization. This is in fact what we experienced in attempting to minimize with respect to σ_1 along the “line of best fit”. The changes in each σ_1 were in almost all cases less than six decimal places. Applying a two parameter Levenberg-Marquardt for fine-tuning afterwards gave slightly better results for some cases, and we display these approximations in Tables 67 and 68. Using these values along with the best estimates so far for our means, we run a four parameter Levenberg-Marquardt search. The results for this procedure are given in Tables 69 through 72.

One may argue that the lack of sensitivity to σ may be due to starting initial estimates too close to the true solution, but this would not explain why, with all other parameters held fixed, gradient based methods take the standard deviation estimates farther from the true solution in half of the cases we tried. It is more likely that the numerical errors involved in the simulations are affecting the gradients. However, the results here reflect attempts that have been made to account for this using varying gradient step sizes (our Levenberg-Marquardt routine exits only after verifying small gradients with three different step sizes as explained in [13]).

Our conclusion for the problem of determining the means and the standard deviations for a bi-gaussian distribution of relaxation times is that the means are very readily determined with reasonable initial estimates, even with significantly incorrect standard deviations. However, the standard deviations are not easily determined even with quite accurate estimates for the means.

Table 59: Initial estimates for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times. For each case, $(\mu_1^*, \mu_2^*) = (10^{-7.80134}, 10^{-7.50031})$. For Case 1, $(\sigma_1^*, \sigma_2^*) = (0.0457575, 0.0457575)$. For Cases 2 and 3, $(\sigma_1^*, \sigma_2^*) = (0.0705811, 0.0222764)$.

case	$\log(\mu_1^0)$	σ_1^0	$\log(\mu_2^0)$	σ_2^0
1	-8.50031	0.036606	-6.80134	0.0571969
2	-7.10237	0.0846973	-8.19928	0.0185637
3	-8.50031	0.0352905	-6.80134	0.0445528

Table 60: The resulting $\log(\mu_1)$ values from performing a one parameter Levenberg-Marquardt optimization in μ_1 for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_1^*) = -7.80134$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.94848	-7.94847	-8.37177
2	-7.10237	-7.10237	-7.39253
3	-7.95035	-7.95034	-8.37264

Table 61: The resulting $\log(\mu_2)$ values from performing a one parameter Levenberg-Marquardt optimization in μ_2 for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_2^*) = -7.50031$).

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-6.80134	-6.80377	-7.36944
2	-7.91764	-7.91763	-8.13781
3	-6.80134	-6.80375	-7.36803

Table 62: The corresponding $\tilde{\lambda}$ and $\tilde{\tau}$ values from performing each one parameter Levenberg-Marquardt optimization in μ_1 and μ_2 for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case. Note that the absolute values of the differences between the estimates and the exact values are shown.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	7.945e-10	5.7109e-05	1.45861e-10
2	5.00444e-10	1.73813e-06	4.24458e-10
3	1.74996e-08	5.57019e-05	1.46209e-10

Table 63: The resulting $\log(\mu_1)$ values from traversing the “line of best fit” for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_1^*) = -7.80134$). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.80682	-7.5052	-7.79549
2	-7.47853	-7.47881	-7.50275
3	-7.81476	-7.48674	-7.79517

Table 64: The resulting $\log(\mu_2)$ values from traversing the “line of best fit” for the inverse problem to determine the mean and standard deviation of a normal distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_2^*) = -7.50031$). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.48957	-7.79775	-7.50821
2	-7.81299	-7.81285	-7.78246
3	-7.47935	-7.8106	-7.50706

Table 65: The resulting $\log(\mu_1)$ values from performing a two parameter Levenberg-Marquardt optimization for fine-tuning of (μ_1, μ_2) for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_1^*) = -7.80134$). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.80682	-7.5052	-7.79549
2	-7.47851	-7.47881	-7.50275
3	-7.81442	-7.48674	-7.79517

Table 66: The resulting $\log(\mu_2)$ values from performing a two parameter Levenberg-Marquardt optimization for fine-tuning of (μ_1, μ_2) for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_2^*) = -7.50031$). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.48957	-7.79775	-7.50821
2	-7.813	-7.81285	-7.78246
3	-7.47962	-7.8106	-7.50706

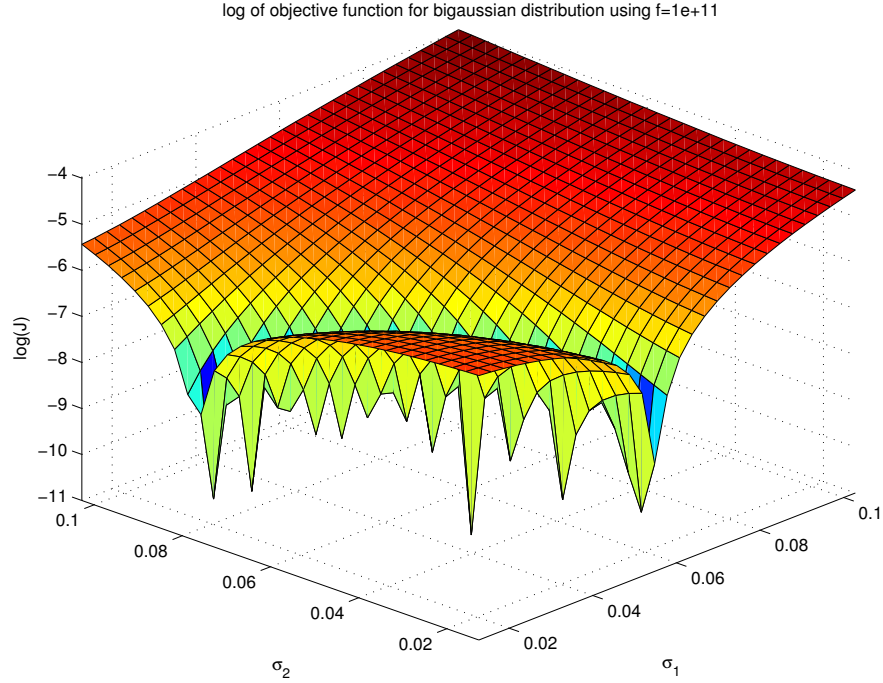


Figure 25: The log of the objective function for the bi-gaussian distribution inverse problem versus σ_1 and σ_2 using a frequency of $10^{11} Hz$.

Table 67: The resulting σ_1 values from a two parameter Levenberg-Marquardt search for σ_1 and σ_2 for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\sigma_1^* = 0.0457575$ for Case 1 and 0.0705811 otherwise). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	0.036606	0.00603311	0.0136791
2	0.0846974	0.0844293	0.0761876
3	0.0352905	0.00857058	0.0374784

Table 68: The resulting σ_2 values from a two parameter Levenberg-Marquardt search for σ_1 and σ_2 for the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\sigma_2^* = 0.0457575$ for Case 1 and 0.0222764 otherwise). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	0.0571969	0.060935	0.0658016
2	0.0185637	0.0192023	0.044288
3	0.0445528	0.0482781	0.0456734

Table 69: The resulting $\log(\mu_1)$ values from performing a four parameter Levenberg-Marquardt optimization for all parameters in the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_1^*) = -7.80134$). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.80682	-7.5052	-7.78543
2	-7.47851	-7.47911	-7.50072
3	-7.81432	-7.48831	-7.78507

Table 70: The resulting $\log(\mu_2)$ values from performing a four parameter Levenberg-Marquardt optimization for all parameters in the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\log(\mu_2^*) = -7.50031$). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	-7.48957	-7.79687	-7.50993
2	-7.813	-7.81279	-7.80691
3	-7.47506	-7.80992	-7.50848

Table 71: The resulting σ_1 values from a four parameter Levenberg-Marquardt search for all parameters in the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\sigma_1^* = 0.0457575$ for Case 1 and 0.0705811 otherwise). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	0.0366059	0.0056312	0.0136811
2	0.0846973	0.0832475	0.0736677
3	0.0352905	0.00879725	0.037513

Table 72: The resulting σ_2 values from a four parameter Levenberg-Marquardt search for all parameters in the inverse problem to determine the means and standard deviations of a bi-gaussian distribution of relaxation times for each frequency and initial value case (recall the exact solution is $\sigma_2^* = 0.0457575$ for Case 1 and 0.0222764 otherwise). The highlighted values denote convergence to the symmetric solution.

case	Frequency (Hz)		
	10^{11}	10^9	10^6
1	0.0571968	0.0604793	0.0663628
2	0.0185637	0.0198981	0.0442772
3	0.0445528	0.0477348	0.0463262

9 Summary

In this report we have discussed theoretical and numerical results for inverse problems involving Maxwell's equations with a general polarization term which includes probability distributions for both dielectric parameters and mechanisms. A theoretical framework to treat such classes of problems was given along with stability results for the corresponding least squares inverse problems. We presented examples of distributions of parameters including discrete, uniform, log-normal, and log-bi-gaussian. In each case the objective function was characterized by a "line of best fit" corresponding to a curve of constant value of either λ or τ depending on the relative values of the interrogating frequency ω , and $1/\tilde{\tau}$, where $\tilde{\tau}$ is the weighted average of the relaxation times. The constant values of these curves turned out to be, in fact, $\tilde{\tau}$ and $\tilde{\lambda}$, which is the weighted average of the inverses of the relaxation times (scaled by the speed of light for convenience). The similarities between the cases were exploited in determining effective optimization procedures for subsequent cases.

For the discrete case involving two atoms, each value of τ was readily determined if the volume proportions were known. Likewise, the volume proportion could be determined given the values of each relaxation time. However, the problem of determining all three parameters simultaneously is under-determined. Although, in certain situations, interrogating with a different frequency (in particular, greater than the critical frequency if the original interrogating frequency is smaller) can provide enough new information for the volume proportions and the relaxation times to be determined *simultaneously*.

The uniform distribution is generally straight-forward to optimize; each endpoint of the distribution was determined, on average, to about three decimal places of accuracy.

The inverse problem involving the log-normal distribution, or rather a gaussian distribution on the logarithms of the relaxation times, behaved similarly to the discrete distribution, and therefore the previous solution methods worked well. The mean of the distribution was determined to about four decimal places in log-space. The standard deviation was also determined quite accurately (about three decimal places), but it was evident in the surface plots that the objective function was particularly insensitive to the standard deviation. Any amount of noise in the system may prevent accurate results.

Finally, we considered a log-bi-gaussian distribution of relaxation times. As in the gaussian distribution problem, the objective function was insensitive to the standard deviations. Therefore, even with quite accurate estimates for the means, no usable estimates for the standard deviations could be found. However, in contrast, rather accurate solutions for both of the means could be determined using substantially incorrect estimates for the standard deviations. Still, the effect of the standard deviation on the objective function is sufficient to be able to distinguish distributions which are fairly broad (large standard deviations) from those that are nearly delta functions (very small standard deviations). This is important as it suggests that systems which truly have continuous distributions for their parameters should not be modeled using discrete distributions. Further, an objective function of the type in this report may be used to create a measure of the error induced in the signal by opting to use a discrete distribution instead of a continuous one when it is sufficiently helpful for computational simplicity to do so.

In fact, homogenization techniques determine single values of relaxation times which produce signals that closely match those generated with more complex mechanisms (such as a distribution of Debye polarization mechanisms)—for example, see [7]. It turns out that from our analysis here that for lower frequencies ($\omega_L < \omega_c = 1/\tilde{\tau}$) the *effective* relaxation time is simply the weighted average of the actual relaxation times (or for continuous distributions, the integral of τ with respect to the distribution). However, for higher frequencies ($\omega_H > \omega_c = 1/\tilde{\tau}$) the *effective* relaxation time is actually the inverse of the weighted average of inverses (or the inverse of the integral of $1/\tau$ with respect to the distribution in the continuous case). For example,

$$\tau_e = \frac{1}{\sum_i \frac{\alpha_i}{\tau_i}},$$

where α_i represents the volume fraction of the material with discrete relaxation times τ_i . Both the low frequency and high frequency estimates are common results in homogenization theory. However, any homogenization procedure which does not take the interrogating frequency into account will not be able to determine which is the correct *effective* relaxation time for the situation.

Lastly, it should be noted that the procedure of traversing a "line of best fit", as done in this report, is only a valid approach to solving non-linear optimization problems (which have relations that are relatively

optimal) if the equation of this curve is known. Other possible solutions methods for more general problems include taking an “arbitrary orthogonal step”, or resorting to a simplex search for those parts of the inverse problem that are particularly difficult for the gradient based methods.

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