Example: Seismic Traveltime Tomography

Various examples of the methods to be discussed in this talk could be provided, but to keep the analysis both simple and concrete we will limit most of our discussion to physical problems based on seismic traveltime tomography. Other choices such ocean acoustic tomography or electrical impedance tomography, among many others, could have been made, but an understanding of seismic traveltime tomography is sufficient to introduce all the main ideas and yet complex enough to be a good introduction to the main difficulties with these methods as well. We will introduce the problem of seismic traveltime tomography now.

A typical problem in seismic traveltime tomography is to infer the (isotropic) compressional-wave velocity distribution of an inhomogeneous medium, given a set of observed first-arrival traveltimes between sources and receivers of known location within the medium. This problem is common for borehole-to-borehole seismic tomography in oil field applications. We might also consider the problem of inverting for wave slowness when the absolute traveltimes are not known, as is normally the case in earthquake seismology but will not discuss this further today.

Slowness Models

We can consider three kinds of slowness (reciprocal of velocity) models. If we permit the slowness to be a general function of position, \(s(\mathbf{x})\), then this is a continuum model, which are explicitly excluding from consideration here. We can instead make one of two more restrictive assumptions that (i) the model comprises homogeneous blocks, or cells, with \(s_j\) denoting the slowness value of the \(j\)th cell, or (ii) the model is composed of a grid with values of slowness assigned at the grid points with some interpolation scheme.
to assign the values between grid points.

When it is not important which type of slowness model is involved, we will refer to
the model abstractly as a vector \( s \) in a vector space. For a block model with \( n \) blocks,
the model becomes a vector in \( \mathbb{R}^n \), the \( n \)-dimensional Euclidean vector space.

**Fermat’s Principle and Traveltime Functionals**

The traveltime of a seismic wave is the integral of slowness along a ray path connecting
the source and receiver. To make this more precise, we will define two functionals for
taveltime.

Let \( P \) denote an arbitrary path connecting a given source and receiver in a slowness
model \( s \). We will refer to \( P \) as a *trial ray path*. We define a functional \( \tau^P \) which yields
the traveltime along path \( P \). Letting \( s \) be the continuous slowness distribution \( s(x) \), we have

\[
\tau^P(s) = \int_P s(x) \, dl^P,
\]

where \( dl^P \) denotes the infinitesimal distance along the path \( P \).

**Fermat’s principle** states that the correct ray path between two points is the one
of least overall traveltime, *i.e.*, it minimizes \( \tau^P(s) \) with respect to path \( P \). [Actually,
Fermat’s principle is the weaker condition that the traveltime integral is *stationary* with
respect to variations in the ray path, but for traveltime tomography using measured first
arrivals it follows that the traveltimes must truly be **minima**.]

Let us define \( \tau^* \) to be the functional that yields the traveltime along the Fermat
(least-time) ray path. Fermat’s principle then states

\[
\tau^*(s) = \min_{P \subseteq \text{Paths}} \tau^P(s),
\]

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where Paths denotes the set of all continuous paths connecting the given source and receiver. The particular path that produces the minimum in (2) is denoted $P^*$. If more than one path produces the same minimum traveltime value, then $P^*$ denotes any particular member in this set of minimizing paths.

To summarize, we have

$$\tau^*(s) = \int_{P^*} s(x) \, dl^{P^*} = \min_P \int_{P} s(x) \, dl^{P} = \min_P \tau^P(s). \quad (3)$$

The traveltime functional $\tau^*(s)$ is stationary with respect to small variations in the path $P^*(s)$.

Snell’s law is a well-known and easily derived consequence of the stationarity of the traveltime functional [Feynman et al., 1963; Whitham, 1974].

**Seismic Inversion**

Suppose we have a set of observed traveltimes, $t_1, \ldots, t_m$, from $m$ source-receiver pairs in a medium of slowness $s(x)$. Let $P_i$ be the Fermat ray path connecting the $i$th source-receiver pair. Neglecting observational errors, we can write

$$\int_{P_i} s(x) \, dl^{P_i} = t_i, \quad i = 1, \ldots, m. \quad (4)$$

Given a block model of slowness, let $l_{ij}$ be the length of the $i$th ray path through the $j$th cell:

$$l_{ij} = \int_{P_i \cap \text{cell}_j} dl^{P_i}. \quad (5)$$

Given a model with $n$ cells with the $j$th cell having constant slowness $s_j$, (4) can then be written

$$\sum_{j=1}^{n} l_{ij}s_j = t_i, \quad i = 1, \ldots, m. \quad (6)$$
Note that for any given $i$, the ray-path lengths $l_{ij}$ are zero for most cells $j$, since a given ray path will in general intersect only a few of the cells in the model. Figure ?? illustrates the ray path intersections for a 2-D block model.

We can rewrite (6) in matrix notation by defining the column vectors $s$ and $t$ and the matrix $M$ as follows:

$$
s = \begin{pmatrix} s_1 \\ s_2 \\ \vdots \\ s_n \end{pmatrix}, \quad t = \begin{pmatrix} t_1 \\ t_2 \\ \vdots \\ t_m \end{pmatrix}, \quad M = \begin{pmatrix} l_{11} & l_{12} & \cdots & l_{1n} \\ l_{21} & l_{22} & \cdots & l_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ l_{m1} & l_{m2} & \cdots & l_{mn} \end{pmatrix}. \tag{7}$$

Equation (6) then becomes

$$
Ms = t. \tag{8}
$$

Equation (8) is just the discretized form of (3).

Equation (8) will be become our generic inversion problem in the remainder of this talk. The right hand side is the data $m$-vector $t$, whose elements are the $m$ traveltime measurements in this application. The left hand side has as one of its factors an $m \times n$ matrix $M$, which for this talk will be assumed known (but for real applications is often another unknown in the problem). The remaining factor is the model $n$-vector $s$, which is the model we seek in the inversion process. In practice, either $m < n$ or $m > n$ may occur in real problems. We refer to the case $m < n$ as underdetermined and the case $m > n$ as overdetermined. Most of the remainder of the article will concentrate on solving this equation approximately for $s$, and simultaneously providing some means of analyzing how well the method of approximate solution has worked.