TWO DIFFICULTIES IN INTERPRETING VECTOR AUTOREGRESSIONS

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Christopher Sims made useful comments on an earlier draft. Danny Quah performed the calculations in Section 2.

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0. Introduction

The equilibrium of a typical dynamic rational expectations model is a covariance stationary \((n \times 1)\) vector stochastic process \(z(t)\). This stochastic process determines the manner in which random shocks to the environment impinge over time on agents' decisions and ultimately upon market prices and quantities. "Surprises," or random shocks to agents' information sets prompt revisions in their contingency plans, thereby impinging on equilibrium prices and quantities.

Every \((n \times 1)\) covariance stationary stochastic process \(z(t)\) can be represented in the form of a vector autoregression (of any finite order). Consequently, it is natural to represent the equilibrium of a dynamic rational expectations model in terms of its vector autoregression. A vector autoregression recovers a vector of innovations which yield characterizations of the vector stochastic process via the "innovation accounting" techniques invented by Christopher Sims.

In interpreting these innovation accountings, it is useful to understand the connections between the innovations recovered by vector autoregressions, on the one hand, and the random shocks to private agents' information sets, on the other hand. From the viewpoint of interpreting vector autoregressions that are estimated without imposing restrictions from formal economic theories, it would be desirable if the innovations recovered by a vector autoregression could generally be expected to equal either the random shocks to agents' information sets, or else some simply interpretable functions of these random shocks.
This paper describes two important classes of theoretical models in which no such simple connections exist. In these contexts, (without explicitly imposing the restrictions implied by the economic theory), it is impossible to make correct inferences about the shocks impinging on agents' information sets. In addition to describing these situations, we briefly indicate in each case how the economic theory can be used to deduce correct inferences about the shocks impinging on agents' information sets.

Let \( z(t) \) be an \((nx1)\) vector, covariance stationary stochastic process. Imagine that \( z(t) \) is observed at discrete points in time separated by the sampling interval \( \Delta \). A vector autoregression is defined by the projection equation

\[
(0.1) \quad z(t) = \sum_{j=1}^{\infty} A_j^\Delta z(t-\Delta j) + a(t-\Delta j) \quad t=0, \pm \Delta, \pm 2\Delta, \ldots
\]

where \( a(t) \) is an \((nx1)\) vector of population residuals from the regression with \( E a(t)a(t)^T = V \), and where the \( A_j^\Delta \)'s are \((nxn)\) matrices that, in general, are uniquely determined by the orthogonality conditions (or normal equations)

\[
(0.2) \quad E z(t-\Delta j)a(t)^T = 0, \quad j \geq 1.
\]

The \( A_j^\Delta \)'s in general are "square summable," that is, they satisfy

\[
(0.3) \quad \sum_{j=1}^{\infty} \text{tr} A_j^\Delta A_j^{\Delta T} < +\infty.
\]

Equations (0.1)-(0.3) imply two important properties of \( a(t) \). First, (0.1) and (0.2) imply that

\[
E a(t)a(t-\Delta j)^T = 0 \quad j \neq 0,
\]

\[
E a(t)a(t-\Delta j)^T = a(t) \quad j = 0.
\]
so that \( a(t) \) is a vector white noise. Second, (0.1) and (0.3) imply that \( a(t) \) is in the closed linear space spanned by \( \{z(t), z(t-A), z(t-2A), \ldots\} \). Further, by successively eliminating all lagged \( z(t)'s \) from (0.1), we obtain the vector moving average representation

\[
(0.4) \quad z(t) = \sum_{j=0}^{\infty} C_j^A a(t-A^j)
\]

where the \( C_j^A \)'s are \( nxn \) matrices that satisfy

\[
- \sum_{j=0}^{\infty} C_j^A A_s = \begin{cases} 
  I & s=0 \\
  0 & s \neq 0
\end{cases}
\]

where \( A_0^A = -I \). The \( C_j^A \)'s in (0.4) satisfy

\[
(0.5) \quad \sum_{j=0}^{\infty} \text{tr} C_j^A C_j^{A^T} < +\infty.
\]

Equations (0.4) and (0.5) imply that \( z(t) \) is in the closed linear space spanned by \( \{a(t), a(t-A), a(t-2A), \ldots\} \). Thus, the closed linear space spanned by \( \{z(t), z(t-A), \ldots\} \) equals the closed linear space spanned by \( \{a(t), a(t-A), \ldots\} \). In effect, \( a(t) \) is a stochastic process that forms an orthogonal basis for the stochastic process \( z(t) \), and which is constructed from \( z(t) \) via a "Gram-Schmidt" process. The property of the vector white noise \( a(t) \) that it is contained in the linear space spanned by current and lagged \( z(t)'s \) is said to mean that "\( a(t) \) is a fundamental white noise for the \( z(t) \) process."

It is a moving average representation for \( z(t) \) in terms of a fundamental white noise which is automatically recovered by vector autoregression.\(^1\) However, there are in addition a variety of other moving average representations for \( z(t) \) of the form
CORRECT REPRESENTATION SUPPLY SHOCK

Legend

Δ QUANTITY
X PRICE
(0.6) \[ z(t) = \sum_{j=0}^{\infty} \hat{C}_j \tilde{a}(t-J) \]

where \( \tilde{a}(t) \) is an \((nx1)\) vector white noise in which the linear space spanned by \( (\tilde{a}(t), \tilde{a}(t-A), \ldots) \) is strictly larger than the linear space spanned by current and lagged \( z(t) \)'s. Current and lagged \( z(t) \)'s fail to be "fully revealing" about the \( \tilde{a}(t) \)'s in such representations.

Representation (0.4) induces the following decomposition of \( j \)-step ahead prediction errors

\[ E(\hat{z}(t)-\hat{E}_{t-j}z(t))(\hat{z}(t)-\hat{E}_{t-j}z(t))^T \]

(0.7) \[ = \sum_{k=0}^{j-1} C_k V C_k^T. \]

By studying versions of decomposition (0.7), Sims has shown how the \( j \)-step ahead prediction error variance can be decomposed into parts attributable to "innovations" in particular components of the vector \( z(t) \).

Christopher Sims has described methods for estimating vector autoregressions and for obtaining alternative fundamental moving average representations. He has also created a useful method known as "innovation accounting" that is based on decomposition (0.7). In the hands of Sims and other skilled analysts, these methods have been used successfully to detect interesting patterns in data, and to suggest possible interpretations of them in terms of the responses of systems of people to surprise events.

This paper focuses on the question of whether dynamic economic theories readily appear in the form of a fundamental moving average representation (0.4), so that the vector white
noises $a(t)$ recovered by vector autoregressions are potentially interpretable in terms of the white noises impinging on the information sets of the agents imagined to populate the economic model. This question is important because it influences the ease with which one can interpret the variance decompositions (or innovation accounts) and the response to innovations $a(t)$ that are associated with the fundamental moving average $(0.4)$.

This paper is organized as follows. Section 1 describes a class of discrete-time models whose equilibria can be represented in the form

$$(0.8) \quad z_t = \sum_{j=0}^{\infty} D_j \varepsilon_{t-j}$$

where $\varepsilon_t$ is an (nx1) vector white noise; $D_j$ is an (nxn) matrix for each $n$; and $\text{tr} \sum_{j=0}^{\infty} D_j D_j^T < \infty$. Here $\varepsilon_t$ represents a set of shocks to agents' information sets. We study how the $\varepsilon_t$ of representation $(0.8)$ are related to the $a(t)$ of the (Wold) representation $(0.4)$, and how the $D_j$'s of $(0.8)$ are related to the $C^\Delta_j$'s of $(0.4)$. We describe contexts in which $a(t)$ fails to match up with $\varepsilon_t$ and $C^\Delta_j$ fails to match up with $D_j$ because $z(t)$ fails to be fully revealing about $w(t)$. Such examples were encountered earlier by Hansen and Sargent [1980], Futia [1981], and Townsend [1981]. The discussion in Section 1 assumes that the sampling interval $\Delta$ equals the sampling interval in terms of which the economic model is correctly specified.

Section 1 describes a class of continuous time model whose equilibria are represented in the form

$$(0.9) \quad z_t = \int_0^{\infty} P(\tau) w(t-\tau)$$
where \( w(t) \) is a continuous time white noise and \( p(\tau) \) is an \((nx1)\) function satisfying \( \text{tr} \int_0^\infty p(\tau)p(\tau)^T \, d\tau < +\infty \). In (0.9), \( w(t) \) represents shocks to agents' information set. It is supposed that economic decisions occur in continuous time according to (0.9), but that the econometrician only possesses data at discrete intervals of time. Section 2 studies the relationship between the \( w(t) \) of (0.9) and the \( a(t) \) of (0.4), and also the relationship between \( p(\tau) \) of (0.9) and the \( C_j^A \) of (0.4). In general, these pairs of objects do not match up in ways that can be determined without the imposition of restrictions from a dynamic economic theory.

1. Unrevealing Stochastic Processes

We consider a class of discrete time linear rational expectations models that can be represented as the solution of the following pair of stochastic difference equations

(1.1) \[ H(L) y_t = E_t J(L^{-1})^{-1} px_t \]

\[ x_t = K(L) \epsilon_t \]

where

\[ H(L) = H_0 + H_1 L + \ldots + H_m L^m \]

(1.2) \[ J(L) = J_0 + J_1 L + \ldots + J_m L^m \]

\[ K(L) = \sum_{j=0}^{\infty} K_j L^j, \quad K_0 = I \]

\[ \epsilon_t = x_t - E x_t | x_{t-1}, x_{t-2}, \ldots \]
In (1.1), \( y_t \) is an \( n_1 \times 1 \) vector, while \( x_t \) is an \( n_2 \times 1 \) vector. In (1.2), \( J_j \) and \( H_j \) are \((n_1 \times n_1)\) matrices, while \( K_j \) is an \((n_2 \times n_2)\) matrix. In (1.1), \( p \) is an \((n_1 \times n_2)\) matrix. We assume that the zeroes of \( \det H(z) \) lie outside the unit circle, that those of \( \det J(z) \) lie inside the unit circle, and that those of \( \det K(z) \) do not lie inside the unit circle.

A wide variety of discrete time linear rational expectations models are special cases of (1.1). For example, interrelated factor demand versions of Lucas-Prescott equilibrium models are special cases with \( J(L^{-1}) = H(L^{-1})^T \) and with \( H(L^{-1})^T H(L) \) being the matrix factorization of the Euler equation that is solved by the fictitious social planner (see Hansen and Sargent [1981] and Eichenbaum [1981] for some examples). Kydland-Prescott equilibria with feedback from market-wide variables to forcing variables that individual agents face parametrically form a class of examples with \( H(L^{-1})^T \neq J(L^{-1}) \) (see Hansen and Sargent [1984]). Other examples with \( H(L^{-1})^T \neq J(L^{-1}) \) arise in the context of various dominant player equilibria of linear quadratic differential games (see Hansen, Epple, and Roberds [1984]). Finally, market equilibrium models of the Kennan [1982]–Sargent [1979] variety, an example of which is studied below, solve a version of (1.1) with \( H(L^{-1})^T \neq J(L) \). Models of this general class are studied by Whiteman [1983].

Hansen and Sargent [1981] have displayed a convenient representation of the solution of models related but not identical to (1.1). To adapt their results, first obtain the partial fractions representation of \( J(z^{-1})^{-1} \). We have \( J(z^{-1})^{-1} = \det J(z^{-1}) \operatorname{adj} J(z^{-1}) \). Let
\[ \det J(z^{-1}) = \lambda_0 (1 - \lambda_1 z^{-1}) \cdots (1 - \lambda_k z^{-1}) \]

where \( k = m_2 \cdot n_1 \) and \( |\lambda_j| < 1 \) for \( j = 1, \ldots, k \). The \( \lambda_j \)'s are the zeroes of \( \det J(z^{-1}) \). Then we have

\[ (1.3) \quad J(z^{-1})^{-1} = \sum_{j=1}^{k} \frac{M_j}{1 - \lambda_j z^{-1}} \]

where

\[ (1.4) \quad M_j = \lim_{z \to \lambda_j} J(z^{-1})^{-1} (1 - \lambda_j z^{-1}). \]

Substitute (1.3) into (1.1) to obtain

\[ (1.5) \quad H(L) y_t = E_t \sum_{k=1}^{k} \frac{M_j}{1 - \lambda_j L^{-1}} p X_t. \]

Hansen and Sargent [1980] establish that

\[ (1.6) \quad E_t \frac{M_j}{1 - \lambda_j L^{-1}} p X_t = M_j p \left( LK(L) - \lambda_j K(\lambda_j) \right) \epsilon_t. \]

Define the operator \( M \) by

\[ (1.7) \quad M(K(L)) = \sum_{j=1}^{k} M_j p \left( LK(L) - \lambda_j K(\lambda_j) \right). \]

Then, using (1.5), (1.6), and (1.7) we have the representation of the solution

\[ (1.8) \quad H(L) y_t = M(K(L)) \epsilon_t \]

\[ x_t = K(L) \epsilon_t. \]

A vector stochastic process \((y_t^T, x_t^T)\) governed by (1.8) generally has a singular spectral density at all frequencies
because \((y_t^T, x_t^T)\) consists of \(n_1 + n_2\) variables being driven by only \(n_2\) white noises. Such a model implies that various of the first \(n_1\) equations of the following model, which is equivalent to (1.8),

\[
H(L) y_t = M(K(L)) K(L)^{-1} x_t
\]

\[
x_t = K(L) \epsilon_t
\]

will fit perfectly (i.e., possess sample \(R^2's\) of 1). To avoid this implication of no errors in various of the equations of the model, while still retaining the model, one path that has been suggested is to assume that the econometrician seeks to estimate (1.8), but that he possesses data only on a subset of the variables in \((y_t, x_t)\). (See Hansen and Sargent [1980].) One common procedure, but not the only one possible, is the following one described by Hansen and Sargent [1980]. Assume that (1.8) holds, but that the econometrician only has data on a subset of observations \(x_{2t}\) of \(x_t\). Further suppose that the second equation of (1.1) can be partitioned and restricted as

\[
x_t = \begin{pmatrix} x_{1t} \\ x_{2t} \end{pmatrix} = \begin{pmatrix} K_1(L) & 0 \\ 0 & K_2(L) \end{pmatrix} \begin{pmatrix} \epsilon_{1t} \\ \epsilon_{2t} \end{pmatrix}.
\]

Then (1.8) assumes a special form which can be represented as

\[
\begin{pmatrix} H(L) & 0 \\ 0 & I \end{pmatrix} \begin{pmatrix} y_t \\ x_{2t} \end{pmatrix} = \begin{pmatrix} M(K_1(L)) & M(K_2(L)) \\ 0 & K_2(L) \end{pmatrix} \begin{pmatrix} \epsilon_{1t} \\ \epsilon_{2t} \end{pmatrix}.
\]

The idea is to imagine that the econometrician is short of observations on a sufficient number of series, those forming \(x_{1t}\), to make the \((y_t^T, x_{2t}^T)\) process described by (1.10) have a nonsingular
spectral density matrix at all frequencies. To accomplish this, it will generally be sufficient if the dimension of the vector of variables of \((y_t^T, x_{2t}^T)\) is less than or equal to the dimension of \((\varepsilon_{1t}^T, \varepsilon_{2t}^T)\). For the argument below, we will consider the case often encountered in practice in which \((y_t^T, x_{2t}^T)\) and \((\varepsilon_{1t}^T, \varepsilon_{2t}^T)\) have equal dimensions. Thus we assume that \(x_{1t}\) is an \((n_1 \times 1)\) vector, so that \(\varepsilon_{1t}\) is an \((n_1 \times 1)\) vector of white noises.\(^5\)

Equation (1.10) implies the moving average representation for \((y_t^T, x_{2t}^T)\)

\[
\begin{pmatrix}
    y_t \\
    x_{2t}
\end{pmatrix}
= \begin{pmatrix}
    (H(L)^{-1}M(K_1(L)) & H(L)^{-1}M(K_2(L)) \\
    0 & K_2(L)
\end{pmatrix}
\begin{pmatrix}
    \varepsilon_{1t} \\
    \varepsilon_{2t}
\end{pmatrix}
\]

Equation (1.11) is a moving average that expresses \((y_t^T, x_{2t}^T)\) in terms of current and lagged values of the white noises \((\varepsilon_{1t}, \varepsilon_{2t})\) that are the innovations in the information sets \((x_{1t}, x_{2t})\) of the agents in the model. Equivalently, \((\varepsilon_{1t}, \varepsilon_{2t})\) are "fundamental for \(x_{1t}, x_{2t}\)" the one step-ahead errors in predicting \(x_{1t}, x_{2t}\) from their own pasts being expressible as linear combinations of \(\varepsilon_{1t}, \varepsilon_{2t}\).

Granted that the linear space spanned by current and lagged \((\varepsilon_{1t}, \varepsilon_{2t})\) equals that spanned by current and lagged values of the agents' information \((x_{1t}, x_{2t})\), there remains the question of whether this space equals that spanned by current and lagged values of the econometrician's information \((y_t^T, x_{2t}^T)\). From (1.11) and as is evident from the theory used to construct (1.11), the former space is included in the latter. The question is whether they are equal. This question is an important one from the view-
point of interpreting vector autoregressions because a vector autoregression by construction would recover a vector moving average for \((y_t, x_{2t})\) that is driven by a vector white noise \(a_t\) that is fundamental for \((y_t, x_{2t})\), i.e., one that is in the linear space spanned by current and lagged values of \((y_t, x_{2t})\). If this space is smaller than the one spanned by current and lagged values of agents' information \((\varepsilon_{1t}, \varepsilon_{2t})\), then the moving average representation recovered by the vector autoregression will in general give a distorted impression of the response of the system to surprises from agents' viewpoint.

The vector white noise \((\varepsilon_{1t}, \varepsilon_{2t})\) is fundamental for \((y_t, x_{2t})\) if and only if the zeroes of

\[
\begin{vmatrix}
H(z)^{-1}M(K_1(z)) & H(z)^{-1}M(K_2(z)) \\
0 & K_2(z)
\end{vmatrix}
= \det H(z)^{-1}M(K(z)) \cdot \det K_2(z)
\]

do not lie inside the unit circle. The zeroes of \(\det K_2(z)\) do not lie inside the unit circle by assumption, and \(\det H(z)^{-1} = 1/\det H(z)\) is a function with all its poles outside the unit circle. Therefore, the necessary and sufficient condition that \((\varepsilon_{1t}, \varepsilon_{2t})\) be fundamental for \((y_t, x_{2t})\) is that

\[
(1.12) \quad \det M(K_1(z^0)) = 0 \Rightarrow |z^0| > 1.
\]

or, equivalently, using (1.7),

\[
(1.12') \quad \det \sum_{j=1}^{k} M_j \frac{z^0 K_1(z^0) - \lambda_j K_1(\lambda_j)}{z^0 - \lambda_j} = 0 \Rightarrow |z^0| > 1.
\]
In general, condition (1.12) is not satisfied. For some specifications of $K_1(L)$ and $J(L^{-1})$, which determines the $M_j$, $\lambda_j$ via (1.3)-(1.4), condition (1.12) is met, while for others, it is not met. Hansen and Sargent [1980] encountered a class of examples where (1.12) isn't met. Furthermore, the class of cases for which (1.12) fails to be met is not thin in any natural sense. Our conclusion is that for the class of models defined by (1.1)-(1.2), the moving average representation (1.11) that is expressed in terms of the white noises that are fundamental for agents' information sets in general cannot be expected to be fundamental for the econometrician's data set $(y_t, x_{2t})$. Equivalently, current and lagged values of $(y_t, x_{2t})$ fail to be fully revealing of current and lagged values of $(\varepsilon_{1t}, \varepsilon_{2t})$.

For convenience, let us rewrite (1.10) as

(1.13) $S(L) z_t = R(L) \varepsilon_t$

where

$S(L) = \begin{pmatrix} H(L) & 0 \\ 0 & I \end{pmatrix}, \quad R(L) = \begin{pmatrix} M(K_1(L)) & M(K_2(L)) \\ 0 & K_2(L) \end{pmatrix}$

$\varepsilon_t = \begin{pmatrix} \varepsilon_{1t} \\ \varepsilon_{2t} \end{pmatrix}, \quad z_t = \begin{pmatrix} y_t \\ x_{2t} \end{pmatrix}$.

The condition that $\varepsilon_t$ be fundamental for $z_t$ is then expressible as the condition that the zeroes of $\det R(z)$ not lie inside the unit circle. If this condition is violated, then a Wold representation for $z_t$, which is what is recovered via vector autoregression, will
be related to representation (1.13) as follows. It is possible to show that there exists a matrix polynomial $G(L)$ satisfying

$$G(L) G(L^{-1})^T = I$$

and such that the matrix

$$R^*(L) = R(L) G(L)$$

has a characteristic polynomial $\det R^*(z)$, none of whose zeroes lie inside the unit circle. Let us define

$$\varepsilon^*_t = G(L^{-1}) \varepsilon_t.$$ 

Then a Wold representation corresponding to (1.13) is

$$(1.14) \quad S(L) z_t = R^*(L) \varepsilon^*_t.$$ 

In general, $R^*(L) = R(L) G(L)$ is quite a complicated transformation of $R(L)$. Further, $\varepsilon^*_t$ is a complicated function of past, present, and future $\varepsilon_t$'s, one that has the effect of diminishing the amount of information in $\{\varepsilon^*_t, \varepsilon^*_t-1, \ldots\}$ relative to that contained in $\{\varepsilon_t, \varepsilon_{t-1}, \ldots\}$.

We now describe a concrete hypothetical numerical example, one in which the econometrician observes no $x$'s, only $y$'s, so that (1.10) takes the special form

$$(1.15) \quad H(L) y_t = M(C_1(L)) \varepsilon_{1t}.$$ 

The model is one of the dynamics of demand and supply, and is similar to that studied by Sargent [1979] and Kennan [1982]. Suppose that there are two types of agents, each of whom
solves a quadratic optimum problem. The first agent, the supplier, maximizes the objective

$$E_0 \sum_{t=0}^{\infty} \beta^t \left[ p_t q_t - s_t q_t - \frac{h_s}{2} q_t^2 - \frac{g_s}{2} [(1-L)q_t]^2 \right]$$

where $p_t$ is the price at time $t$, $q_t$ is the quantity supplied at time $t$, $s_t$ is a supply shock at time $t$, $\beta$ is a discount factor between zero and one, and $h_s$ and $g_s$ are parameters of the cost function. The term $[(1-L)q_t]^2$ is introduced to capture costs in adjusting the output from period $t-1$ to period $t$. The supplier views the price and the supply shock processes as if they were uncontrollable. All variables are treated as deviations from their population means. The stochastic Euler equation for the supplier's optimum is given by

$$-E_t \left\{ [h_s + g_s (1-L)(1-\beta^{-1})]q_t \right\} + p_t = s_t.$$ 

The second agent, the demander, maximizes the objective

$$E_0 \sum_{t=0}^{\infty} \beta^t \left[ p_t q_t - d_t q_t - \frac{h_d}{2} q_t^2 - \frac{g_d}{2} [G^*(L)q_t]^2 \right]$$

where $d_t$ is a shock to preferences, $h_d$ and $g_d$ are preference parameters, and $G^*(L) = 1 + .8L + .6L^2 + .4L^3 + .2L^4$. The price term enters into the demander's objective because of an implicit substitution from the demander's budget constraint. The term $G^*(L)$ is introduced to capture the notion that purchases of $q$ in recent past time periods give rise to services today. The demander treats the price and demand shock processes as if they were uncontrollable. The stochastic Euler equation for the demander's optimum problem is given by
To complete model specification, we specify the stochastic law of motion for the forcing processes, i.e., the demand and supply shocks. These shocks are assumed to satisfy

\begin{equation}
(1.18) \quad s_t = B_s(L)w_{st}
\end{equation}

\begin{equation}
(1.19) \quad d_t = B_d(L)w_{dt}
\end{equation}

where \( B_s(z) \) and \( B_d(z) \) are scalar polynomials with zeroes that are outside the unit circle. The \( w_{st} \) and \( w_{dt} \) processes are mutually uncorrelated white noises so that \( w_{st} \) is the innovation in the supply shock, \( s_t \), and \( w_{dt} \) is the innovation in demand shock, \( d_t \). Economic agents are presumed to observe current and past values of both shocks and hence also the innovations in both shocks.

This model fits into our general set up (1.1) as follows. Let

\[ y_t = \begin{pmatrix} q_t \\ p_t \end{pmatrix}, \quad \epsilon_t = \begin{pmatrix} w_{st} \\ w_{dt} \end{pmatrix}. \]

\[ K_2(L) = 0, \quad x_{2t} = 0. \]

\[ K_1(L) = \begin{pmatrix} B_s(L) & 0 \\ 0 & B_d(L) \end{pmatrix}. \]

Define a matrix polynomial

\[ E(L) = \begin{pmatrix} -(h_s + g_s(1-L)(1-\beta L^{-1})) & 1 \\ -(h_d + g_d G(L) G(\beta L^{-1})) & -1 \end{pmatrix}. \]
Then polynomial matrices $H(L)$ and $J(L)$ that are one-sided in nonnegative powers of $L$ can be found such that

$$J(L^{-1}) H(L) = E(L)$$

and such that the zeroes of $\det J(z)$ lie inside the unit circle, while the zeroes of $\det H(z)$ lie outside the unit circle. (See Whiteman [1983], and Gohberg, Lancaster, and Rodman [1982] for proofs of the existence of such a matrix factorization, and for description of algorithms for achieving the factorization.)

In this model, it is possible for the demand and supply shocks to generate an information set that is strictly larger than that generated by current and past quantities and prices. So an econometrician using innovation accounts derived from observations on quantities and prices may not obtain innovations that are linear combinations of the contemporaneous innovations to the demand and supply shocks.

The equilibrium of the model has representation

$$(1.18) \quad S(L) \begin{pmatrix} q_t \\ p_t \end{pmatrix} = R(L) \begin{pmatrix} w_{st} \\ v_{dt} \end{pmatrix}$$

where $S(L) = H(L)$ and $R(L) = M(C_1(L))$, and where $S(z)$ is a $(2\times2)$ fourth-order matrix polynomial in $L$, with the zeroes of $\det S(z)$ outside the unit circle. The zeroes of $\det R(z)$ can be on either side of the unit circle in this example. Only when the zeroes of $\det R(z)$ are not inside the unit circle can the one-step ahead forecast errors from the vector autoregression of prices and quantities be expressed as linear combinations of the contemporaneous demand and supply shock innovations $(w_{st}, v_{dt})$. 

$\frac{8}{8}$
In our numerical example, we chose the parameter specification

\[ h_s = h_d = 1 \]

\[ \varepsilon_s = \varepsilon_d = 4 \]

\[ B_s(L) = (1+\lambda_1 L)(1+\lambda_2 L)(1+\lambda_3 L) \]

\[ B_d(L) = (1+\mu_1 L)(1+\mu_2 L)(1+\mu_3 L). \]

The parameters \( \lambda_j \) and \( \mu_j \) for \( j = 1, 2, 3 \) were selected among the values \{.2, .4, .6, .8\}. In our systematic search among various combinations of \( \lambda_j \)'s and \( \mu_j \)'s, we found that there always was at least one zero of \( \text{det} \, R(z) \) that was inside the unit circle. For instance when

\[ (\lambda_1, \lambda_2, \lambda_3) = (.6, .4, .2) \]

and

\[ (\mu_1, \mu_2, \mu_3) = (.8, .6, .4) \]

det \( R(z) \) had a zero at -.339.

For this example, Table 1 displays both \( R(L) \) and \( R^*(L) \) for a fundamental representation, as well as \( S(L) \). Figure 1 graphs the moving average coefficients \( S(L)^{-1} R^*(L) \) for two alternative orthogonalization orders\(^9\) used to normalize \( R^*(L) \). The "innovation in price" for the fundamental representation for \( (q_t, p_t) \) traces out a moving average response that mimics fairly well the responses of the \( (q_t, p_t) \) system to a supply shock. This is true for either orthogonalization order. The response to an
FUNDAMENTAL REPRESENTATION: QUANTITY FIRST INNOVATION IN PRICE
CORRECT REPRESENTATION DEMAND SHOCK

Legend

QUANTITY

PRICE
FUNDAMENTAL REPRESENTATION: QUANTITY FIRST INNOVATION IN QUANTITY

Legend
\(\Delta\) QUANTITY
\(\times\) PRICE
FUNDAMENTAL REPRESENTATION: PRICE FIRST INNOVATION IN QUANTITY
innovation in quantity fails to resemble the response of the system to a demand shock.

In examples such as this one, in which some of the zeroes of det $R(z)$ lie inside the unit circle, the innovations $(w_{st}, w_{dt})$ cannot be expressed as linear combinations of the innovations $a_t$ in the vector autoregression. If all of the zeroes of det $R(z)$ are outside the unit circle, then $w_t = (w_{st}, w_{dt})$ is related to $a_t$ by an exact linear relation, $w_t = Pa_t$, although to know the matrix $P$, one needs to know the parameters of the agents' Euler equations, and the laws of motion of the supply and demand shocks. When the zeroes of det $R(z)$ are all outside the unit circle, there is at least some comfort in the fact that there exists some linear transformation of the $a_t$'s that would recover the innovations to the agents in the model.

Remedies in Discrete Time

The preceding difficulty can be circumvented if a sufficiently restrictive dynamic economic theory is imposed during estimation. Hansen and Sargent [1980] describe methods for estimating $S(L)$ and $R(L)$ subject to extensive cross-equation restrictions of the rational expectations variety. The approach is to use the method of maximum likelihood to estimate free parameters of preferences and constraint sets, of which the parameters of $S(L)$ and $R(L)$ are in turn functions. These methods do not require that the zeroes of det $R(z)$ be restricted, and in particular are capable of recovering good estimates of $R(z)$ even when some of the zeroes of det $R(z)$ are inside the unit circle.
Given such an estimate of \((S(L), R(L))\), it would be possible to recover consistent estimates of the \(\varepsilon_t\) that are innovations to agents' information sets by using the following procedure. Let \(R(L)\) be factored as \(R(L) = R_1(L) R_2(L)\) where the zeroes of \(\det R_1(z)\) lie inside the unit circle, while those of \(\det R_2(z)\) lie outside the unit circle. Such a factorization exists by results described in Gohberg, Lancaster, and Rodman [1982]. Then, using the above decomposition in \(S(L)z_t = R(L)\varepsilon_t\), we can obtain

\[
R_2(L)\varepsilon_t = R_1(L)^{-1} S(L)z_t,
\]

where \(R_1(L)^{-1}\) is interpreted as the "stable inverse" in nonpositive powers of \(L\). This equation expresses \(\varepsilon_t\) as a square summable sum of past, present, and future values \(z_t\). By using this equation, it would be possible to recover estimates of the innovations \(\varepsilon_t\) to agents' information sets from a time series record on \(\{z_t\}\). This equation once again makes the point that \(\varepsilon_t\) fails to lie in the linear space spanned by current and lagged \(z_t\) only.

2. Time Aggregation

Consider a linear economic model that is formulated in continuous time, and which can be represented as

\[
(2.1) \quad z(t) = \int_0^\infty p(\tau)w(t-\tau)d\tau
\]

where \(z(t)\) is an \((n \times 1)\) vector stochastic process, \(w(t)\) is an \((m \times 1)\) vector white noise with \(Ew(t)w(t-s)^T = \delta(t-s)I\), \(\delta\) is the Dirac delta generalized function, and \(p(\tau)\) is an \((n \times m)\) matrix function that satisfies \(\int \text{trace } p(\tau)p(\tau)^T d\tau < +\infty\). We let \(P(s) = \int_0^\infty e^{-s\tau}p(\tau)d\tau\), i.e., \(P(s)\) is the Laplace transform of \(p(\tau)\). Some-
times we shall find it convenient to write (2.1) in operator notation

\[ z(t) = P(D)w(t) \]

where \( D \) is the derivative operator. We shall assume that the Laplace transform \( P(s) \) has no zeroes in the right half of the complex plane. This guarantees that square integrable functionals of \((z(t-s), s > 0)\) and of \((w(t-s), s > 0)\) span the same linear space, and is equivalent to specifying that (2.1) is a Wold representation for (2.1).

A variety of continuous time stochastic linear rational expectations models have equilibria that assume the form of the representations (2.1) or (2.2). Hansen and Sargent [1981] provide some examples. In these examples, the continuous time white noises \( w(t) \) often have interpretations as innovations in the uncontrollable processes that agents care about forecasting, and which stochastically drive the model. These include processes that are imagined to be observable to both the econometrician and the private agent (e.g., various relative prices and quantities) and also those which are observable to private agents but are hidden from the econometrician (e.g., random disturbances to technologies, preferences, and maybe even particular factors of production such as "effort" or capital of specific kinds). The \( w(t) \) process is economically interpretable as the continuous time innovation to private agents, because the forecast error of the variables in the model over any horizon \( t + \tau \) which the private agents are assumed to make at \( t \) can be expressed as a weighted sum
of w(s), t < s < t + \tau. Thus, to private agents the w(t) process represents "news" or "surprises."

In rational expectations models, typically there are extensive restrictions across the rows of P(D). In general these restrictions leave open the possibility that the current and lagged values of the w(t) process span a larger linear space than do current and lagged values of the z(t) process. This outcome can possibly occur even if the dimension m of the w(t) process is less than or equal to the dimension n of the z(t) process. This is the continuous time version of the phenomenon that we treated for discrete time in the previous section. In the present section, we ignore this phenomenon, by assuming that P(s) has no zeroes in the right half of the complex plane.

For this continuous time specification, there exists a discrete time moving average representation

\[(2.3) \quad z_t = C(L)a_t\]

where C(L) is an infinite order, (nxn) polynomial in the lag operator L, where a_t is a vector white noise with Ea_t a_t^T = W, and where a_t = z_t - \hat{E}[z_t|z_{t-1},...]. The lag operator C(L) and the positive semi-definite matrix W solve the following equation, subject to the side condition that the zeroes of det C(z) do not lie inside the unit circle: 10/

\[(2.4) \quad C(e^{-iw}) W C(e^{iw})^T = \sum_{j=-\infty}^{+\infty} P(iw) P(-iw)^T.\]

When z_t has a discrete time autoregressive representation, the discrete time innovations a_t are related to the w(t) process by the formula
\[ a_t = C(L)^{-1} P(D) w(t) \]

or

\[ (2.5) \quad a_t = V(L) P(D) w(t) = V(L) \int_0^\infty p(\tau) w(t-\tau) d\tau \]

where we have defined \( V(L) = C(L)^{-1} = \sum_{j=0}^\infty V_j L^j \), \( V_0 = I \). Here \(-V_j\) is the \( n \times n \) matrix of coefficients on the j\(^{th}\) lag in the vector autoregression for \( z \). It follows directly upon writing out (2.5) that

\[ (2.6) \quad a_t = \int_0^\infty f(\tau) w(t-\tau) d\tau \]

where \( f(\tau) = \sum_{j=0}^\infty V_j p(\tau-j) \)

It also follows from (2.6) and the identity for integer \( t \), \( C(L) a_t = P(D) w_t \), that

\[ (2.8) \quad p(\tau) = \sum_{j=0}^\infty C_j f(\tau-j). \]

Equations (2.6) and (2.7) show how the discrete time innovation \( a_t \) in general reflects all past values of the continuous time innovation \( w(t) \).

Analyses of vector autoregressions often proceed by summarizing the shape of \( C(L) \) in various ways, and attempting to interpret that shape. The innovation accounting methods of Sims, based on decomposition (0.7), are good examples of procedures that summarize the shape of \( C(L) \). From the viewpoint of interpreting discrete time vector autoregressions in terms of the economic forces acting on individual agents, it would be desirable if the
discrete time and continuous time moving average representations were to match up in some simple and interpretable ways. In particular, the following two distinct but related features would be desirable. First, it would be desirable if the discrete time innovations $a_t$ closely reflected the behavior of $w(s)$ near $t$. Probably the most desirable outcome would be if $a_t$ could be expressed as

$$a_t = \int_0^1 f(\tau) w(t-\tau) d\tau \tag{2.9}$$

so that in (2.6), $f(\tau) = 0$ for $\tau > 1$. In that case, $a_t$ would be a weighted sum of the continuous time innovations over the unit forecast interval. It would be even more desirable if (2.9) were to hold with $f(\tau) = p(\tau)$, for then $a_t$ would equal the one step ahead forecast error from the continuous time system. Second, assuming a smooth $p(\tau)$ function, it would be desirable if the discrete time moving average coefficients $\{C_0, C_1, C_2, \ldots\}$ resemble a sampled version of the continuous time moving average kernel $\{p(\tau), \tau \geq 0\}$. This is desirable because the pattern of the $C_j$'s would then faithfully reflect the response of the system to innovations in continuous time. We shall consider each of these desiderata in turn.

We first study conditions under which $f(\tau) = 0$ for $\tau > 1$. Consider the decomposition

$$a_t = z(t) - \hat{E}[z(t)|w(t-s), s \geq 1]$$

$$+ \hat{E}[z(t)|w(t-s), s \geq 1] - \hat{E}[z_t|z_{t-1}, \ldots]$$

$$= \int_0^1 p(\tau) w(t-\tau) d\tau + \int_1^\infty p(\tau) w(t-\tau) d\tau$$
This last equality implies that if (2.9) is to hold it must be the case that

\[(2.10) \quad \hat{E}[z(t)|w(t-s), s>1] = \hat{E}[z_t|z_{t-1}, \ldots],\]

which in turn implies that \( p(\tau) = f(\tau) \) for \( 0 < \tau < 1 \). The interpretation of requirement (2.10) is that the discrete time and continuous time forecasts of \( z(t) \) over a unit time interval coincide.

When condition (2.9) is met, the link between \( P(D) \) and \( C(L) \) is particularly simple. Using \( f(\tau) = 0 \) for \( \tau > 1 \), equation (2.8) becomes

\[(2.11) \quad p(\tau) = c_j f(\tau-j) \text{ for } j < \tau < j+1.\]

Equation (2.11) implies that for the particular class of continuous time processes for which \( f(\tau) = 0 \) for \( \tau > 1 \), the continuous time moving average coefficients are completely determined by the discrete time moving average coefficients and the function \( f(\tau) \) defined on the unit interval. The aliasing problem is manifested in this relationship since \( f(\tau) \) cannot be inferred from discrete time data. In the absence of additional restrictions, all functions \( f(\tau) \) that satisfy

\[\int_0^1 f(\tau)f(\tau)^T d\tau = W\]

are observationally equivalent. Relation (2.11) also implies that in general, without some more restrictions on \( p(\tau) \), condition (2.9) does not place any restrictions on the discrete time moving average coefficients.
However, in many (if not most) applications, it is usual to impose the additional requirement that the continuous time moving average coefficients be a continuous function of $t$.\textsuperscript{12} This requirement together with (2.11) then imposes a very stringent restriction on the discrete time moving average representation. In particular, (2.11) then implies that

\begin{equation}
C_j f(0) = C_{j-1} f(1)
\end{equation}

where $f(t)$ is now a continuous function on the unit interval. When $w(t)$ and $z(t)$ have the same dimension ($m=n$) and $f(0)$ is nonsingular, relation (2.12) implies that

\begin{equation}
C_j = [f(1)f(0)^{-1}]^j
\end{equation}

and

\begin{equation}
C(L) = [I-f(1)f(0)^{-1}L]^{-1}.
\end{equation}

Hence, if (2.9) is to hold, the discrete time process must have a first order autoregressive representation. We have therefore established that condition (2.9) and the continuity requirement on $p(t)$ substantially restrict not only the admissible continuous time moving average coefficients but the admissible discrete time moving average coefficients as well.

Thus, with a continuous $p(t)$ function, in general, relation (2.9) does not hold. Instead, $a_t$ given by (2.6) is a function of all current and past $w(t)$'s, a function whose complications can pose problems in several interrelated ways for interpreting $a_t$ in terms of the continuous time noises $w(t)$ that are
imagined to impinge on agents in the model. First, as in the
discrete time case, the process \( w(t) \) need not be fundamental for
\( z(t) \) in continuous time. Second, the matrix function \( f(\tau) \) in (6)
is not usually diagonal, so that each component of \( a_t \) in general
is a function of all of the components of \( w(t) \). This is a version
of what Geweke has characterized as "contamination," which occurs
in the context of the aggregation over time of several inter-
related distributed lags. It is also related to the well-known
phenomenon that aggregation over time generally leads to Granger-
causality of discrete sampled \( y \) to \( x \) even though \( y \) fails to
Granger-cause \( x \) in continuous time. Third, the matrix function
\( f(\tau) \) in (2.5) in general is nonzero for all values of \( \tau > 0 \), so
that \( a_t \) in general depends on values of \( w(t-\tau) \) in the remote past.

We now turn to our second desideratum, namely that the
sequence \( \{C_j\}_{j=0}^\infty \) resemble a sampled version of the function
\( p(\tau) \). For studying this matter, we set \( m = n \), because we are
interested in the circumstances under which \( \{C_j\} \) fails to reflect
\( p(\tau) \) even when the number of white noises \( n \) in \( a_t \) equals the
number \( m \) in \( w(t) \). We can represent most of the issues here with a
univariate example, and so set \( m = n = 1 \) in most of our discus-
sion. It is also convenient to study the case in which \( z_t \) has a
rational spectral density in continuous time. Thus we assume that

\[
(2.13) \quad \theta(D)z_t = \psi(D) w(t)
\]

where \( z_t \) is a scalar stochastic process, and \( \theta(s) = (s-\lambda_1) (s-\lambda_2)
\ldots (s-\lambda_r) \), \( \psi(s) = \psi_0 + \psi_1 s + \ldots + \psi_{r-1} s^{r-1} \). We assume that
the real parts of \( \lambda_1, \ldots, \lambda_r \), which are the zeroes of \( \theta(s) \), are
less than zero, but that the real parts of the zeroes of $\psi(s)$ are unrestricted. Only if the real parts of the zeroes of $\psi(s)$ are less than zero do current and past values of $z(t)$ and $w(t)$ span the same linear space. If any zeroes of $\psi(s)$ exceed zero in real past, then current and lagged $w(t)$ span a larger space than do current and lagged $z(t)$. The above equation can be expressed as

$$z_t = P(D) w(t) \tag{2.14}$$

where $P(D) = \psi(D)/\theta(D)$. A partial fraction representation of $p(D)$ is

$$P(D) = \sum_{j=1}^{r} \frac{\delta_j}{D - \lambda_j} \tag{2.15}$$

where

$$\delta_j = \lim_{s \to \lambda_j} P(s) (s - \lambda_j) \tag{2.16}$$

We therefore have

$$p(\tau) = \sum_{j=1}^{r} \delta_j e^{\lambda_j \tau} \tag{2.17}$$

Thus, the weighting function $p(\tau)$ in the continuous time moving average representation is a sum of $r$ exponentially decaying functions. Our object will now be to get an analogous expression to (2.17) for the discrete time coefficients $B_k$.

It is known that the discrete time process $z_t$ implied by (2.13) is an $r^{th}$ order autoregressive, $(r-1)$ order moving average process. Let this be $z_t = \frac{c(L)}{d(L)} a_t$ where $c(L) = \sum_{j=0}^{r-1} c_j L^j$, $d(L) = \sum_{j=0}^{r} d_j L^j$. To find this representation, we must use (2.4). Hansen
and Sargent [1981] show that for the process (2.13), the term on the left side of (2.4) can be represented

\[
\sum_{j=-\infty}^{\infty} P(iw) P(-iw) = \sum_{j=1}^{r} \left[ \frac{w_j}{1-e^{j\lambda e^{-i\omega}}} + \frac{w_j e^{j\lambda e^{i\omega}}}{1-e^{j\lambda e^{i\omega}}} \right]
\]

where

\[ w_j = \lim_{s \to \lambda_j} P(s) P(-s) (s - \lambda_j). \]

Letting \( z = e^{-i\omega} \), to find the required mixed moving average autoregressive representation, we must solve

\[
\frac{c(z)c(z^{-1})}{d(z)d(z^{-1})} = \sum_{j=1}^{r} \left[ \frac{w_j}{1-e^{j\lambda z}} + \frac{w_j e^{j\lambda z^{-1}}}{1-e^{j\lambda z^{-1}}} \right]
\]

subject to the condition that the zeroes of \( c(z) \) and \( d(z) \) all lie outside the unit circle. The term on the right side of (2.8) can be expressed as

\[
\sum_{j=1}^{r} w_j \prod_{k \neq j} (1-\alpha_j z) \prod_{k \neq j} (1-\alpha_j z^{-1}) + \sum_{j=1}^{r} w_j \alpha_j \prod_{j=1}^{r} \prod_{k \neq j} (1-\alpha_j z) \prod_{k \neq j} (1-\alpha_j z^{-1}) z^{-1}
\]

(2.19)

where \( \alpha_j = e^{\lambda_j} \). Note that \( |\alpha_j| < 1 \) by virtue of the assumption that \( \text{re}(\lambda_j) < 0 \). Thus, the denominator is already factored as required, so that

\[
d(z) = \prod_{j=1}^{r} (1-\alpha_j z).
\]

(2.20)
The numerator must be factored to find $c(z)$. Standard procedures to find the zeroes of scalar polynomials can be used to achieve this factorization, as described by Hansen and Sargent [\textsuperscript{[1]}].

Thus we have that

$$z_t = \frac{c(L)}{d(L)} a_t \equiv C(L) a_t.$$  

Proceeding in a similar fashion as we did for the continuous time moving average representation, we can find a partial fraction representation for $C(L)$, namely

$$C(L) = \sum_{j=1}^{r} \frac{\gamma_j}{1-\alpha_j L}$$

where

$$\gamma_j = \lim_{z+\alpha_j} C(z) (1-\alpha_j z).$$

Recalling that $\alpha_j = e^{\lambda_j}$, equation (2.22) implies that

$$C_k = \sum_{j=1}^{r} \gamma_j e^{\lambda_j k}.$$  

Collecting and comparing the key results, we have that

$$p(\tau) = \sum_{j=1}^{r} \delta_j e^{\lambda_j \tau} \quad \tau \in [0, \infty).$$

$$C_k = \sum_{j=1}^{r} \gamma_j e^{\lambda_j k} \quad k = 0, 1, 2, \ldots.$$  

Equations (2.17) and (2.24) imply that $C_k$ will be (proportional to) a sampled version of $p(\tau)$ if and only if $\gamma_j/\delta_j = \gamma_1/\delta_1$ for all $j = 2, \ldots, r$. It can be shown directly by using (2.17) and (2.24) in (2.7) and (2.8) that this condition will not be met for
any r \geq 2. Thus, only if \( z(t) \) is a first-order autoregressive process does \( C_k \) turn out to be a sampled version of \( p(t) \).

Table 2 presents a numerical example that illustrates the preceding ideas. For the univariate process \((D^3+.6D^2+.4D+.2)z(t) = w(t)\), we have calculated \( p(t) \), \( f(t) \), \( c(L) \), \( d(L) \), \( B(L) = c(L)/d(L) \), \( \lambda_j \), \( \gamma_j \) for \( j = 1, 2, 3 \). In this example, we have that \( \gamma_j/\gamma_1 \neq \delta_j/\delta_1 \) for \( j \geq 2 \), so that the shapes of the moving averages in continuous and discrete time, \( p(t) \) and \( C_k \), respectively, are different. We also have that \( f(t) \neq 0 \) for some \( t \)'s greater than 1. In particular, notice that \( f(t) \) is larger in absolute value over most of the interval \([1,2]\) than it is over the interval \([0,1]\). The failure of \( f(t) \) to be concentrated on \([0,1]\) and the failure of \( B_k \) to resemble a sampled version of \( p(t) \) are both consequences of the fact that this is a third order autoregressive system in continuous time, rather than a first order one.

The preceding results and the example generalize readily to the case of a vector stochastic process \( z_t \). Matrix versions of (2.17) and (2.24) hold, where the \( \lambda_j \)'s are the zeroes of \( \text{det} \theta(s) \) and the \( \delta_j \)'s and \( \gamma_j \)'s are \((nxn)\) matrices given by (2.16) and (2.23).

Locally Unpredictable Processes and Linear Quadratic Models

The stochastic process \( z(t) \) in Table 2 is mean square differentiable, as evidenced by the fact that \( p(0) = 0 \). A stochastic process of the form (2.1) is \( j \) times mean square differentiable if \( p(0) = p'(0) = p''(0) = \ldots = p^{(j-1)}(0) = 0 \) (see Hansen and Sargent [ ] for a proof). Consequently, the process \((D^3 +
.6D^2 + .4D + .2)x(t) = w(t) can be verified to be twice (but not three times) mean square differentiable. It is the smoothness and proximity to zero near $\tau = 0$ of $p(\tau)$ that makes it difficult for $C_j$ to resemble a sampled version of $p(\tau)$, and that makes $a(t)$ a poor estimator of $\int_0^1 p(\tau)w(t-\tau)d\tau$.

Sims [1984] has argued that there is a class of economic variables that are best modeled as failing to be mean square differentiable. For these processes, $p(0) \neq 0$. Processes of the form (2.1) in which $p(0) \neq 0$ are said to be locally unpredictable because if $p(0) \neq 0$, then

$$\lim_{\delta \to 0} \frac{E_t(x(t+\delta) - E_t x(t+\delta))^2}{E_t(x(t+\delta) - x(t))^2} = 1.$$  

(2.25)

Here $E_t$ is the linear least squares projection operator, conditioned on $\{x(t-s), s > 0\}$. Now condition (2.25) can readily be shown to imply that

$$\lim_{\delta \to 0} \frac{E_t(x(t+\delta) - E_t x(t+\delta))^2}{E_t(x(t+\delta) - E(x(t+\delta)|x(t),x(t-\delta),x(t-2\delta),\ldots))^2} = 1$$

(2.26)

where $E$ is the linear least squares projection operator. In (2.26), $E_t x(t+\delta)$ is the linear least squares projection of $x(t+\delta)$ conditioned on $\{x(t-s), s > 0\}$, while $(E(x(t+\delta)|x(t),x(t-\delta),\ldots)$ is the projection of $x(t+\delta)$ on the discrete time sample $x(t), x(t-\delta), \ldots$. Condition (2.26) holds for any locally unpredictable process, and states that for small enough sampling interval $\delta$, the $\delta$-ahead projection error from the continuous time process is close in the mean square error sense to the $\delta$-ahead projection error.
from the $\delta$-discrete time data. Thus, when $p(0) \neq 0$, for small enough $\delta$, the innovation $a_t$ in the $\delta$-counterpart to (2.21) is arbitrarily close to $\int_{0}^{\delta} p(s)w(t-s)ds$ in the mean square sense.

Now suppose that $z(t)$ is given by (2.1), with $p(0) = 0$, so that $z(t)$ is mean square differentiable. Following Sims [1980], suppose that the economist is interested in studying the expectational variable $x^*(t)$ given by

$$(2.27) \quad x^*(t) = E\left[ \int_{0}^{\infty} e^{\rho s}z(t+s)ds \mid (z(t-T), x) > 0 \right]$$

where $\rho < 0$. Hansen and Sargent [1981] have shown that

$$(2.28) \quad x^*(t) = \left[ \frac{-P(D)+P(-\rho)}{D+\rho} \right]w(t) = G(D)w(t) = \int_{0}^{\infty} g(s)w(t-s)ds.$$

where $P(s) = \int_{0}^{\infty} e^{-\tau}p(\tau)d\tau$ is the Laplace transform of $p(\tau)$. Now if $G(s)$ is the Laplace transform of $g(\tau)$, with support $[0, \infty)$, the initial value theorem for Laplace transforms states that

$$g(0) = \lim_{s \to \infty} sG(s).$$

Using the initial value theorem together with (2.28), we find that

$$g(0) = \lim_{s \to \infty} s[\frac{-P(s)+P(-\rho)}{s+\rho}] = P(-\rho) \neq 0.$$

(We know that $P(-\rho) \neq 0$ because $P(s)$ is assumed to have no zeroes in the right half of the complex plane by the assumption that $p(s)$ is the kernel associated with a Wold representation for $z(t)$.) Therefore, even if $p(0) = 0$, $g(0) \neq 0$, so that the geometric expectational variable $x^*(t)$ fails to be mean square differenti-
able and therefore is locally unpredictable. For such expectational variables, (2.26) holds. Therefore, for such variables, for small enough sampling interval 6, the discrete time innovation a(t) corresponding to (2.21) is close to \( \int_0^\delta p(s)w(t-s)ds \) in the mean squared sense.

These results imply that for a variable \( x^*(t) \) and sufficiently small sampling interval 6, the situation is not as bad as is depicted by the example in Table 2. As Sims has pointed out, there are theories of consumption and asset pricing which imply that consumption or asset prices behave like \( x^*(t) \) and are governed by a version of (2.27). For example, with \( x^*(t) \) being consumption and \( z(t) \) income, (2.27) is a version of the permanent income theory. Alternatively, with \( x^*(t) \) being a stock price and \( z(t) \) being the dividend process, (2.27) is a simple version of an asset-pricing formula.

However, there is a wide class of generalized adjustment cost models discussed by Hansen and Sargent [1981] in which observable variables are such smoothed versions of \( x^*(t) \) that they are mean square continuous. In adjustment cost models, decisions are driven by convolutions of \( x^*(t) \), not by \( x^*(t) \) alone. For example, the stochastic Euler equation for a typical quadratic adjustment cost problem is

\[
(D-\rho)x(t) = E_t \left( \frac{1}{D+\rho} \right) z(t),
\]

where \( \rho > 0 \),

or

\[
(D-\rho)x(t) = x^*(t).
\]
Here \( k(t) \) is "capital." The solution for capital is then

\[
k(t) = \frac{1}{D - \rho} x^*(t)
\]

or

\[
k(t) = \left( \frac{1}{D - \rho} \right) \left[ \frac{-P(D) + P(-\rho)}{D + \rho} \right] w(t)
\]

where \( z(t) = \int_0^\infty \rho(s) w(t-s) \).

Let

\[
k(t) = \int_0^\infty h(\tau) w(t-\tau) d\tau
\]

where

\[
H(s) = \int_0^\infty e^{-\tau s} h(\tau) d\tau
\]

\[
H(s) = \left( \frac{1}{s - \rho} \right) \left[ \frac{-P(s) + P(-\rho)}{s + \rho} \right].
\]

Using the initial value theorem to calculate \( h(0) \) we have

\[
h(0) = \lim_{s \to \infty} sH(s) = 0.
\]

Thus, \( k(t) \) is mean square differentiable and so is not locally predictable. (The convolution integration required to transform \( x^*(t) \) to \( k(t) \) "smoothes" \( k(t) \) relative to \( x^*(t) \).)

More generally, the endogenous dynamics of adjustment cost models typically leads to mean square differentiable endogenous variables, provided that the agent is posited to be facing mean square differentiable forcing processes \( (z(t)) \). This means that for such models, the difficulties of interpretation that are illustrated in Table 2 cannot be eluded by appealing to an approximation based on the limit (2.26).
Table 2

An Example

\[ \psi(D) = 1 \]

\[ \theta(D) = 0.2 + 0.4D + 0.6D^2 + D^3 \]

\[ \lambda_j \text{ (zeroes of } \theta(s)) = 0.5424, 0.0288 \pm 0.6066 \]

\[ \delta_j \text{ in Partial Fraction Representation of } \psi(D)/\theta(D)^{-1} : \]

<table>
<thead>
<tr>
<th>( \delta_j )</th>
<th>( \delta_j/\delta_1 )</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
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<tr>
<td>2</td>
<td>-0.7915</td>
</tr>
<tr>
<td>3</td>
<td>-0.7915</td>
</tr>
</tbody>
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Zeroes of Spectral Factorization of Numerator Polynomial:

<table>
<thead>
<tr>
<th>Real Part of Zero</th>
<th>Imaginary Part of Zero</th>
<th>Modulus</th>
</tr>
</thead>
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<tr>
<td>1</td>
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</tr>
<tr>
<td>2</td>
<td>-0.4359</td>
<td>0</td>
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</tbody>
</table>

\[ \lambda_j \text{ in Partial Fraction Representation of } C(L) : \]

<table>
<thead>
<tr>
<th>Real (( \lambda_j ))</th>
<th>Imaginary (( \lambda_j ))</th>
<th>( \delta_j/\delta_1 )</th>
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<tr>
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<tr>
<td>3</td>
<td>-0.3992</td>
<td>-2.0310</td>
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</table>

Discrete Time Mixed Moving Average, Autoregressive Representation:

\[ d(L) = 1 - 2.1779L + 1.8722L^2 - 0.5485L^3 \]

\[ c(L) = 1 + 0.4800L + 0.0192L^2 \]
<table>
<thead>
<tr>
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<th>$p(t)$</th>
<th>$C_k$</th>
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<td>0.004900</td>
<td>0.004900</td>
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<tr>
<td>.300</td>
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<td>0.073563</td>
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<tr>
<td>.500</td>
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The preceding problems of interpretation are results of estimating vector autoregressions while foregoing the imposition of any explicit economic theory in estimation. These problems can be completely overcome if a sufficiently restrictive and reliable dynamic economic theory is available to be imposed during estimation. For example, Hansen and Sargent [1981,1982] have described how the function $p(\tau)$ can be identified and estimated from observations on discrete time data in the context of a wide class of linear rational expectations models. The basic idea is that the rich body of cross-equation restrictions that characterize dynamic linear rational expectations models can be used to identify a unique continuous time model from discrete time data.

If an estimate of $p(\tau)$ is available, then by using only discrete time data on $\{z_t\}$, it is even possible to recover an estimate of the one-step ahead prediction error that agents are making in continuous time. This is accomplished by treating the continuous time forecast error as a hidden variable whose covariances with the discrete time process $\{z_t\}$ are known. Thus, given estimates of $p(\tau)$, let us define the one-step ahead prediction error from continuous time data as $e_t^* = \int_0^1 p(\tau) w(t-\tau) \, d\tau$. Then it is straightforward to calculate the following second moments:

$$E[z_t z_{t-j}^T] = \int_0^\infty p(\tau+j)p(\tau)^T \, d\tau = \sum_{k=0}^\infty C_{k+j} \, W^{k+T}, \quad j > 0$$

$$E[e_t^* z_{t+j}^T] = \begin{cases} \int_0^1 p(\tau)p(\tau+j)^T \, d\tau & j > 0 \\ 0 & j < 0 \end{cases}$$
We can estimate the projection \( \sum_{j=-m_1}^{m_2} D_j z_{t-j} \) in the projection equation

\[
e^*_t = \sum_{j=-m_1}^{m_2} D_j z_{t-j} + u_t
\]

where \( u_t \) is orthogonal to \( z_{t-j} \) for all \( j = -m_1 \ldots, m_2 \). The \( D_j \)'s can be computed from the normal equations

\[
E e^*_t z_{t+k}^T = \sum_{j=-m_1}^{m_2} D_j E z_{t-j} z_{t+k}^T, \quad k = -m_2, \ldots, m_1.
\]

These calculations could be of use if one's aim were truly to extract and to interpret estimates of the forecast errors made by agents. In continuous time versions of various models, such as those of Lucas [1973] or Barro [1977], agents' forecasting errors are an important source of impulses, so that it is of interest to have this method for characterizing their stochastic properties and estimating them.
This is, after all, the construction used in Wold's decomposition theorem.

Representations of the moving average form (0.4) are not in general unique, once one relaxes the restriction in (0.1) that \( A_0^\Delta = -I \), which in turn implies that \( C_0^\Delta = I \). If this restriction is relaxed, then any representation generated by slipp­ing a \( UU^T \) in between \( C_j^\Delta \) and \( a(t-A_j) \) in (0.4), where \( U \) is a unitary matrix (\( UU^T=I \)), is also a fundamental moving average representa­tion. That is,

\[
z(t) = \sum_{j=0}^{\infty} (C_j^\Delta U) (U^T a(t-A_j))
\]

is also a fundamental moving average representation, since \( U^T a(t) \) spans the same linear space as \( a(t) \). In terms of such a representation, the decomposition of prediction error covariance becomes

\[
E(z(t)-\hat{E}_{t-j}z(t)) (z(t)-\hat{E}_{t-j}z(t))^T = \sum_{k=0}^{j-1} C_k^\Delta U V U^T C_j^\Delta T,
\]

which is altered by alternative choices of \( U \). Sims' choice of orthogonalization order amounts to a choice of \( U \).

An earlier version of this paper considered four classes of examples, the other two being nonlinearities and aggregation across agents. Due to length constraints, we decided to restrict this paper to the two classes of examples studied here.

Danny Quah has conveyed to us the viewpoint that implicit in the desire to match the \{w\} process of the economic
model (8) with the \{a\} process of the vector autoregression (1) must be a decision problem that concerns the data analyst. For example, on the basis of variance decompositions based on (0.7), the analyst might want to predict the consequences of "interventions" in the form of alterations in various diagonal elements of the innovation covariance matrix \( V \), interpreting these alterations, e.g., as changes in the predictability by agents of various economic process, e.g., the money supply.

3/ This assumption is made in the interests of providing the best possible chance that the process \( a(t) \) and \( w(t) \) described in the introduction match up. If \( \varepsilon_{1t} \) is a vector of dimension greater than \( n_1 \), then in general current and lagged values of \( (\varepsilon_{1t}, \varepsilon_{2t}) \) span a larger linear space than do current and lagged values of \( (y_t, x_{2t}) \).

6/ See Rozanov [ ] or Townsend [ ].

7/ It is interesting to note that although this system is one in which there are no strictly econometrically exogenous variables, or even any variables that are not Granger-caused by any others, its parameters are in principle identifiable. Identification is achieved through the cross-equation restrictions. Even when \( (w_{st}, w_{dt}) \) lie in the space spanned by the one-step ahead errors in predicting \( (q_t, p_t) \) from their own pasts, it is necessary to know the structural parameters of the model in order to deduce the former from the latter innovations.

8/ See Sims [ ] for a treatment of orthogonalization orders. Different "orthogonalization orders" in the sense of Sims
amount to different triangular choices of the orthogonal matrix $U$ that appears in footnote 2. If $U^T$ is chosen to be upper triangular, then the first component of $a_t$ corresponds to the first component of the new (basis) fundamental noise $U^T a_t$. On the other hand, if $U^T$ is chosen to be lower triangular, the last component of $a_t$ gets to go first in the Gram-Schmidt process that is used to create $U^T a_t$.

10/Practical methods for solving this equation for the case in which $P(s)$ is rational are discussed by Phillips [ ], Christiano [ ], and Hansen and Sargent [ ].

11/An alternative derivation of (2.7) uses operational calculus. Setting $L = e^{-D}$, express (2.5) as $a_t = V(e^{-D}) P(D) w(t) = f(D) w(t)$. Here the function $f(t)$ is the inverse Fourier transform of $F(iw)$, which is defined by

$$F(iw) = C(e^{-iw})^{-1} P(iw).$$

Equation (2.7) follows from the above equation by the convolution property of Fourier transforms.

12/For example, the function $p(t)$ will be continuous whenever $P(D)$ is rational, a common specification in applied work. The functions $p(t)$ and $f(t)$ are only defined up to an $L^2$ equivalence. Consequently, we can only impose continuity on one version of the continuous time moving average coefficients.
REFERENCES


