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Nonparametric identification of dynamic models with unobserved state variables

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cemmap working paper CWP13/08

Nonparametric Identification of Dynamic Models with Unobserved State Variables*

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Comments welcome

Abstract

We consider the identification of a Markov process $\{W_t, X_t^*\}$ for $t = 1, 2, \dots, T$ when only $\{W_t\}$ for $t = 1, 2, \dots, T$ is observed. In structural dynamic models, W_t denotes the sequence of choice variables and observed state variables of an optimizing agent, while X_t^* denotes the sequence of serially correlated unobserved state variables. The Markov setting allows the distribution of the unobserved state variable X_t^* to depend on W_{t-1} and X_{t-1}^* . We show that the joint distribution $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is identified from the observed distribution $f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}}$ under reasonable assumptions. Identification of $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is a crucial input in methodologies for estimating dynamic models based on the “conditional-choice-probability (CCP)” approach pioneered by Hotz and Miller.

1 Introduction

In this paper, we consider the identification of a Markov process $\{W_t, X_t^*\}$ for $t = 1, 2, \dots, T$ when only $\{W_t\}$ for $t = 1, 2, \dots, T$ is observed. The variable W_t describes the observed behavior and status of agent i at period t . X_t^* consists of latent variables, which are observed by the agent, but unobserved to the econometrician. The common interpretation of the latent variable X_t^* is an unobserved state variable at period t .

We show that the distribution $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is identified from the observed distribution $f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}}$ under reasonable assumptions. In most applications, W_t consists of

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two elements $W_t = (Y_t, M_t)$, where Y_t denotes the agent's action in period t , and M_t denotes the period- t observed state variable. X_t^* are persistent unobserved state variables (USV for short), which are observed by agents and affect their choice of Y_t , but are unobserved by the econometrician. In turn, the realization of the USV X_t^* can also be affected by Y_{t-1} or M_{t-1} , in addition to X_{t-1}^* . We begin by giving two motivating examples of well-known Markovian dynamic discrete-choice models which have been estimated in the existing literature.

Example 1: Rust (1987) In Rust's bus engine replacement model, Y_t is an indicator for whether Harold Zurcher (the bus depot manager) decides to replace the bus engine in week t . M_t is the accumulated mileage of the bus since the last engine replacement, in week t . Although Rust's original paper had no persistent unobserved state variable X_t^* , it is reasonable to extend the model to allow for them. For example, X_t^* could be Harold Zurcher's health, or weather or road conditions during week t .¹ ■

Example 2: Pakes (1986) Pakes estimates an optimal stopping model of the year-by-year renewal decision on European patents. In his model, the decision variable Y_t is an indicator for whether a patent is renewed in year t , and the unobserved state variable X_t^* is the profitability from the patent in year t , which is not observed by the econometrician. The observed state variable M_t could be other time-varying factors, such as the stock price or total sales of the firm holding the patent, which affect the renewal decision. ■

The main result in this paper concerns the identification of the joint density $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$. This implies that the conditional density $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ is also identified. Once this is known, it can be factorized into conditional and marginal distributions of economic interest. For Markov dynamic choice models (such as the two examples given above), an interesting factorization is

$$\begin{aligned} f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*} &= f_{Y_t, M_t, X_t^* | Y_{t-1}, M_{t-1}, X_{t-1}^*} \\ &= \underbrace{f_{Y_t | M_t, X_t^*}}_{\text{CCP}} \cdot \underbrace{f_{M_t, X_t^* | Y_{t-1}, M_{t-1}, X_{t-1}^*}}_{\text{state transition}}. \end{aligned} \quad (1)$$

The second term is the Markovian transition probabilities for the state variables (M_t, X_t^*) . The first term denotes the conditional choice probability for the agent's optimal choice in period t . Note that this setting accommodates quite general feedback in the unobserved

¹See Norets (2006), who likewise considers an example of the Rust (1987) model extended to accommodate persistent unobserved state variables.

state variable process from previous values W_{t-1}, X_{t-1}^* to X_t^* .

Once the CCP's and the state transitions are recovered, it is straightforward to use them as inputs in a CCP-based approach for estimating dynamic discrete-choice models. This approach was pioneered in Hotz and Miller (1993) and Hotz, Miller, Sanders, and Smith (1994), and subsequent methodological developments in this vein include Aguirregabiria and Mira (2002), Pesendorfer and Schmidt-Dengler (2003), Bajari, Benkard, and Levin (2008), Aguirregabiria and Mira (2007), Pakes, Ostrovsky, and Berry (2007), and Hong and Shum (2007).² Alternatively, it is possible to use our identification results for the CCP's and state transition densities as a "first-step" in an argument for identification of the per-period utility functions, in the spirit of Magnac and Thesmar (2002) and Bajari, Chernozhukov, Hong, and Nekipelov (2007), who considered the case of dynamic discrete-choice models without unobserved state variables.

A general criticism of these CCP-based methods is that they cannot accommodate unobservables which are persistent over time. However, there are some recent papers focusing on the CCP-based estimation of dynamic discrete-choice models, in the presence of the latent state variable X_t^* . Buchinsky, Hahn, and Hotz (2004) and Houde and Imai (2006) consider the case where X_t^* is time-invariant, corresponding to the case of unobserved heterogeneity, and discrete. Arcidiacono and Miller (2006) develop a CCP-based approach to estimate dynamic discrete models where X_t^* can vary over time according to an exogenous and discrete first-order Markov process.

Several recent papers have focused on the estimation of parametric dynamic models with unobserved state variables, using non-CCP-based approaches. Imai, Jain, and Ching (2005) and Norets (2006) consider Bayesian MCMC estimation. Fernandez-Villaverde and Rubio-Ramirez (2007) develop an efficient simulation procedure (based on particle filtering) for estimation these models via simulation.

While these papers have focused on estimation, our focus is on identification. Kasahara and Shimotsu (2007) considers the nonparametric identification of dynamic models when the latent variable X_t^* is time-invariant and discrete. In section 3.2 of their paper, Kasahara and Shimotsu prove the nonparametric identification of the Markov kernel $W_{t+1}|W_t, X^*$ in this setting, using six periods of data. In this paper, we build upon these results to the case

²Applications applying the CCP insights to dynamic settings have grown quickly in recent years, and include Collard-Wexler (2006), Ryan (2006), and Dunne, Klimer, Roberts, and Xu (2006). See the discussion in Pakes (2008, section 3) and Akerberg, Benkard, Berry, and Pakes (2007). All of these papers apply the CCP insight to dynamic games, which are more complex multi-agent generalizations of the single-agent dynamic setting consider in this paper.

where X_t^* is continuous, and can vary over time and evolve depending on (W_{t-1}, X_{t-1}^*) .

Finally, Cunha, Heckman, and Schennach (2006) apply the result of Hu and Schennach (2008) to show nonparametric identification of a nonlinear factor model containing $(W_t, W_t', W_t'', X_t^*)$, where the observed processes $\{W_t\}_{t=1}^T$, $\{W_t'\}_{t=1}^T$, $\{W_t''\}_{t=1}^T$ constitute noisy measurements of the latent process $\{X_t^*\}_{t=1}^T$, contaminated with random disturbances. In contrast, we consider a setting where (W_t, X_t^*) jointly evolves as a dynamic Markov process. We use observations of W_t in different periods t to identify the joint distribution of $(W_t, X_t^*, W_{t-1}, X_{t-1}^*)$. Thus, our model and identification strategy are different from theirs.

The paper is organized as follows. Section 2 contains our main identification result, which we prove for the case where X_t^* is continuous. We discuss the implications of the identification assumptions in the context of Rust's (1987) bus engine replacement model in Section 3. Section 4 discusses the nonparametric identification of DDC models given the results in section 2. We conclude in Section 5. The appendix includes the proof of the theorem, remarks, and a special case where the unobserved state variable X_t^* is discrete.

2 Nonparametric identification with unobserved state variables

Consider an i.i.d. sample of the dynamic process $\{(W_{t+1}, X_{t+1}^*), (W_t, X_t^*), \dots, (W_1, X_1^*)\}_i$ for agent $i \in \{1, 2, \dots, n\}$. The researcher observes an i.i.d. random sample of the dynamic process $\{W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}\}_i$ for many agents i . The variable W_t describes the observed behavior and status of the agent i at period t . The variable X_t^* stands for the unobserved state variable at period t . We assume that for each agent i , $\{(W_{t+1}, X_{t+1}^*), (W_t, X_t^*), \dots, (W_1, X_1^*)\}_i$ is an independent random draw from the identical distribution $f_{W_{t+1}, W_t, \dots, W_1, X_{t+1}^*, X_t^*, \dots, X_1^*}$. Let $\mathcal{W}_t \subseteq \mathbb{R}^d$ be the support of W_t and $\mathcal{X}_t^* \subseteq \mathbb{R}$ be the support of X_t^* . Define $\Omega_{<t} = \{W_{t-1}, \dots, W_1, X_{t-1}^*, \dots, X_1^*\}$.

We assume the dynamic process satisfies:

Assumption 1 (i) *First-order Markov*:

$$f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*, \Omega_{<t-1}} = f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}; \quad (2)$$

(ii) *Limited feedback*:

$$f_{W_t|W_{t-1}, X_t^*, X_{t-1}^*} = f_{W_t|W_{t-1}, X_t^*}. \quad (3)$$

Assumption 1(i) is just a first-order Markov assumption, which is assumed in most empirical applications of dynamic discrete-choice models, and holds for both the Pakes and Rust examples. Assumption 1(ii) is a “limited feedback” assumption, because it rules out direct feedback from the last period’s USV, X_{t-1}^* , on the current value of the observed component W_t . When $W_t = (Y_t, M_t)$, where Y_t denotes the agent’s action in period t , and M_t denotes the period- t observed state variable, Assumption 1 implies that:

$$\begin{aligned} f_{W_t|W_{t-1}, X_t^*, X_{t-1}^*} &= f_{Y_t, M_t|Y_{t-1}, M_{t-1}, X_t^*, X_{t-1}^*} \\ &= f_{Y_t|M_t, Y_{t-1}, M_{t-1}, X_t^*, X_{t-1}^*} \cdot f_{M_t|Y_{t-1}, M_{t-1}, X_t^*, X_{t-1}^*} \\ &= f_{Y_t|M_t, X_t^*, Y_{t-1}, M_{t-1}} \cdot f_{M_t|Y_{t-1}, M_{t-1}, X_t^*}. \end{aligned} \quad (4)$$

In the bottom line of the above display, the limited feedback assumption eliminates X_{t-1}^* as a conditioning variable in both terms. In most applications of Markov dynamic choice models, the first term (corresponding to the CCP) can be further simplified to $f_{Y_t|M_t, X_t^*}$, because the Markovian transition probabilities for the state variables M_t, X_t^* imply that the optimal policy function depends just on the current state variables, but not past realizations. (See Rust (1994, section 2) for a discussion of optimal policy functions in Markovian dynamic decision models.)

optimal policy function depends just on the current state variables, which are (M_t, X_t^*) . Hence, the above display shows that Assumption 1 imposes weaker restrictions on the first term than typical dynamic optimization models. Moreover, if we move outside the class of dynamic optimization models, Eq. (4) also shows that Assumption 1 does not rule out the dependence of Y_t on Y_{t-1} or M_{t-1} , which corresponds to some models of state dependence.³

In the second term of the above display, the limited feedback condition rules out direct feedback from last period’s unobserved state variable X_{t-1}^* to the current observed state variable X_t^* . However, it allows indirect effects via X_{t-1}^* ’s influence on Y_{t-1} or M_{t-1} . Indeed,

³These may include linear or nonlinear panel data models with lagged dependent variables, and serially correlated errors, cf. Arellano and Honore (2000). Arellano (2003, chs. 7–8) considers linear panel models with lagged dependent variables and persistent unobservables, which is also related to our framework.

most empirical applications of dynamic optimization models with unobserved state variables satisfy the Markov and limited feedback conditions above. Examples of models in the industrial organization setting satisfying these conditions include Pakes (1986), Akerberg (2003), Erdem, Imai, and Keane (2003), Crawford and Shum (2005), Das, Roberts, and Tybout (2007), Xu (2007), and Hendel and Nevo (2007). Finally, note that when X_t^* is time invariant, so that $X_t^* = X_{t-1}^*$, the limited feedback assumption is trivial.

Our goal is to identify the density

$$f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}.$$

Since W_{t+1} is usually a vector and X_t^* is a scalar, we first reduce the dimensionality of W_{t+1} by defining

$$V_{t+1} \equiv g(W_{t+1})$$

where the function $g : \mathbb{R}^d \rightarrow \mathbb{R}$ is known. (When W_{t+1} is a scalar, we may just let $g(w) = w$.) Another advantage of introducing V_{t+1} is that the identification still holds with a discrete X_t^* if we let $g : \mathcal{W}_{t+1} \rightarrow \mathcal{X}_t^*$. The restrictions imposed later on the function g guarantee that the scalar random variable V_{t+1} still contains enough information to identify X_t^* . Similarly, we reduce the dimensionality of W_{t-2} by defining

$$Z_{t-2} \equiv q(W_{t-2}),$$

with a known function $q : \mathbb{R}^d \rightarrow \mathbb{R}$. When X_t^* is discrete, we may let $q : \mathcal{W}_{t-2} \rightarrow \mathcal{X}_{t-2}^*$. We introduce the function q only for the reason of avoiding technical complications. As shown later, we may just let $q(w) = w$ by using the generalized inverse of an operator.

The identification argument consists of four steps. The discussion in this section omits the derivation of some equations. A complete proof, including all derivations, is given in the Appendix.

Step 1: Identification of $f_{V_{t+1}|W_t, X_t^*}$. The most substantial step of the argument is the first step, which demonstrate the identification of $f_{V_{t+1}|W_t, X_t^*}$. Consider the joint density of $\{V_{t+1}, W_t, W_{t-1}, Z_{t-2}\}$. One can show that assumption 1 implies that for any

$$(x, w_t, w_{t-1}, z) \in g(\mathcal{W}_{t+1}) \times \mathcal{W}_t \times \mathcal{W}_{t-1} \times q(\mathcal{W}_{t-2})$$

$$\begin{aligned} & f_{V_{t+1}, W_t | W_{t-1}, Z_{t-2}}(x, w_t | w_{t-1}, z) \\ = & \int f_{V_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) f_{X_t^* | W_{t-1}, Z_{t-2}}(x_t^* | w_{t-1}, z) dx_t^*. \end{aligned} \quad (5)$$

The density on the left hand side is observed in the data. Let $\mathcal{L}^p(\mathcal{X})$, $1 \leq p < \infty$ stand for the space of function $h(\cdot)$ with $\int_{\mathcal{X}} |h(x)|^p dx < \infty$, and let $\mathcal{L}^\infty(\mathcal{X})$ denote the space of function $h(\cdot)$ with $\sup_{x \in \mathcal{X}} |h(x)| < \infty$. We let $p = 2$ when an inner product is introduced later. For any $1 \leq p \leq \infty$, we define the integral operator $L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} : \mathcal{L}^p(q(\mathcal{W}_{t-2})) \rightarrow \mathcal{L}^p(g(\mathcal{W}_{t+1}))$ for any given $(w_t, w_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1}$ and any $h \in \mathcal{L}^p(q(\mathcal{W}_{t-2}))$,

$$(L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} h)(x) = \int f_{V_{t+1}, W_t | W_{t-1}, Z_{t-2}}(x, w_t | w_{t-1}, z) h(z) dz.$$

Notice that we treat (w_t, w_{t-1}) as fixed and $L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}}$ is a mapping from $\mathcal{L}^p(q(\mathcal{W}_{t-2}))$ to $\mathcal{L}^p(g(\mathcal{W}_{t+1}))$.

For any given $w_t \in \mathcal{W}_t$, we also define the operator corresponding to the unobserved density $f_{V_{t+1} | W_t, X_t^*}$, i.e., $L_{V_{t+1} | w_t, X_t^*} : \mathcal{L}^p(\mathcal{X}_t^*) \rightarrow \mathcal{L}^p(g(\mathcal{W}_{t+1}))$, as follows:

$$(L_{V_{t+1} | w_t, X_t^*} h)(x) = \int f_{V_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) h(x_t^*) dx_t^*.$$

As shown in Hu and Schennach (2008), the identification of an operator, e.g., $L_{V_{t+1} | w_t, X_t^*}$, is equivalent to that of its corresponding density, e.g., $f_{V_{t+1} | W_t, X_t^*}$. Define for any given $w_t \in \mathcal{W}_t$

$$\mathcal{A}(w_t) = \{w_{t-1} \in \mathcal{W}_{t-1} : L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} \text{ is one-to-one}\}.$$

Identification of $L_{V_{t+1} | w_t, X_t^*}$ from the observed $L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}}$ requires

Assumption 2 for any $w_t \in \mathcal{W}_t$,

- (i) $L_{V_{t+1} | w_t, X_t^*}$ is one-to-one ;
- (ii) $\Pr \{\mathcal{A}(w_t)\} > 0$.

Assumption 2(i) implies that the function g reduces the dimension of W_t but $V_{t+1} = g(W_{t+1})$ still contains enough information on X_t^* . A sufficient condition for assumption 2(i) is that $L_{V_{t+1} | w_t, X_t^*} h = 0$ implies $h = 0$. A detailed discussion on one-to-one operators can be found

in Carrasco, Florens, and Renault (2005) and Hu and Schennach (2008). In the case where W_{t+1} are discrete and X_t^* is continuous, the assumption 2(i) fails. Notice that assumption 2(ii) is imposed on the observables.

REMARK: The one-to-one assumptions on $L_{V_{t+1}|w_t, X_t^*}$ and $L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}}$ rule out cases where X_t^* has a continuous support, but W_{t+1} has only discrete components. Hence, dynamic discrete-choice models with a continuous unobserved state variable X_t^* , but only discrete observed state variables M_t , fail this assumption, and may be nonparametrically underidentified without further assumptions. Moreover, models where the W_t and X_t^* processes evolve independently will also fail the one-to-one assumption. ■

REMARK: When we just use W_{t-2} instead of Z_{t-2} , it is possible that the corresponding operator $L_{V_{t+1}, w_t|w_{t-1}, W_{t-2}}$ may be surjective. In this case, there are extra instruments for X_t^* , and Assumption 2(ii) may be replaced by the condition that

$$\Pr \left\{ w_{t-1} : L_{V_{t+1}, w_t|w_{t-1}, W_{t-2}} L_{V_{t+1}, w_t|w_{t-1}, W_{t-2}}^* \text{ is one-to-one} \right\} > 0.$$

where L^* denotes an adjoint operator.⁴ We would then need to use the generalized inverse of $L_{V_{t+1}, w_t|w_{t-1}, W_{t-2}}$ instead of the inverse of $L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}}$. By using $Z_{t-2} = q(W_{t-2})$ and reducing the dimensionality of W_{t-2} to that of X_t^* , we avoid the technical complications of stating assumptions in terms of inner products or adjoint operators. ■

Inspired by the identification strategies in Carroll, Chen, and Hu (2008), Hu (2007), and Hu and Schennach (2008), we assume, in assumption 3 below, that for any given $w_t \in \mathcal{W}_t$ there exists $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1} \times \mathcal{W}_{t-1}$ such that $\bar{w}_t \neq w_t$, $\bar{w}_{t-1} \neq w_{t-1}$, and that $L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}}$, $L_{V_{t+1}, \bar{w}_t|w_{t-1}, Z_{t-2}}$, $L_{V_{t+1}, w_t|\bar{w}_{t-1}, Z_{t-2}}$ and $L_{V_{t+1}, \bar{w}_t|\bar{w}_{t-1}, Z_{t-2}}$ are all one-to-one mappings. One can show that equation 5 implies

$$\begin{aligned} & L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}} L_{V_{t+1}, \bar{w}_t|w_{t-1}, Z_{t-2}}^{-1} \left(L_{V_{t+1}, w_t|\bar{w}_{t-1}, Z_{t-2}} L_{V_{t+1}, \bar{w}_t|\bar{w}_{t-1}, Z_{t-2}}^{-1} \right)^{-1} \\ \equiv & L_{V_{t+1}|w_t, X_t^*} D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} L_{V_{t+1}|w_t, X_t^*}^{-1} \end{aligned} \quad (6)$$

where $D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} : \mathcal{L}^P(\mathcal{X}_t^*) \rightarrow \mathcal{L}^P(\mathcal{X}_t^*)$

$$(D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} g)(x_t^*) = k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) g(x_t^*),$$

⁴Let $L_{x,z}^* : \mathcal{L}^2(\mathcal{Z}) \rightarrow \mathcal{L}^2(\mathcal{X})$ denotes the adjoint operator of operator $L_{x,z} : \mathcal{L}^2(\mathcal{X}) \rightarrow \mathcal{L}^2(\mathcal{Z})$ such that $\langle L_{x,z} \varphi, \phi \rangle_{\mathcal{Z}} = \langle \varphi, L_{x,z}^* \phi \rangle_{\mathcal{X}}$, where the inner product is defined as $\langle \varphi, \phi \rangle = \int \varphi(t) \phi(t) dt$.

$$k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) = \frac{f_{W_t|W_{t-1}, X_t^*}(w_t|w_{t-1}, x_t^*) f_{W_t|W_{t-1}, X_t^*}(\bar{w}_t|\bar{w}_{t-1}, x_t^*)}{f_{W_t|W_{t-1}, X_t^*}(\bar{w}_t|w_{t-1}, x_t^*) f_{W_t|W_{t-1}, X_t^*}(w_t|\bar{w}_{t-1}, x_t^*)}.$$

Notice that the operator $D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*}$ is a "diagonal" or multiplication operator with a given $(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1})$. This equation implies that the observed operator on the left hand side, which is a mapping from $\mathcal{L}^p(\mathcal{X}_t^*) \rightarrow \mathcal{L}^p(\mathcal{X}_t^*)$, has an eigenvalue-eigenfunction decomposition. The eigenfunctions are $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$, which are normalized by $\int f_{V_{t+1}|W_t, X_t^*} dx_{t+1} = 1$. Notice that the eigenfunction in $L_{V_{t+1}|w_t, X_t^*}$ does not depend on $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1})$, while the eigenvalue in $D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*}$ may be different for a different $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1})$. The identification of $f_{V_{t+1}|W_t, X_t^*}$ then relies on the uniqueness of such a decomposition.

Formally, define a set $\mathcal{B}(w_t)$ for a given w_t such that any $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1}) \in \mathcal{B}(w_t)$ satisfies the following conditions:

1. $\bar{w}_t \in \mathcal{W}_t$, $w_{t-1} \in \mathcal{A}(\bar{w}_t)$, $\bar{w}_{t-1} \in \mathcal{A}(w_t) \cap \mathcal{A}(\bar{w}_t)$, $\bar{w}_t \neq w_t$, and $\bar{w}_{t-1} \neq w_{t-1}$;
2. $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) < \infty$ for all $x_t^* \in \mathcal{X}_t^*$.

Essentially, for a given $w_t \in \mathcal{W}_t$, the set $\mathcal{B}(w_t)$ contains triples of points $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1} \times \mathcal{W}_{t-1}$ such that $\bar{w}_t \neq w_t$, $\bar{w}_{t-1} \neq w_{t-1}$, and that $L_{V_{t+1}, \bar{w}_t|w_{t-1}, Z_{t-2}}$, $L_{V_{t+1}, w_t|\bar{w}_{t-1}, Z_{t-2}}$ and $L_{V_{t+1}, \bar{w}_t|\bar{w}_{t-1}, Z_{t-2}}$ are all one-to-one mappings. Notice that $L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}}$ in equation 6 is not required to be one-to-one. Furthermore, at these points, the eigenvalues $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*)$ are bounded away from $+\infty$. The boundedness of the eigenvalues allows us to use the results on the spectral decomposition of bounded linear operators in Dunford and Schwartz (1971).

A sufficient condition for $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) < \infty$ for all $x_t^* \in \mathcal{X}_t^*$ is that, for all $(w_t, w_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1}$, there exist functions $L(w_t, w_{t-1})$ and $U(w_t, w_{t-1})$ such that for all $x_t^* \in \mathcal{X}_t^*$

$$0 < L(w_t, w_{t-1}) \leq f_{W_t|W_{t-1}, X_t^*}(w_t|w_{t-1}, x_t^*) \leq U(w_t, w_{t-1}) < \infty. \quad (7)$$

The existence and uniqueness of the decomposition in equation 6 requires

Assumption 3 for any given $w_t \in \mathcal{W}_t$,

(i) $\Pr \{\mathcal{B}(w_t)\} > 0$;

(ii) for any $\hat{x}_t^* \neq \tilde{x}_t^* \in \mathcal{X}_t^*$, there exists $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1}) \in \mathcal{B}(w_t)$ such that $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, \hat{x}_t^*) \neq k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, \tilde{x}_t^*)$.

Part (i) of this assumption guarantees that for any given $w_t \in \mathcal{W}_t$, there exists more than one $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1}) \in \mathcal{B}(w_t)$ such that $\bar{w}_t \neq w_t$, $\bar{w}_{t-1} \neq w_{t-1}$, and that $L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}}$, $L_{V_{t+1}, \bar{w}_t | w_{t-1}, Z_{t-2}}$, $L_{V_{t+1}, w_t | \bar{w}_{t-1}, Z_{t-2}}$ and $L_{V_{t+1}, \bar{w}_t | \bar{w}_{t-1}, Z_{t-2}}$ are all one-to-one. This validates taking inverses of the operators in equation 6.

Part (ii) implies that all the eigenvalues are finite and distinctive for some $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1})$ in equation 6. Notice that $\ln k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*)$ can be treated as a second order difference of $\ln f_{W_t | W_{t-1}, X_t^*}$ with respect to w_t and w_{t-1} . Therefore, a sufficient condition for part (ii) is that for any $x_t^* \in \mathcal{X}_t^*$ and $w_t \in \mathcal{W}_t$, there exists $w_{t-1} \in \mathcal{W}_{t-1}$ such that

$$\frac{\partial^3}{\partial w_t \partial w_{t-1} \partial x_t^*} \ln f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) \neq 0. \quad (8)$$

REMARK: Since condition 2 in the definition of $\mathcal{B}(w_t)$ must be satisfied for all $w_t \in \mathcal{W}_t$, it will be violated if $f_{W_t | W_{t-1}, X_t^*}$ is identically zero for all X_t^* , and all W_{t-1} . However, in practice, most empirical applications of dynamic models avoid this possibility by including i.i.d. shocks which smooth out the CCP's and state transitions in order to avoid zeros, which are inconvenient from a computational point of view. In section 3 and Appendix B, we present examples of $f_{W_t | W_{t-1}, X_t^*}$ which satisfy assumption 3. ■

REMARK: Given the forgoing discussion, assumptions 2 and 3 may be replaced by the following sufficient conditions:

1. For any $w_t \in \mathcal{W}_t$ and $w_{t-1} \in \mathcal{W}_{t-1}$, $L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}}$ and $L_{V_{t+1} | w_t, X_t^*}$ are one-to-one ;
2. For any $w_t \in \mathcal{W}_t$ and $w_{t-1} \in \mathcal{W}_{t-1}$, there exist functions $L(w_t, w_{t-1})$ and $U(w_t, w_{t-1})$ such that the density $f_{W_t | W_{t-1}, X_t^*}$ satisfies Eq. (7) for all $x_t^* \in \mathcal{X}_t^*$;
3. For any $w_t \in \mathcal{W}_t$ and $x_t^* \in \mathcal{X}_t^*$, there exists $w_{t-1} \in \mathcal{W}_{t-1}$ such that the density $f_{W_t | W_{t-1}, X_t^*}$ satisfies Eq. (8). ■

Without further assumptions, an eigenfunction $f_{V_{t+1} | W_t, X_t^*}(\cdot | w_t, x_t^*)$ for a given w_t is only identified up to the value of the index x_t^* . Since the value of x_t^* is not observed anywhere, there is no difference between x_t^* and its monotone transformation. We may make the following assumption:

Assumption 4 for any given $w_t \in \mathcal{W}_t$,

- (i) There exist a known functional G such that $G \left[f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*) \right]$ is monotonic in x_t^* ;
- (ii) Without loss of generality, we normalize x_t^* as $x_t^* = G \left[f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*) \right]$.

This assumption pins down the value of x_t^* identified from each eigenfunction $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$. This normalization allows x_t^* to depend on w_t , which accommodates the fact that X_t^* may be correlated with W_t . Assumption 4 also provides an approach to estimate the model following the identification procedure. As shown in Hu and Schennach (2008), such a normalization may be very flexible. Because this procedure holds for all $w_t \in \mathcal{W}_t$, the density $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$ and operator $L_{V_{t+1}|w_t, X_t^*}$ are nonparametrically identified.

REMARK: The functional G may map the density f to its mean, mode, median or other quantile, for example, $G[f] = \int x f(x) dx$ or $G[f] = \inf \left\{ \tilde{x}^* : \int_{-\infty}^{\tilde{x}^*} f(x) dx \geq \tau \right\}$. Moreover, the functional G may depend on w_t . When G corresponds to a quantile, Matzkin (2003) suggests that for a fixed w_t one may have $V_{t+1} = h_{w_t}(X_t^*, \varepsilon)$, where ε is independent of X_t^* and has a standard uniform distribution. The function h_{w_t} can be interpreted as the inverse of the cdf $F_{V_{t+1}|W_t=w_t, X_t^*}$. That implies the τ -th quantile of $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$ is $h_{w_t}(x_t^*, \tau)$. Assumption 4 then requires that $h_{w_t}(x_t^*, \tau)$ is monotonic in x_t^* for a known τ . We may then normalize x_t^* as $x_t^* = h_{w_t}(x_t^*, \tau)$ without loss of generality. ■

Step 2: Identification of $\mathbf{f}_{W_{t+1}|W_t, X_t^*}$. In order to identify the density $f_{W_{t+1}|W_t, X_t^*}$, we define the following operators $L_{W_{t+1}, w_t|w_{t-1}, Z_{t-2}} : \mathcal{L}^p(q(\mathcal{W}_{t-2})) \rightarrow \mathcal{L}^p(\mathcal{W}_{t+1})$ and $L_{W_{t+1}|w_t, X_t^*} : \mathcal{L}^p(\mathcal{X}_t^*) \rightarrow \mathcal{L}^p(\mathcal{W}_{t+1})$ as

$$\begin{aligned} (L_{W_{t+1}, w_t|w_{t-1}, Z_{t-2}} h)(x) &= \int f_{W_{t+1}, W_t|W_{t-1}, Z_{t-2}}(x, w_t|w_{t-1}, z) h(z) dz, \\ (L_{W_{t+1}|w_t, X_t^*} h)(x) &= \int f_{W_{t+1}|W_t, X_t^*}(x|w_t, x_t^*) h(x_t^*) dx_t^*. \end{aligned}$$

One may show that $L_{W_{t+1}|w_t, X_t^*}$ is identified from

$$L_{W_{t+1}|w_t, X_t^*} = L_{W_{t+1}, w_t|w_{t-1}, Z_{t-2}} L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}}^{-1} L_{V_{t+1}|w_t, X_t^*}. \quad (9)$$

for any $w_t \in \mathcal{W}_t$. To see this, note that

$$\begin{aligned} L_{W_{t+1}, w_t|w_{t-1}, Z_{t-2}} &= L_{W_{t+1}|W_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}}^{-1} \\ L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}} &= L_{V_{t+1}|W_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}}^{-1} \end{aligned} \quad (10)$$

which leads to equation 9. Consequently, the density $f_{W_{t+1}|W_t, X_t^*}$ is identified.

REMARK: In the time-invariant case where $X_t^* = X^*$, $\forall t$, the conditional density $f_{W_{t+1}|W_t, X^*}$ is the main object of interest, and is enough to permit CCP-based estimation of dynamic discrete-choice models. However, when X_t^* varies over time, knowing $f_{W_{t+1}|W_t, X^*}$ is not enough to permit CCP-based estimation. ■

Step 3: Identification of $\mathbf{f}_{W_t, X_t^*, W_{t-1}, Z_{t-2}}$. As an intermediate step, we show that the density $f_{W_t, X_t^*|W_{t-1}, Z_{t-2}}$ is also identified. With $L_{V_{t+1}|w_t, X_t^*}$ identified in the first step, the density $f_{W_t, X_t^*, W_{t-1}, W_{t-2}}$ may also be identified as

$$f_{W_t=w_t, X_t^*, W_{t-1}, W_{t-2}} = L_{V_{t+1}|w_t, X_t^*}^{-1} f_{V_{t+1}, W_t=w_t, W_{t-1}, W_{t-2}}.$$

for any given $w_t \in \mathcal{W}_t$. Given the known mapping from W_{t-2} to Z_{t-2} , the identification of $f_{W_t, X_t^*, W_{t-1}, W_{t-2}}$ implies that of $f_{W_t, X_t^*|W_{t-1}, Z_{t-2}}$. Moreover, because the density of W_{t-1}, Z_{t-2} is identified from the data, the conditional density $f_{W_t, X_t^*|W_{t-1}, W_{t-2}}$ is also identified.

Step 4: Identification of $\mathbf{f}_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$. Finally, we show that the density of interest $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is identified. Assumption 1(i) implies

$$f_{W_t, X_t^*|W_{t-1}, Z_{t-2}} = \int f_{W_t, X_t^*|W_{t-1}, X_{t-1}^*} f_{X_{t-1}^*|W_{t-1}, Z_{t-2}} dx_{t-1}^*. \quad (11)$$

The density $f_{W_t, X_t^*|W_{t-1}, Z_{t-2}}$ on the left hand side of equation 11 is identified in the previous step. Thus far, we have only used the four observations $\{W_{t+1}, W_t, W_{t-1}, W_{t-2}\}$. In order to identify the density $f_{X_{t-1}^*|W_{t-1}, Z_{t-2}}$ on the right hand side of equation 11, we use one more period of the data W_{t-3} . Replacing t by $t-1$ in the previous three steps implies that the density of $\{W_t, W_{t-1}, W_{t-2}, W_{t-3}\}$ identifies $f_{W_{t-1}, X_{t-1}^*|W_{t-2}, Z_{t-3}}$ for $Z_{t-3} = q(W_{t-3})$. In turn, we can identify the density $f_{X_{t-1}^*|W_{t-1}, W_{t-2}}$ from $f_{W_{t-1}, X_{t-1}^*, W_{t-2}, Z_{t-3}}$ as

$$f_{X_{t-1}^*, W_{t-1}, W_{t-2}} = \int f_{W_{t-1}, X_{t-1}^*|W_{t-2}, Z_{t-3}} f_{W_{t-2}, Z_{t-3}} dz_{t-3}.$$

Given the known mapping q from W_{t-2} to Z_{t-2} , we can identify $f_{X_{t-1}^*|W_{t-1}, Z_{t-2}}$.

Now that the densities $f_{W_t, X_t^*|W_{t-1}, Z_{t-2}}$ and $f_{X_{t-1}^*|W_{t-1}, Z_{t-2}}$ in equation 11 have been identified, the density of interest $f_{W_t, X_t^*|W_{t-1}, X_{t-1}^*}$ may be identified under the following assumption

tion. Define $L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} : \mathcal{L}^p(\mathcal{X}_{t-1}^*) \rightarrow \mathcal{L}^p(\mathcal{X}_t^*)$,

$$\left(L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} h \right) (x_t^*) = \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}(w_t, x_t^* | w_{t-1}, x_{t-1}^*) h(x_{t-1}^*) dx_{t-1}^*.$$

We assume

Assumption 5 *for any $w_{t-1} \in \mathcal{W}_{t-1}$ and $w_t \in \mathcal{W}_t$, $L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*}$ is one-to-one.*

When X_t^* is discrete, this assumption requires that the support of X_t^* is time-invariant, i.e., $\mathcal{X}_{t-1}^* = \mathcal{X}_t^*$. Under assumption 5, one can show that the density $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ is identified from equation 11. We summarize the main identification results as follows:

Theorem 1 *Under the assumptions 1, 2, 3, 4, and 5, the density $f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}}$ uniquely determines the density $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$.*

Proof. see appendix. ■

In the case where X_t^* is discrete, the whole identification procedure still holds, which is presented in detail in the appendix. This result implies that the whole dynamic process $\{W_t, X_t^*\}$ is identified even if we only observe $\{W_t\}$. Moreover, the density f_{W_{t-1}, X_{t-1}^*} is identified from $f_{W_{t-1}, X_{t-1}^*, W_{t-2}, X_{t-2}^*}$, so that the unconditional density $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is also identified from $f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}}$.

3 Comments on Assumptions in Specific Example: Rust's (1987) Engine Replacement Model

Because some of the assumptions that we made for our identification argument are quite abstract, in this section we discuss these assumptions in the context of a version of Rust's (1987) bus-engine replacement model, augmented to allow for persistent unobserved state variables. As we remarked before, in this model, $W_t = (Y_t, M_t)$, where Y_t is the indicator that the bus engine was replaced in week t , and M_t is the mileage since the last engine replacement.

Because the stylized model we consider here is fully parametric, it may be identified without needing our identification results. However, what we focus on here is not the identifiabil-

ity of this model, but rather whether data generated from this model would allow us to nonparametrically identify the Markov kernel $W_t, X_t^* | W_{t-1}, X_{t-1}^*$.

We introduce two specifications of the model, which differ in how the unobserved state variable X_t^* enters. In both specifications, we assume that X_t^* evolves as a first-order Markov process, which can depend on past realizations of Y_t and M_t . For technical reasons (as will be clear below), we will restrict X_t^* to have a bounded support: for $[L, U]$ such that $-\infty < L < U < +\infty$,

$$X_t^* = \begin{cases} 0.5X_{t-1}^* + 0.3\psi(M_{t-1}) + 0.2\nu_t & \text{if } Y_{t-1} = 0 \\ 0.8X_{t-1}^* + 0.2\nu_t & \text{if } Y_{t-1} = 1 \end{cases} \quad (12)$$

with

$$\psi(M_{t-1}) = L + (U - L) \frac{\exp(M_{t-1}) - 1}{\exp(M_{t-1}) + 1},$$

where ν_t is a truncated standard normal shock over the interval $[L, U]$, distributed independently over weeks t , and the $\psi(\cdot)$ function maps mileage $M_{t-1} \in [0, +\infty)$ into $[L, U]$. We also assume that the support of the initial value X_0^* is $[L, U]$, which guarantees that the support of X_t^* is $[L, U]$ for all t . Hence, $X_t^* | X_{t-1}^*, Y_{t-1}, M_{t-1}$ is distributed with density determined by $f_{\nu_t}(\cdot)$. Furthermore, we assume that the characteristic function of ν_t satisfies that $\phi_{\nu_t}(s) \neq 0$ for any real s , which simply requires $L + U \neq 0$. This restriction on ϕ_{ν_t} guarantees that the operator corresponding to the density $f_{X_t^* | X_{t-1}^*, Y_{t-1}, M_{t-1}}$ is injective.

Let $S_t \equiv (M_t, X_t^*)$ denote the persistent state variables in this model. Following Rust (1987), we assume that the single-period utility from each choice is additive in a function of the state variables S_t , and a choice-specific non-persistent preference shock:

$$u_t = \begin{cases} u_0(S_t) + \epsilon_{0t} & \text{if } Y_t = 0 \\ u_1(S_t) + \epsilon_{1t} & \text{if } Y_t = 1 \end{cases}$$

where ϵ_{0t} and ϵ_{1t} are i.i.d. Type I Extreme Value shocks, which are also independent over time, and also independent of the state variables S_t .

Specification A In this specification, the choice-specific utility functions are:

$$\begin{aligned} u_0(S_t) &= -c(M_t) + X_t^* \\ u_1(S_t) &= -RC. \end{aligned} \tag{13}$$

In the above, $c(M_t)$ denotes the maintenance cost function, which is increasing in mileage M_t , and $0 < RC < +\infty$ denotes the cost of replacing the engine. We also assume that the maintenance cost function $c(\cdot)$ is bounded below and above:

$$c(0) = 0; \quad \lim_{M \rightarrow +\infty} c(M) = \bar{c} < +\infty.$$

Mileage evolves as:

$$M_{t+1} = \begin{cases} M_t + \eta_{t+1} & \text{if } Y_t = 0 \\ \eta_{t+1} & \text{if } Y_t = 1 \end{cases} \tag{14}$$

where the incremental mileage $\eta_{t+1} > 0$ is uniformly distributed $U[0, 1]$,⁵ independent across weeks, and independent of $(X_t^*, \epsilon_{0t}, \epsilon_{1t})$. ■

Specification B In this second specification, the agent's per-period utility functions are given by:

$$\begin{aligned} u_0(S_t) &= -c(M_t) \\ u_1(S_t) &= -RC. \end{aligned} \tag{15}$$

with the same assumptions on RC and $c(\cdot)$ as in Specification A. Mileage evolves as:

$$M_{t+1} = \begin{cases} M_t + \eta_{t+1} \cdot \exp(X_{t+1}^*) & \text{if } Y_t = 0 \\ \eta_{t+1} \cdot \exp(X_{t+1}^*) & \text{if } Y_t = 1. \end{cases} \tag{16}$$

Here, the incremental mileage $\eta_{t+1} \cdot \exp(X_{t+1}^*)$ is distributed as a mixture of a uniform and truncated lognormal distribution. ■

Finally, for the dimension-reducing mappings $g(\cdot)$ and $q(\cdot)$ introduced at the beginning of

⁵For this to be reasonable, assume that mileage is measured in units of 10,000 miles.

Section 3, we use:

$$\begin{aligned} V_{t+1} &= g(W_{t+1}) = M_{t+1} \\ Z_{t-2} &= q(W_{t-2}) = M_{t-2}. \end{aligned}$$

That is, the $g(\cdot)$ and $q(\cdot)$ mappings pick out the continuous component of W_t , which is just the mileage M_t .

The main difference between the two specifications is that in Specification A, the unobserved state variable X_t^* affects utilities directly (and therefore the CCP's), but not the mileage process. In Specification B, X_t^* directly affects the evolution of mileage, but not the agent's utilities. We will see that these two specifications differ in how well they satisfy the assumptions of the identification proof.

Given the assumptions so far, the conditional choice probabilities take the multinomial logit form (for $Y_t = 0, 1$):

$$P(Y_t|S_t) = \frac{\exp(V_{Y_t}(S_t))}{\sum_{y=0}^1 \exp(V_y(S_t))}$$

where $V_y(S_t)$ is the choice-specific value function in period t , which is defined recursively by

$$V_y(S_t) = u_y(S_t) + \beta E \left[\log \left\{ \sum_{y'=0}^1 \exp(V_{y'}(S_{t+1})) \right\} \mid Y_t = y, S_t \right].$$

Assumption 1 has already been discussed in much detail thus far, and it is satisfied for both specifications. We now comment on each remaining assumption in turn.

Assumption 2 Assumption 2 contains two “injectivity” (or one-to-one) assumptions, and we consider both in some detail. The first requirement is that: for all $w_t \in \mathcal{W}_t$, there exists w_{t-1} such that $L_{M_{t+1}, w_t | w_{t-1}, M_{t-2}}$ is one-to-one. (Note that we have substituted M_{t+1} for $g(W_{t+1})$, and M_{t-2} for $q(W_{t-2})$.)

Consider Specification A, and consider w_t such that $Y_t = 1$ (so that the engine is replaced in period t). In this case, $M_{t+1} | Y_t = 1$ is uniformly distributed on $[0, 1]$, and does not depend stochastically on either w_{t-1} or M_{t-2} . Hence, the one-to-one assumption fails.

Now consider Specification B, using the same w_t such that $Y_t = 1$. Because X_t^* directly enters the mileage process, the distribution of M_{t+1} depends on X_{t+1}^* . Similarly, M_{t-2} is a mixture of a truncated lognormal with a uniform random variable, and this distribution depends on X_{t-2}^* . Since (X_{t+1}^*, X_{t-2}^*) are correlated, conditional on w_{t-1} (which does not include X_{t-1}^*), the one-to-one assumption should be satisfied.

The second requirement in Assumption 2 requires that, for all w_t , the mapping $L_{M_{t+1}|w_t, X_t^*}$ is one-to-one. As before, consider a value w_t such that $Y_t = 1$. In Specification A, $M_{t+1}|w_t, X_t^*$ is uniformly distributed on $[0, 1]$, regardless of the value of X_t^* . Hence, the one-to-one requirement fails. For Specification B, however, M_{t+1} is distributed according to a mixture distribution which depends on X_{t+1}^* . Given the serial correlation between X_{t+1}^* and X_t^* , the one-to-one assumption should be satisfied.

Assumption 3 Assumption 3 concerns the behavior of $f_{W_t|W_{t-1}, X_{t-1}^*}$, at fixed values of w_t, w_{t-1} but holding for all values of X_t^* . We focus here on the sufficient condition (7), given right before Assumption 3, that for given (w_t, w_{t-1}) , the density $f_{W_t|W_{t-1}, X_t^*}$ must be bounded strictly between 0 and $+\infty$. We note that

$$f_{W_t|W_{t-1}, X_t^*} = f_{Y_t|M_t, X_t^*} \cdot f_{M_t|X_t^*, Y_{t-1}, M_{t-1}}.$$

The mileage transition $f_{M_t|X_t^*, Y_{t-1}, M_{t-1}}$ is a uniform distribution, so it is bounded away from zero and $+\infty$. Moreover, the CCP $f_{Y_t|M_t, X_t^*}$ is a logit probability. Because the per-period utilities (under both specification A and B), net of the ϵ 's, are bounded away from $-\infty$ and $+\infty$, the logit choice probabilities are also bounded away from zero.

The bounded support assumption on the observed state variable M_t is crucial here. However, in practice, these assumptions on M_t imply very little loss in generality, because typically in estimating these models, one can take the upper and lower bounds on M_t from the observed data.

Assumption 4 Assumption 4 requires that there exist a known functional G such that $G \left[f_{M_{t+1}|M_t, Y_t, X_t^*}(\cdot | m_t, y_t, x_t^*) \right]$ is monotonic in x_t^* . Let the functional G map a density f to its median, i.e., $G[f] = \inf \left\{ \tilde{x}^* : \int_{-\infty}^{\tilde{x}^*} f(x) dx \geq 0.5 \right\}$. Equations 12 and 16 imply that

$$M_{t+1} = \begin{cases} M_t + \eta_{t+1} \cdot \exp(0.2\nu_{t+1}) \cdot \exp(0.3\psi(M_t)) \cdot \exp(0.5X_t^*) & \text{if } Y_t = 0 \\ \eta_{t+1} \cdot \exp(0.2\nu_{t+1}) \cdot \exp(0.8X_t^*) & \text{if } Y_t = 1. \end{cases} \quad (17)$$

Let constant C_{med} stand for the median of the random variable $\eta_{t+1} \cdot \exp(0.2\nu_{t+1})$, which is a product of a uniform and a truncated lognormal random variable. Given the distribution of η_{t+1} and ν_{t+1} and the value of (y_t, m_t) , we have

$$G \left[f_{M_{t+1}|M_t, Y_t, X_t^*}(\cdot | m_t, y_t, x_t^*) \right] = \begin{cases} m_t + C_{med} \cdot \exp(0.3\psi(m_t)) \cdot \exp(0.5x_t^*) & \text{if } y_t = 0 \\ C_{med} \cdot \exp(0.8x_t^*) & \text{if } y_t = 1, \end{cases}$$

which is monotonic in x_t^* . The normalization just requires redefining x_t^* according to the equation above in the whole identification procedure.

Assumption 5 This assumption requires that, for a given pair w_t, w_{t-1} , the Markov transition kernel $L_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ is one-to-one. For both specifications of the Rust model, we can factorize

$$f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*} = f_{Y_t | M_t, X_t^*} \cdot f_{M_t | X_t^*, Y_{t-1}, M_{t-1}} \cdot f_{X_t^* | X_{t-1}^*, Y_{t-1}, M_{t-1}}.$$

Because $f_{X_t^* | X_{t-1}^*, Y_{t-1}, M_{t-1}}$ is a truncated normal density which is differentiable and positive everywhere on its support, the one-to-one requirement is satisfied unless there are w_{t-1}, w_t such that the CCP $f_{Y_t | M_t, X_t^*}$ and the mileage transition $f_{M_t | X_t^*, Y_{t-1}, M_{t-1}}$ are equal to zero for multiple values of X_t^* . However, our support assumptions (see discussion under Assumption 3 before) already imply that both of these quantities are bounded away from zero.

4 Using the Markov Kernel $W_t, X_t^* | W_{t-1}, X_{t-1}^*$ to Identify DDC models

The identification of the Markov kernel $W_t, X_t^* | W_{t-1}, X_{t-1}^*$ is only the first step in establishing nonparametric identification of the underlying dynamic model. However, once $W_t, X_t^* | W_{t-1}, X_{t-1}^*$ can be identified, nonparametric identification of the remaining parts of the models – particularly, the per-period utility functions – can proceed by straightforward application of the identification results in Magnac and Thesmar (2002) and Bajari, Chernozhukov, Hong, and Nekipelov (2007), which were developed for dynamic models without persistent latent variables X_t^* . In this section, we use the identification arguments in Bajari, Chernozhukov, Hong, and Nekipelov (2007) to show nonparametric identification of the per-period utility functions once the nonparametric identification of $W_t, X_t^* | W_{t-1}, X_{t-1}^*$

has been established.

We make the following assumptions, which are standard in this literature (except for the inclusion of X_t^* as the unobserved state variable):

1. Agents are optimizing in an infinite-horizon, stationary setting. Hence, $W_t, X_t^* | W_{t-1} X_{t-1}^*$ is identical for all periods t . Therefore, in the rest of this section, we use primes 's to denote next-period values.
2. Actions Y are chosen from the set $\mathcal{Y} = \{0, 1, \dots, K\}$.
3. The state variables are $S \equiv (M, X^*)$.
4. The per-period utility from taking action $y \in \mathcal{Y}$ in period t is:

$$u_y(S_t) + \epsilon_{y,t}, \quad \forall y \in \mathcal{Y}.$$

The $\epsilon_{y,t}$'s are utility shocks which are independent of S_t , and distributed i.i.d with known distribution $F(\epsilon)$ across periods t and actions y . Let $\vec{\epsilon}_t \equiv (\epsilon_{0,t}, \epsilon_{1,t}, \dots, \epsilon_{K,t})$.

5. From the data, the CCP's

$$p_y(S) \equiv \text{Prob}(Y = 1|S),$$

and the Markov transition kernel for S , denoted $p(S'|Y, S)$, are identified. Nonparametric identification of these two elements was the main result demonstrated in Section 2 of this paper.

6. $u_0(S)$, the per-period utility from $Y = 0$, is normalized to zero, for all S .
7. β , the discount factor, is known.⁶

Following the arguments in Magnac and Thesmar (2002) and Bajari, Chernozhukov, Hong, and Nekipelov (2007), we will show the nonparametric identification of $u_y(\cdot)$, $y = 1, \dots, K$, the per-period utility functions for all action except $Y = 0$.

The Bellman equation for this dynamic optimization problem is

$$V(S, \vec{\epsilon}) = \max_{y \in \mathcal{Y}} \left(u_y(S) + \epsilon_y + \beta E_{S', \vec{\epsilon}' | Y, S} V(S', \vec{\epsilon}') \right)$$

⁶Magnac and Thesmar (2002) discuss the possibility of identifying β via exclusion restrictions, but we do not pursue that here.

where $V(S, \bar{\epsilon})$ denotes the value function. We define the choice-specific value function as

$$V_y(S) \equiv u_y(S) + \beta E_{S', \bar{\epsilon}' | Y, S} V(S', \bar{\epsilon}').$$

Given these definitions, an agent's optimal choice when the state is S is given by

$$y^*(S) = \operatorname{argmax}_{y \in \mathcal{Y}} (V_y(S) + \epsilon_y).$$

Hotz and Miller (1993) and Magnac and Thesmar (2002) show that in this setting, there is a known one-to-one mapping, $q(S) : \mathbb{R}^K \rightarrow \mathbb{R}^K$, which maps the K -vector of choice probabilities $(p_1(S), \dots, p_K(S))$ to the K -vector $(\Delta_1(S), \dots, \Delta_K(S))$, where $\Delta_y(S)$ denotes the difference in choice-specific value functions

$$\Delta_y(S) \equiv V_y(S) - V_0(S).$$

Let the i -th element of $q(p_1(S), \dots, p_K(S))$, denoted $q_i(S)$, be equal to $\Delta_i(S)$. The known mapping q derives just from $F(\epsilon)$, the known distribution of the utility shocks.

Hence, since the choice probabilities can be identified from the data, and the mapping q is known, the value function differences $\Delta_1(S), \dots, \Delta_K(S)$ is also known.

Next, we note that the choice-specific value function also satisfies a Bellman-like equation:

$$\begin{aligned} V_y(S) &= u_y(S) + \beta E_{S' | S, Y} \left[E_{\bar{\epsilon}'} \max_{y' \in \mathcal{Y}} (V_{y'}(S') + \epsilon'_{y'}) \right] \\ &= u_y(S) + \beta E_{S' | S, Y} [G(\Delta_1(S'), \dots, \Delta_K(S')) + V_0(S')] \end{aligned} \tag{18}$$

where $G(\dots)$ denotes McFadden's "social surplus" function, for random utility models (cf. Rust (1994, pp. 3104ff)). Like the q mapping, G is a known function, which depends just on $F(\epsilon)$, the known distribution of the utility shocks.

From the normalization assumption that $u_0(S) = 0, \forall S$, we can write the Bellman equation for $V_0(S)$ as

$$V_0(S) = \beta E_{S' | S, Y} [G(\Delta_1(S'), \dots, \Delta_K(S')) + V_0(S')]. \tag{19}$$

In this equation, everything is known (including, importantly, the distribution of $S' | S, Y$), except the $V_0(\cdot)$ function. Hence, by iterating over Eq. (19), we can recover the $V_0(S)$

function. Once $V_0(\cdot)$ is known, the other choice-specific value functions can be recovered as

$$V_y(S) = \Delta_y(S) + V_0(S), \quad \forall y \in \mathcal{Y}, \quad \forall S.$$

Finally, the per-period utility functions $u_y(S)$ can be recovered from the choice-specific value functions as

$$u_y(S) = V_y(S) - \beta E_{S'|S,Y} [G(\Delta_1(S'), \dots, \Delta_K(S')) + V_0(S')], \quad \forall y \in \mathcal{Y}, \quad \forall S,$$

where everything on the right-hand side is known.

REMARK: For the case where $F(\epsilon)$ is the Type 1 Extreme Value distribution, the social surplus function is

$$G(\Delta_1(S), \dots, \Delta_K(S)) = \log \left[1 + \sum_{y=1}^K \exp(\Delta_y(S)) \right]$$

and the mapping q is such that

$$q_y(S) = \Delta_y(S) = \log(p_y(S)) - \log(p_0(S)), \quad \forall y = 1, \dots, K,$$

where $p_0(S) \equiv 1 - \sum_{y=1}^K p_y(S)$. ■

5 Concluding remarks

In this paper, we have considered the identification of a Markov process $\{W_t, X_t^*\}$ for $t = 1, 2, \dots, T$ when only $\{W_t\}$ for $t = 1, 2, \dots, T$ is observed. We showed that the joint distribution $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is identified from the observed distribution of the five observations $W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}$ under reasonable assumptions. Identification of $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ is a crucial input in methodologies for estimating dynamic models based on the “conditional-choice-probability (CCP)” approach pioneered by Hotz and Miller.

In the identification arguments, we have not invoked a stationarity assumption, which would require that the $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ be invariant across periods t . Because of this, our identification argument works in both stationary and non-stationary settings. One caveat is that, because we require the five observations $W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}$ to identify $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ for every t , we would only be able to identify $f_{W_t, X_t^*, W_{t-1}, X_{t-1}^*}$ from period

$t = 4, \dots, T - 1$.

Another assumption we made is that the unobserved state variable X_t^* is scalar-valued. We believe the proof can be extended to cases where X_t^* is a multivariate process. This may enable our identification procedure to be applied to dynamic game settings, where M_t and X_t^* may contain the set of, respectively, observed and unobserved state variables for all agents in the game.

Finally, this paper has focused completely on identification, but not estimation. While our identification proof is constructive, and can be mimicked directly for estimation, it is cumbersome to invert the functional operators computationally. For this reason, it may be more convenient to estimate using a semi-nonparametric sieve maximum likelihood procedure (Carroll, Chen, and Hu (2008)). In ongoing work, we are applying our identification results to estimate dynamic discrete-choice models with unobserved state variables.

A Proofs

Proof. (Theorem 1)

We prove the identification result in six steps. First, we show an equation which links the observed and the unobserved densities; Second, such an equation implies a relationship between corresponding linear operators; Third, we reduce the number of unknown by showing that an observed operator has an inherent eigenvalue-eigenfunction decomposition; Fourth, the uniqueness of the decomposition implies the identification of $f_{V_{t+1}|W_t, X_t^*}$; Fifth, we show that $f_{W_{t+1}|W_t, X_t^*}$ and $f_{W_t, X_t^*, W_{t-1}, Z_{t-2}}$ are also identified; Finally, the identification of the density $f_{W_t, X_t^*|W_{t-1}, X_{t-1}^*}$ follows.

First, we show that assumption 1 implies

$$f_{W_{t+1}, W_t|W_{t-1}, W_{t-2}} = \int f_{W_{t+1}|W_t, X_t^*} f_{W_t|W_{t-1}, X_t^*} f_{X_t^*|W_{t-1}, W_{t-2}} dx_t^*. \quad (20)$$

For simplicity, we omit all the arguments in the density functions. Assumption 1(i) implies that

$$\begin{aligned} & f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}} \\ = & \int \int f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}, X_t^*, X_{t-1}^*} dx_t^* dx_{t-1}^* \\ = & \int \int f_{W_{t+1}|W_t, W_{t-1}, W_{t-2}, X_t^*, X_{t-1}^*} f_{W_t, X_t^*|W_{t-1}, W_{t-2}, X_{t-1}^*} f_{W_{t-1}, W_{t-2}, X_{t-1}^*} dx_t^* dx_{t-1}^* \\ = & \int \int f_{W_{t+1}|W_t, X_t^*} f_{W_t, X_t^*|W_{t-1}, X_{t-1}^*} f_{W_{t-1}, W_{t-2}, X_{t-1}^*} dx_t^* dx_{t-1}^* \\ = & \int \int f_{W_{t+1}|W_t, X_t^*} f_{W_t|W_{t-1}, X_t^*, X_{t-1}^*} f_{X_t^*|W_{t-1}, X_{t-1}^*} f_{W_{t-1}, W_{t-2}, X_{t-1}^*} dx_t^* dx_{t-1}^* \\ = & \int \int f_{W_{t+1}|W_t, X_t^*} f_{W_t|W_{t-1}, X_t^*, X_{t-1}^*} f_{X_t^*|W_{t-1}, W_{t-2}, X_{t-1}^*} f_{W_{t-1}, W_{t-2}, X_{t-1}^*} dx_t^* dx_{t-1}^* \\ = & \int \int f_{W_{t+1}|W_t, X_t^*} f_{W_t|W_{t-1}, X_t^*, X_{t-1}^*} f_{X_t^*, X_{t-1}^*, W_{t-1}, W_{t-2}} dx_t^* dx_{t-1}^*. \end{aligned}$$

Assumption 1(ii) then implies that

$$\begin{aligned}
& f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}} \\
&= \int f_{W_{t+1}|W_t, X_t^*} f_{W_t|W_{t-1}, X_t^*} \left(\int f_{X_t^*, X_{t-1}^*, W_{t-1}, W_{t-2}} dx_{t-1}^* \right) dx_t^* \\
&= \int f_{W_{t+1}|W_t, X_t^*} f_{W_t|W_{t-1}, X_t^*} f_{X_t^*, W_{t-1}, W_{t-2}} dx_t^*.
\end{aligned}$$

Second, we show that equation 20 implies an equality between corresponding operators. Let $\mathcal{L}^p(\mathcal{X})$, $1 \leq p < \infty$ stand for the space of function $h(\cdot)$ with $\int_{\mathcal{X}} |h(x)|^p dx < \infty$, and let $\mathcal{L}^\infty(\mathcal{X})$ denote the space of function $h(\cdot)$ with $\sup_{x \in \mathcal{X}} |h(x)| < \infty$. For any $1 \leq p \leq \infty$, we define operators as follows: for any function $h \in \mathcal{L}^p(\mathcal{W}_{t-m})$

$$\begin{aligned}
L_{W_{t+1}, w_t | w_{t-1}, W_{t-2}} & : \mathcal{L}^p(\mathcal{W}_{t-2}) \rightarrow \mathcal{L}^p(\mathcal{W}_{t+1}), \\
(L_{W_{t+1}, w_t | w_{t-1}, W_{t-2}} h)(x) &= \int f_{W_{t+1}, W_t | W_{t-1}, W_{t-2}}(x, w_t | w_{t-1}, z) h(z) dz,
\end{aligned}$$

$$\begin{aligned}
L_{W_{t+1} | w_t, X_t^*} & : \mathcal{L}^p(\mathcal{X}_t^*) \rightarrow \mathcal{L}^p(\mathcal{W}_{t+1}), \\
(L_{W_{t+1} | w_t, X_t^*} h)(x) &= \int f_{W_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) h(x_t^*) dx_t^*,
\end{aligned}$$

$$\begin{aligned}
L_{X_t^* | w_{t-1}, W_{t-2}} & : \mathcal{L}^p(\mathcal{W}_{t-2}) \rightarrow \mathcal{L}^p(\mathcal{X}_t^*), \\
(L_{X_t^* | w_{t-1}, W_{t-2}} h)(x_t^*) &= \int f_{X_t^* | W_{t-1}, W_{t-2}}(x_t^* | w_{t-1}, z) h(z) dz,
\end{aligned}$$

$$\begin{aligned}
D_{w_t | w_{t-1}, X_t^*} & : \mathcal{L}^p(\mathcal{X}_t^*) \rightarrow \mathcal{L}^p(\mathcal{X}_t^*), \\
(D_{w_t | w_{t-1}, X_t^*} h)(x_t^*) &= f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) h(x_t^*).
\end{aligned}$$

Notice that the operator $D_{w_t | w_{t-1}, X_t^*}$ is a "diagonal" or multiplication operator. As shown in Hu and Schennach (2008), the identification of an operator, e.g, $L_{V_{t+1} | w_t, X_t^*}$, is equivalent to that of its corresponding density, e.g., $f_{V_{t+1} | W_t, X_t^*}$. For any given $(w_t, w_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1}$,

we have for any function $h \in \mathcal{L}^p(\mathcal{W}_{t-2})$

$$\begin{aligned}
& \left(L_{W_{t+1}, w_t | w_{t-1}, W_{t-2}} h \right) (x) \\
&= \int f_{W_{t+1}, W_t | W_{t-1}, W_{t-2}}(x, w_t | w_{t-1}, z) h(z) dz \\
&= \int f_{W_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) \left(\int f_{X_t^* | W_{t-1}, W_{t-2}}(x_t^* | w_{t-1}, z) h(z) dz \right) dx_t^* \\
&= \int f_{W_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) \left(L_{X_t^* | w_{t-1}, W_{t-2}} h \right) (x_t^*) dx_t^* \\
&= \int f_{W_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) \left(D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, W_{t-2}} h \right) (x_t^*) dx_t^* \\
&= \left(L_{W_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, W_{t-2}} h \right) (x).
\end{aligned}$$

Therefore, equation 20 is equivalent to

$$L_{W_{t+1}, w_t | w_{t-1}, W_{t-2}} = L_{W_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, W_{t-2}}. \quad (21)$$

Let $g, q : \mathbb{R}^d \rightarrow \mathbb{R}$, and

$$\begin{aligned}
V_{t+1} &= g(W_{t+1}), \\
Z_{t-2} &= q(W_{t-2}).
\end{aligned}$$

We may apply the same procedure to the joint density of $\{V_{t+1}, W_t, W_{t-1}, Z_{t-2}\}$ for any $(x, w_t, w_{t-1}, z) \in g(\mathcal{W}_{t+1}) \times \mathcal{W}_t \times \mathcal{W}_{t-1} \times q(\mathcal{W}_{t-2})$ to obtain

$$L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} = L_{V_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}}, \quad (22)$$

where

$$\begin{aligned}
L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} &: \mathcal{L}^p(q(\mathcal{W}_{t-2})) \rightarrow \mathcal{L}^p(g(\mathcal{W}_{t+1})), \\
(L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} h)(x) &= \int f_{V_{t+1}, W_t | W_{t-1}, Z_{t-2}}(x, w_t | w_{t-1}, z) h(z) dz,
\end{aligned}$$

$$\begin{aligned}
L_{V_{t+1} | w_t, X_t^*} &: \mathcal{L}^p(\mathcal{X}_t^*) \rightarrow \mathcal{L}^p(g(\mathcal{W}_{t+1})), \\
(L_{V_{t+1} | w_t, X_t^*} h)(x) &= \int f_{V_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) h(x_t^*) dx_t^*,
\end{aligned}$$

$$\begin{aligned}
L_{X_t^*|w_{t-1}, Z_{t-2}} & : \mathcal{L}^p(q(\mathcal{W}_{t-2})) \rightarrow \mathcal{L}^p(\mathcal{X}_t^*), \\
\left(L_{X_t^*|w_{t-1}, Z_{t-2}} h\right)(x_t^*) & = \int f_{X_t^*|W_{t-1}, Z_{t-2}}(x_t^*|w_{t-1}, z) h(z) dz.
\end{aligned}$$

Notice that the operator $L_{V_{t+1}|w_t, X_t^*}$ does not depend on w_{t-1} and $L_{X_t^*|w_{t-1}, Z_{t-2}}$ does not depend on w_t . This important fact may help the identification of $L_{V_{t+1}|w_t, X_t^*}$ in equation 22.

Third, we show that an observed operator may have an inherent eigenvalue-eigenfunction decomposition, where the eigenfunctions are $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$. For any $w_t \in \mathcal{W}_t$, we consider with $\bar{w}_t \in \mathcal{W}_t$, $w_{t-1}, \bar{w}_{t-1} \in \mathcal{A}(w_t) \cap \mathcal{A}(\bar{w}_t)$, $\bar{w}_{t-1} \neq w_{t-1}$, and $\bar{w}_t \neq w_t$,

$$\text{for } (w_t, w_{t-1}) : L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}} = L_{V_{t+1}|w_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}}, \quad (23)$$

$$\text{for } (\bar{w}_t, w_{t-1}) : L_{V_{t+1}, \bar{w}_t|w_{t-1}, Z_{t-2}} = L_{V_{t+1}|\bar{w}_t, X_t^*} D_{\bar{w}_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}}, \quad (24)$$

$$\text{for } (w_t, \bar{w}_{t-1}) : L_{V_{t+1}, w_t|\bar{w}_{t-1}, Z_{t-2}} = L_{V_{t+1}|w_t, X_t^*} D_{w_t|\bar{w}_{t-1}, X_t^*} L_{X_t^*|\bar{w}_{t-1}, Z_{t-2}}, \quad (25)$$

$$\text{for } (\bar{w}_t, \bar{w}_{t-1}) : L_{V_{t+1}, \bar{w}_t|\bar{w}_{t-1}, Z_{t-2}} = L_{V_{t+1}|\bar{w}_t, X_t^*} D_{\bar{w}_t|\bar{w}_{t-1}, X_t^*} L_{X_t^*|\bar{w}_{t-1}, Z_{t-2}}. \quad (26)$$

Assumptions 2 and 3 guarantee that the inverse of the operators on the left hand side exist. Eliminating $L_{X_t^*|w_{t-1}, Z_{t-2}}$ in equations 23 and 24 leads to

$$\begin{aligned}
\mathbf{A} & \equiv L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}} L_{V_{t+1}, \bar{w}_t|w_{t-1}, Z_{t-2}}^{-1} \\
& = L_{V_{t+1}|w_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}} \left(L_{V_{t+1}|\bar{w}_t, X_t^*} D_{\bar{w}_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}} \right)^{-1} \\
& = L_{V_{t+1}|w_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} D_{\bar{w}_t|w_{t-1}, X_t^*}^{-1} L_{V_{t+1}|\bar{w}_t, X_t^*}^{-1}. \quad (27)
\end{aligned}$$

Similarly, eliminating $L_{X_t^*|\bar{w}_{t-1}, Z_{t-2}}$ in equations 25 and 26 results in

$$\begin{aligned}
\mathbf{B} & \equiv L_{V_{t+1}, w_t|\bar{w}_{t-1}, Z_{t-2}} L_{V_{t+1}, \bar{w}_t|\bar{w}_{t-1}, Z_{t-2}}^{-1} \\
& = L_{V_{t+1}|w_t, X_t^*} D_{w_t|\bar{w}_{t-1}, X_t^*} D_{\bar{w}_t|\bar{w}_{t-1}, X_t^*}^{-1} L_{V_{t+1}|\bar{w}_t, X_t^*}^{-1}. \quad (28)
\end{aligned}$$

We then eliminate $L_{V_{t+1}|\bar{w}_t, X_t^*}^{-1}$ in equations 27 and 28 to obtain

$$\begin{aligned}
\mathbf{AB}^{-1} &\equiv L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}} L_{V_{t+1}, \bar{w}_t | w_{t-1}, Z_{t-2}}^{-1} \left(L_{V_{t+1}, w_t | \bar{w}_{t-1}, Z_{t-2}} L_{V_{t+1}, \bar{w}_t | \bar{w}_{t-1}, Z_{t-2}}^{-1} \right)^{-1} \\
&= L_{V_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} D_{\bar{w}_t | w_{t-1}, X_t^*}^{-1} L_{V_{t+1} | \bar{w}_t, X_t^*}^{-1} \times \\
&\quad \times \left(L_{V_{t+1} | w_t, X_t^*} D_{w_t | \bar{w}_{t-1}, X_t^*} D_{\bar{w}_t | \bar{w}_{t-1}, X_t^*}^{-1} L_{V_{t+1} | \bar{w}_t, X_t^*}^{-1} \right)^{-1} \\
&= L_{V_{t+1} | w_t, X_t^*} \left(D_{w_t | w_{t-1}, X_t^*} D_{\bar{w}_t | w_{t-1}, X_t^*}^{-1} D_{\bar{w}_t | \bar{w}_{t-1}, X_t^*} D_{w_t | \bar{w}_{t-1}, X_t^*}^{-1} \right) L_{V_{t+1} | w_t, X_t^*}^{-1} \\
&\equiv L_{V_{t+1} | w_t, X_t^*} D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} L_{V_{t+1} | w_t, X_t^*}^{-1}, \tag{29}
\end{aligned}$$

where

$$\begin{aligned}
& \left(D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} h \right) (x_t^*) \\
&= \left(D_{w_t | w_{t-1}, X_t^*} D_{\bar{w}_t | w_{t-1}, X_t^*}^{-1} D_{\bar{w}_t | \bar{w}_{t-1}, X_t^*} D_{w_t | \bar{w}_{t-1}, X_t^*}^{-1} h \right) (x_t^*) \\
&= k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) h(x_t^*),
\end{aligned}$$

$$k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) = \frac{f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) f_{W_t | W_{t-1}, X_t^*}(\bar{w}_t | \bar{w}_{t-1}, x_t^*)}{f_{W_t | W_{t-1}, X_t^*}(\bar{w}_t | w_{t-1}, x_t^*) f_{W_t | W_{t-1}, X_t^*}(w_t | \bar{w}_{t-1}, x_t^*)}.$$

This equation implies that the observed operator \mathbf{AB}^{-1} on the left hand side of equation 29 has an inherent eigenvalue-eigenfunction decomposition. The eigenfunctions are $f_{V_{t+1} | W_t, X_t^*}(\cdot | w_t, x_t^*)$, which is normalized by $\int f_{V_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) dx = 1$. Notice that the eigenfunction in $L_{V_{t+1} | W_t, X_t^*}$ does not depend on \bar{w}_t , w_{t-1} , or \bar{w}_{t-1} , while the eigenvalue in $D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*}$ may be different for a different \bar{w}_t , w_{t-1} , or \bar{w}_{t-1} .

Fourth, we show that the uniqueness of the decomposition in equation 29. Notice that the decomposition in equation 29 is similar to but more complicated than the decomposition in Hu and Schennach (2008) or Carroll, Chen, and Hu (2008). Their results imply that such a decomposition is unique under assumptions 3 and 4. We may show the reasoning as follows. Suppose that for two indices $\hat{x}_t^* \neq \tilde{x}_t^*$ the two eigenvalues are the same, i.e., $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, \hat{x}_t^*) = k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, \tilde{x}_t^*)$ for some $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1} \times \mathcal{W}_{t-1}$. Therefore, we can't identify the two corresponding eigenfunctions. But assumption 3 guarantees that there exist another $(\hat{w}_t, \tilde{w}_{t-1}, \hat{w}_{t-1}) \in \mathcal{W}_t \times \mathcal{W}_{t-1} \times \mathcal{W}_{t-1}$ such that $k(w_t, \hat{w}_t, \tilde{w}_{t-1}, \hat{w}_{t-1}, \hat{x}_t^*) \neq k(w_t, \hat{w}_t, \tilde{w}_{t-1}, \hat{w}_{t-1}, \tilde{x}_t^*)$, which are two eigenvalues corresponding to the same eigenfunctions $f_{V_{t+1} | W_t, X_t^*}(\cdot | w_t, \hat{x}_t^*)$ and $f_{V_{t+1} | W_t, X_t^*}(\cdot | w_t, \tilde{x}_t^*)$. Therefore, the eigenfunction $f_{V_{t+1} | W_t, X_t^*}(\cdot | w_t, x_t^*)$ is identified up to the value of x_t^* for any

given $w_t \in \mathcal{W}_t$. Moreover, assumption 4 reveals the value of x_t^* in each eigenfunction $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$. Hence, the density $f_{V_{t+1}|W_t, X_t^*}$ or $L_{V_{t+1}|w_t, X_t^*}$ is nonparametrically identified for any given $w_t \in \mathcal{W}_t$.

Fifth, we show the identification of the density $f_{W_{t+1}|W_t, X_t^*}$. Equations 20 and 22 imply for any given $w_t \in \mathcal{W}_t$

$$\begin{aligned} & L_{W_{t+1}, w_t|w_{t-1}, Z_{t-2}} L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}}^{-1} L_{V_{t+1}|w_t, X_t^*} \\ &= L_{W_{t+1}|w_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}} \left(L_{V_{t+1}|w_t, X_t^*} D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}} \right)^{-1} L_{V_{t+1}|w_t, X_t^*} \\ &= L_{W_{t+1}|w_t, X_t^*}, \end{aligned}$$

where the left hand side is identified. Moreover, the following equation

$$f_{V_{t+1}, W_t, W_{t-1}, W_{t-2}} = \int f_{V_{t+1}|W_t, X_t^*} f_{W_t, X_t^*, W_{t-1}, W_{t-2}} dx_t^*$$

implies that for any given $w_t \in \mathcal{W}_t$,

$$f_{V_{t+1}, W_t=w_t, W_{t-1}, W_{t-2}} = L_{V_{t+1}|w_t, X_t^*} f_{W_t=w_t, X_t^*, W_{t-1}, W_{t-2}}.$$

Therefore, we identify $f_{W_t=w_t, X_t^*, W_{t-1}, W_{t-2}}$ for any given $w_t \in \mathcal{W}_t$ through

$$f_{W_t=w_t, X_t^*, W_{t-1}, W_{t-2}} = L_{V_{t+1}|w_t, X_t^*}^{-1} f_{V_{t+1}, W_t=w_t, W_{t-1}, W_{t-2}}.$$

In summary, the densities $f_{W_{t+1}|W_t, X_t^*}$ and $f_{W_t, X_t^*, W_{t-1}, W_{t-2}}$ are identified from $f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}}$. Given the known function q in $Z_{t-2} = q(W_{t-2})$, the identification of $f_{W_t, X_t^*, W_{t-1}, W_{t-2}}$ implies that of $f_{W_t, X_t^*, W_{t-1}, Z_{t-2}}$.

Furthermore, we show that the operator corresponding to $f_{W_t, X_t^*|W_{t-1}, Z_{t-2}}$ is in fact $D_{w_t|w_{t-1}, X_t^*} L_{X_t^*|w_{t-1}, Z_{t-2}}$.

Similar to the proof of equation 20, we have

$$\begin{aligned}
f_{W_t, X_t^* | W_{t-1}, Z_{t-2}} &= \int f_{W_t, X_t^*, X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^* \\
&= \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*, Z_{t-2}} f_{X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^* \\
&= \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*} f_{X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^* \\
&= \int f_{W_t | W_{t-1}, X_t^*, X_{t-1}^*} f_{X_t^* | W_{t-1}, X_{t-1}^*} f_{X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^* \\
&= \int f_{W_t | W_{t-1}, X_t^*} f_{X_t^*, X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^* \\
&= f_{W_t | W_{t-1}, X_t^*} f_{X_t^* | W_{t-1}, Z_{t-2}}.
\end{aligned} \tag{30}$$

The corresponding operator of $f_{W_t | W_{t-1}, X_t^*} f_{X_t^* | W_{t-1}, Z_{t-2}}$ is $D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}}$, which is identified through equation 22 as follows:

$$D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}} = L_{V_{t+1} | w_t, X_t^*}^{-1} L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}}. \tag{31}$$

Finally, we show the identification of the density $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$. Assumption 1 implies

$$\begin{aligned}
f_{W_t, X_t^* | W_{t-1}, Z_{t-2}} &= \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*, Z_{t-2}} f_{X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^* \\
&= \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*} f_{X_{t-1}^* | W_{t-1}, Z_{t-2}} dx_{t-1}^*.
\end{aligned} \tag{32}$$

The left hand side has been identified in the previous step. Thus far, we have only used the four observations $W_{t+1}, W_t, W_{t-1}, W_{t-2}$. In order to identify the density $f_{X_{t-1}^* | W_{t-1}, Z_{t-2}}$ on the right hand side of equation 32, we use one more period of the data W_{t-3} or Z_{t-3} . Replacing t by $t-1$ in the previous procedure implies that the observed density $f_{V_t, W_{t-1} | W_{t-2}, Z_{t-3}}$ uniquely determines $f_{W_{t-1}, X_{t-1}^* | W_{t-2}, Z_{t-3}}$. Therefore, the density $f_{X_{t-1}^* | W_{t-1}, Z_{t-2}}$ in equation 32 is identified from

$$f_{X_{t-1}^*, W_{t-1}, W_{t-2}} = \int f_{W_{t-1}, X_{t-1}^* | W_{t-2}, Z_{t-3}} f_{W_{t-2}, Z_{t-3}} dz_{t-3}$$

with $Z_{t-2} = q(W_{t-2})$ for the known function q .

Given that $f_{W_t, X_t^* | W_{t-1}, Z_{t-2}}$ and $f_{X_{t-1}^* | W_{t-1}, Z_{t-2}}$ are identified, equation 32 implies that the density of interest $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ may be identified as follows. For any given $w_{t-1} \in \mathcal{W}_{t-1}$

and $w_t \in \mathcal{W}_t$, we define

$$\begin{aligned} L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} & : \mathcal{L}^p(\mathcal{X}_{t-1}^*) \rightarrow \mathcal{L}^p(\mathcal{X}_t^*), \\ \left(L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} h \right) (x_t^*) & = \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}(w_t, x_t^* | w_{t-1}, x_{t-1}^*) h(x_{t-1}^*) dx_{t-1}^* \end{aligned}$$

and

$$\begin{aligned} L_{X_{t-1}^* | w_{t-1}, Z_{t-2}} & : \mathcal{L}^p(q(\mathcal{W}_{t-2})) \rightarrow \mathcal{L}^p(\mathcal{X}_{t-1}^*), \\ \left(L_{X_{t-1}^* | w_{t-1}, Z_{t-2}} h \right) (x_{t-1}^*) & = \int f_{X_{t-1}^* | W_{t-1}, Z_{t-2}}(x_{t-1}^* | w_{t-1}, z) h(z) dz. \end{aligned}$$

As shown above, the corresponding operator of $f_{W_t | W_{t-1}, X_t^*} f_{X_t^* | W_{t-1}, Z_{t-2}}$ is $D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}}$, which has been identified in equation 31. We then show that equation 32 implies

$$D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}} = L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} L_{X_{t-1}^* | w_{t-1}, Z_{t-2}} \quad (33)$$

as follows. For any function $h \in \mathcal{L}^p(q(\mathcal{W}_{t-2}))$, equation 30 implies

$$\begin{aligned} & \left(D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}} h \right) (x_t^*) \\ & = \int f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) f_{X_t^* | W_{t-1}, Z_{t-2}}(x_t^* | w_{t-1}, z) h(z) dz \\ & = \int f_{W_t, X_t^* | W_{t-1}, Z_{t-2}}(w_t, x_t^* | w_{t-1}, z) h(z) dz. \end{aligned}$$

Equation 32 then implies

$$\begin{aligned} & \left(D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, Z_{t-2}} h \right) (x_t^*) \\ & = \int \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}(w_t, x_t^* | w_{t-1}, x_{t-1}^*) f_{X_{t-1}^* | W_{t-1}, Z_{t-2}}(x_{t-1}^* | w_{t-1}, z) dx_{t-1}^* h(z) dz \\ & = \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}(w_t, x_t^* | w_{t-1}, x_{t-1}^*) \left(\int f_{X_{t-1}^* | W_{t-1}, Z_{t-2}}(x_{t-1}^* | w_{t-1}, z) h(z) dz \right) dx_{t-1}^* \\ & = \int f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}(w_t, x_t^* | w_{t-1}, x_{t-1}^*) \left(L_{X_{t-1}^* | w_{t-1}, Z_{t-2}} h \right) (x_{t-1}^*) dx_{t-1}^* \\ & = \left(L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} L_{X_{t-1}^* | w_{t-1}, Z_{t-2}} h \right) (x_t^*). \end{aligned}$$

Combining equations 31 and 33 leads to

$$L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} L_{X_{t-1}^* | w_{t-1}, Z_{t-2}} = L_{V_{t+1} | w_t, X_t^*}^{-1} L_{V_{t+1}, w_t | w_{t-1}, Z_{t-2}}.$$

Notice that assumptions 2, 3, and 5 imply that $L_{X_{t-1}^*|w_{t-1}, Z_{t-2}}$ is one-to-one. Since $L_{X_{t-1}^*|w_{t-1}, Z_{t-2}}$ and $f_{X_{t-1}^*|W_{t-1}, Z_{t-2}}$ are identified, $f_{W_t=w_t, X_t^*|W_{t-1}=w_{t-1}, X_{t-1}^*}$ and $L_{w_t, X_t^*|w_{t-1}, X_{t-1}^*}$ are identified as

$$L_{w_t, X_t^*|w_{t-1}, X_{t-1}^*} = \left(L_{V_{t+1}|w_t, X_t^*}^{-1} L_{V_{t+1}, w_t|w_{t-1}, Z_{t-2}} \right) L_{X_{t-1}^*|w_{t-1}, Z_{t-2}}^{-1}$$

for any given $w_{t-1} \in \mathcal{W}_{t-1}$ and $w_t \in \mathcal{W}_t$. Hence, the density $f_{W_t, X_t^*|W_{t-1}, X_{t-1}^*}$ is identified. ■

B Remarks

Example of $f_{W_t|W_{t-1}, X_t^*}$ satisfying Assumption 3 Because the second condition defining the set $\mathcal{B}(w_t)$ is not completely obvious, here we present an example which satisfies the condition. We seek a density $f_{W_t|W_{t-1}, X_t^*}$ such that $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) \in (0, \infty)$. Let

$$f_{W_t|W_{t-1}, X_t^*}(w_t|w_{t-1}, x_t^*) = \phi(w_t - F(x_t^*)),$$

where $\phi(\cdot)$ is the pdf of the standard normal and $F(\cdot)$ is a strictly increasing cdf. Therefore,

$$\begin{aligned} \frac{f_{W_t|W_{t-1}, X_t^*}(w_t|w_{t-1}, x_t^*)}{f_{W_t|W_{t-1}, X_t^*}(\bar{w}_t|w_{t-1}, x_t^*)} &= \frac{\phi(w_t - F(x_t^*))}{\phi(\bar{w}_t - F(x_t^*))} \\ &= \exp \left[-\frac{1}{2} \left[(w_t - F(x_t^*))^2 - (\bar{w}_t - F(x_t^*))^2 \right] \right] \end{aligned}$$

Let

$$f_{W_t|W_{t-1}, X_t^*}(w_t|\bar{w}_{t-1}, x_t^*) = \frac{1}{0.5} \phi \left(\frac{w_t - F(x_t^*)}{0.5} \right),$$

and therefore,

$$\begin{aligned} \frac{f_{W_t|W_{t-1}, X_t^*}(w_t|\bar{w}_{t-1}, x_t^*)}{f_{W_t|W_{t-1}, X_t^*}(\bar{w}_t|\bar{w}_{t-1}, x_t^*)} &= \frac{\phi \left(\frac{w_t - F(x_t^*)}{0.5} \right)}{\phi \left(\frac{\bar{w}_t - F(x_t^*)}{0.5} \right)} \\ &= \exp \left[-2 \left[(w_t - F(x_t^*))^2 - (\bar{w}_t - F(x_t^*))^2 \right] \right]. \end{aligned}$$

The kernel function then becomes

$$\begin{aligned}
k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*) &= \frac{f_{W_t|W_{t-1}, X_t^*}(w_t|w_{t-1}, x_t^*) f_{W_t|W_{t-1}, X_t^*}(\bar{w}_t|\bar{w}_{t-1}, x_t^*)}{f_{W_t|W_{t-1}, X_t^*}(\bar{w}_t|w_{t-1}, x_t^*) f_{W_t|W_{t-1}, X_t^*}(w_t|\bar{w}_{t-1}, x_t^*)} \\
&= \frac{\exp\left[-\frac{1}{2}\left[(w_t - F(x_t^*))^2 - (\bar{w}_t - F(x_t^*))^2\right]\right]}{\exp\left[-2\left[(w_t - F(x_t^*))^2 - (\bar{w}_t - F(x_t^*))^2\right]\right]} \\
&= \exp\left[\frac{3}{2}(w_t^2 - \bar{w}_t^2) - 3(w_t - \bar{w}_t)F(x_t^*)\right]
\end{aligned}$$

Thus, for any given $(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1})$, the kernel function $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*)$ is in $(0, \infty)$ for any $x_t^* \in \mathcal{X}_t^*$ because $F(x_t^*) \in [0, 1]$. Assumption 3(ii) also holds in this example because $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, x_t^*)$ is monotonic in x_t^* for any given $(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1})$. ■

C Special case: a discrete unobserved state variable

In this section, we illustrate our identification strategy in the special case where X_t^* is discrete:

$$\forall t, X_t^* \in \mathcal{X}^* \equiv \{1, 2, \dots, J\}.$$

The main difference between this discrete case and the previous continuous case is that the linear integral operators are replaced by matrices, which may be more straightforward.

Since we assume the unobserved state variable X_t^* is discrete in this section, we first discretize the observed variable W_t and then use the discretized W_t to identify the distribution involving the latent X_t^* . Let \mathcal{W}_t be the support of W_t and $\mathcal{W}_t^1, \mathcal{W}_t^2, \dots, \mathcal{W}_t^J$ be a known partition of \mathcal{W}_t . We define a discrete variable $V_t \in \mathcal{X}_t^*$ such that $V_t = j$ if $W_t \in \mathcal{W}_t^j$, i.e.,

$$V_t = \sum_{j=1}^J j \times I(W_t \in \mathcal{W}_t^j),$$

where $I(\cdot)$ is the indicator function. This mapping corresponds to the known functions g and q in the continuous case, which also implies we use X_{t-2} as Z_{t-2} in the continuous case.

Given the proof of theorem 1, a number of equations and derivations are stated without proof in this section.

Step 1: Identification of $f_{V_{t+1}|W_t, X_t^*}$. Equations 2 and 3 implies for any $x, z \in \mathcal{X}_t^*$, $w_t \in \mathcal{W}_t$, and $w_{t-1} \in \mathcal{W}_{t-1}$,

$$\begin{aligned} & f_{V_{t+1}, W_t | W_{t-1}, V_{t-2}}(x, w_t | w_{t-1}, z) \\ = & \sum_{x_t^* \in \mathcal{X}_t^*} f_{V_{t+1} | W_t, X_t^*}(x | w_t, x_t^*) f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, x_t^*) f_{X_t^* | W_{t-1}, V_{t-2}}(x_t^* | w_{t-1}, z) \end{aligned} \quad (34)$$

Define the J -by- J matrices

$$\begin{aligned} L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}} &= [f_{V_{t+1}, W_t | W_{t-1}, V_{t-2}}(i, w_t | w_{t-1}, j)]_{i,j}, \\ L_{V_{t+1} | w_t, X_t^*} &= [f_{V_{t+1} | W_t, X_t^*}(i | w_t, j)]_{i,j}, \\ L_{X_t^* | w_{t-1}, V_{t-2}} &= [f_{X_t^* | W_{t-1}, V_{t-2}}(i | w_{t-1}, j)]_{i,j}, \end{aligned}$$

for $i, j = 1, 2, \dots, J$ and a J -by- J diagonal matrix

$$D_{w_t | w_{t-1}, X_t^*} = \begin{bmatrix} f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, 1) & 0 & & 0 \\ 0 & \dots & & 0 \\ 0 & 0 & f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, J) & \end{bmatrix}$$

for $w_t \in \mathcal{W}_t$, $w_{t-1} \in \mathcal{W}_{t-1}$. Given these definitions, we can write equation 34 in matrix notation as

$$L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}} = L_{V_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, V_{t-2}}. \quad (35)$$

for any $w_t \in \mathcal{W}_t$, $w_{t-1} \in \mathcal{W}_{t-1}$. Obviously, the unknown matrices on the right hand side are not uniquely determined by the observed matrix on the left hand side without further assumptions. Notice, however, that the matrix $L_{V_{t+1} | w_t, X_t^*}$ does not depend on w_{t-1} and $L_{X_t^* | w_{t-1}, V_{t-2}}$ does not depend on w_t . This important fact in equation 35 may help the identification of $L_{V_{t+1} | w_t, X_t^*}$.

We assume that for any given $w_t \in \mathcal{W}_t$ there exists $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1})$ with $\bar{w}_t \neq w_t$ and $\bar{w}_{t-1} \neq w_{t-1} \in \mathcal{W}_{t-1}$ such that the matrices $L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}}$, $L_{V_{t+1}, \bar{w}_t | w_{t-1}, V_{t-2}}$, $L_{V_{t+1}, w_t | \bar{w}_{t-1}, V_{t-2}}$, and $L_{V_{t+1}, \bar{w}_t | \bar{w}_{t-1}, V_{t-2}}$ are all invertible. This assumption is testable from the data. Equation

35 then implies

$$\text{for } (w_t, w_{t-1}): \quad L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}} = L_{V_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, V_{t-2}}, \quad (36)$$

$$\text{for } (\bar{w}_t, w_{t-1}): \quad L_{V_{t+1}, \bar{w}_t | w_{t-1}, V_{t-2}} = L_{V_{t+1} | \bar{w}_t, X_t^*} D_{\bar{w}_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, V_{t-2}}, \quad (37)$$

$$\text{for } (w_t, \bar{w}_{t-1}): \quad L_{V_{t+1}, w_t | \bar{w}_{t-1}, V_{t-2}} = L_{V_{t+1} | w_t, X_t^*} D_{w_t | \bar{w}_{t-1}, X_t^*} L_{X_t^* | \bar{w}_{t-1}, V_{t-2}}, \quad (38)$$

$$\text{for } (\bar{w}_t, \bar{w}_{t-1}): \quad L_{V_{t+1}, \bar{w}_t | \bar{w}_{t-1}, V_{t-2}} = L_{V_{t+1} | \bar{w}_t, X_t^*} D_{\bar{w}_t | \bar{w}_{t-1}, X_t^*} L_{X_t^* | \bar{w}_{t-1}, V_{t-2}}, \quad (39)$$

where all the left hand side matrices are observed. The key identification procedure includes three eliminations. First, eliminating matrix $L_{X_t^* | w_{t-1}, V_{t-2}}$ in equations 36 and 37 leads to

$$\begin{aligned} \mathbf{A} &\equiv L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}} L_{V_{t+1}, \bar{w}_t | w_{t-1}, V_{t-2}}^{-1} \\ &= L_{V_{t+1} | w_t, X_t^*} D_{w_t | w_{t-1}, X_t^*} D_{\bar{w}_t | w_{t-1}, X_t^*}^{-1} L_{V_{t+1} | \bar{w}_t, X_t^*}^{-1}. \end{aligned} \quad (40)$$

Second, eliminating $L_{X_t^* | \bar{w}_{t-1}, V_{t-2}}$ in equations 38 and 39 results in

$$\begin{aligned} \mathbf{B} &\equiv L_{V_{t+1}, w_t | \bar{w}_{t-1}, V_{t-2}} L_{V_{t+1}, \bar{w}_t | \bar{w}_{t-1}, V_{t-2}}^{-1} \\ &= L_{V_{t+1} | w_t, X_t^*} D_{w_t | \bar{w}_{t-1}, X_t^*} D_{\bar{w}_t | \bar{w}_{t-1}, X_t^*}^{-1} L_{V_{t+1} | \bar{w}_t, X_t^*}^{-1}. \end{aligned} \quad (41)$$

Notice that matrices \mathbf{A} and \mathbf{B} are still directly estimable from the data. Third, we eliminate $L_{V_{t+1} | \bar{w}_t, X_t^*}$ in equation 40 and 41 to obtain

$$\mathbf{AB}^{-1} = L_{V_{t+1} | w_t, X_t^*} D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} L_{V_{t+1} | w_t, X_t^*}^{-1}, \quad (42)$$

where

$$D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*} = D_{w_t | w_{t-1}, X_t^*} D_{\bar{w}_t | w_{t-1}, X_t^*}^{-1} D_{\bar{w}_t | \bar{w}_{t-1}, X_t^*} D_{w_t | \bar{w}_{t-1}, X_t^*}^{-1}.$$

Since $D_{w_t | w_{t-1}, X_t^*}$ is diagonal, the matrix $D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*}$ is also diagonal with j -th diagonal entry equal to

$$k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j) = \frac{f_{W_t | W_{t-1}, X_t^*}(w_t | w_{t-1}, j) f_{W_t | W_{t-1}, X_t^*}(\bar{w}_t | \bar{w}_{t-1}, j)}{f_{W_t | W_{t-1}, X_t^*}(\bar{w}_t | w_{t-1}, j) f_{W_t | W_{t-1}, X_t^*}(w_t | \bar{w}_{t-1}, j)}.$$

Therefore, equation 42 implies that the observed matrix \mathbf{AB}^{-1} on the left hand side has an eigenvalue-eigenvector decomposition. Each value on the diagonal of $D_{w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, X_t^*}$ is an eigenvalue and each corresponding column of $L_{V_{t+1} | w_t, X_t^*}$ is a corresponding eigenvector. An eigenvector is automatically normalized because the sum of each column of $L_{V_{t+1} | w_t, X_t^*}$

is 1.

One ambiguity left is the possibility that the eigenvalues may not be distinctive. Therefore, we need to assume that for any $w_t \in \mathcal{W}_t$, there exists a $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1})$ such that $0 < k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j) < \infty$ for all $j \in \mathcal{X}_t^*$ and $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j_1) \neq k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j_2)$ for $j_1 \neq j_2$. This assumption can be relaxed to assumption 3 when equation 42 holds for another $(\hat{w}_t, \hat{w}_{t-1}, \hat{w}_{t-1})$. Then all the unknowns on the right hand side of equation 42 are uniquely determined by the decomposition of the observed matrix on the left hand side. This matrix $L_{V_{t+1}|w_t, X_t^*}$ is identified up to the permutation of its columns, which implies the identification of $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$ up to the value of x_t^* .

In order to identify how the USV X_t^* changes, it is still useful to reveal its value. As shown in Hu (2007), there are various ways to fix the value of x_t^* . For example, we may normalize the value of x_t^* be the median or another quantile of the distribution $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$. As required in assumption 4, such a quantile needs to be different for a different value of x_t^* . In summary, the conditional density $f_{V_{t+1}|W_t, X_t^*}(\cdot|w_t, x_t^*)$ is identified for any $w_t \in \mathcal{W}_t$.

Step 2: Identification of $\mathbf{f}_{W_{t+1}|W_t, X_t^*}$. We then show that the identification of $f_{V_{t+1}|W_t, X_t^*}$ implies that of $f_{W_{t+1}|W_t, X_t^*}$. Define for any given $w_{t+1} \in \mathcal{W}_{t+1}$, $w_t \in \mathcal{W}_t$, and $w_{t-1} \in \mathcal{W}_{t-1}$,

$$\begin{aligned} \vec{f}_{W_{t+1}, W_t|W_{t-1}, V_{t-2}} &= [f_{W_{t+1}, W_t|W_{t-1}, V_{t-2}}(w_{t+1}, w_t|w_{t-1}, 1), \dots, f_{W_{t+1}, W_t|W_{t-1}, V_{t-2}}(w_{t+1}, w_t|w_{t-1}, J)], \\ \vec{f}_{W_{t+1}|W_t, X_t^*} &= [f_{W_{t+1}|W_t, X_t^*}(w_{t+1}|w_t, 1), \dots, f_{W_{t+1}|W_t, X_t^*}(w_{t+1}|w_t, J)]. \end{aligned}$$

One can show that for any $w_t \in \mathcal{W}_t$

$$\vec{f}_{W_{t+1}|W_t, X_t^*} = \vec{f}_{W_{t+1}, W_t|W_{t-1}, V_{t-2}} (L_{V_{t+1}, w_t|w_{t-1}, V_{t-2}})^{-1} L_{V_{t+1}|w_t, X_t^*}.$$

Therefore, the density $f_{W_{t+1}|W_t, X_t^*}$ is identified.

Step 3: Identification of $\mathbf{f}_{W_t, X_t^*, W_{t-1}, V_{t-2}}$. Moreover, the identification of $f_{V_{t+1}|W_t, X_t^*}$ also implies that of $f_{W_t, X_t^*, W_{t-1}, V_{t-2}}$. Equation 34 also implies

$$f_{V_{t+1}, W_t, W_{t-1}, V_{t-2}} = \sum_{X_t^* \in \mathcal{X}_t^*} f_{V_{t+1}|W_t, X_t^*} f_{W_t, X_t^*, W_{t-1}, V_{t-2}}. \quad (43)$$

Define for any given $w_t \in \mathcal{W}_t$ and $w_{t-1} \in \mathcal{W}_{t-1}$,

$$L_{w_t, X_t^*, w_{t-1}, V_{t-2}} = [f_{W_t, X_t^*, W_{t-1}, V_{t-2}}(w_t, i|w_{t-1}, j)]_{i,j}.$$

Equation 43 is equivalent to⁷

$$L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}} = L_{V_{t+1} | w_t, X^*} L_{w_t, X_t^*, w_{t-1}, V_{t-2}}.$$

Therefore, the identification of $L_{V_{t+1} | w_t, X^*}$ implies that $L_{w_t, X_t^*, w_{t-1}, V_{t-2}}$ is identified as $L_{V_{t+1} | w_t, X^*}^{-1} L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}}$ for any $w_t \in \mathcal{W}_t$. Consequently, the density $f_{w_t, X_t^*, w_{t-1}, V_{t-2}}$ is identified.

Step 4: Identification of $\mathbf{f}_{w_t, X_t^*, w_{t-1}, X_{t-1}^*}$. So far, we have only used the four observations $W_{t+1}, W_t, W_{t-1}, W_{t-2}$. In the last step, we use one more period of the data W_{t-3} to identify the desired joint density $f_{w_t, X_t^*, w_{t-1}, X_{t-1}^*}$.

Replacing t by $t-1$ in the previous three steps implies that the additional information from $\{W_t, W_{t-1}, W_{t-2}, W_{t-3}\}$ or the density $f_{W_t, W_{t-1}, W_{t-2}, W_{t-3}}$ identifies $f_{w_{t-1}, X_{t-1}^*, w_{t-2}, V_{t-3}}$. In turn, we can identify the density $f_{X_{t-1}^* | w_{t-1}, V_{t-2}}$ given the known mapping from W_{t-2} to V_{t-2} .

We then use the identified densities $f_{w_t, X_t^*, w_{t-1}, V_{t-2}}$ and $f_{X_{t-1}^* | w_{t-1}, V_{t-2}}$ to identify $f_{w_t, X_t^* | w_{t-1}, X_{t-1}^*}$. The Markov property implies

$$\begin{aligned} & f_{w_t, X_t^* | w_{t-1}, V_{t-2}}(w_t, x_t^* | w_{t-1}, z) \\ = & \sum_{X_{t-1}^* \in \mathcal{X}_{t-1}^*} f_{w_t, X_t^* | w_{t-1}, X_{t-1}^*}(w_t, x_t^* | w_{t-1}, x_{t-1}^*) f_{X_{t-1}^* | w_{t-1}, V_{t-2}}(x_{t-1}^* | w_{t-1}, z). \end{aligned} \quad (44)$$

Define for any $w_t \in \mathcal{W}_t$, and $w_{t-1} \in \mathcal{W}_{t-1}$

$$\begin{aligned} L_{w_t, X_t^* | w_{t-1}, V_{t-2}} &= \left[f_{w_t, X_t^* | w_{t-1}, V_{t-2}}(w_t, i | w_{t-1}, j) \right]_{i,j}, \\ L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} &= \left[f_{w_t, X_t^* | w_{t-1}, X_{t-1}^*}(w_t, i | w_{t-1}, j) \right]_{i,j}, \\ L_{X_{t-1}^* | w_{t-1}, V_{t-2}} &= \left[f_{X_{t-1}^* | w_{t-1}, V_{t-2}}(i | w_{t-1}, j) \right]_{i,j}, \end{aligned}$$

for $i, j = 1, 2, \dots, J$. Then it is straightforward to show that equation 44 implies

$$L_{w_t, X_t^* | w_{t-1}, V_{t-2}} = L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} L_{X_{t-1}^* | w_{t-1}, V_{t-2}}, \quad (45)$$

where the invertibility of $L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}}$ in equation 35 implies that of $L_{w_t, X_t^* | w_{t-1}, V_{t-2}}$. That mean all the matrices in equation 45 are invertible. Therefore, $L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*}$ is

⁷In fact, $L_{w_t, X_t^*, w_{t-1}, X_{t-2}} = D_{w_t | w_{t-1}, X_t^*} L_{X_t^* | w_{t-1}, X_{t-2}}$.

identified as

$$L_{w_t, X_t^* | w_{t-1}, X_{t-1}^*} = L_{w_t, X_t^* | w_{t-1}, V_{t-2}} L_{X_{t-1}^* | w_{t-1}, V_{t-2}}^{-1}.$$

This results hold for any $w_t \in \mathcal{W}_t$, and $w_{t-1} \in \mathcal{W}_{t-1}$, and therefore, the density $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ is identified. Notice that the identification of $f_{W_{t-1}, X_{t-1}^* | W_{t-2}, V_{t-3}}$ implies that of f_{W_{t-1}, X_{t-1}^*} . Hence, $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ is identified. ■

In summary, the observed density $f_{W_{t+1}, W_t, W_{t-1}, W_{t-2}, W_{t-3}}$ uniquely determines $f_{W_{t+1} | W_t, X_t^*}$ and $f_{W_t, X_t^* | W_{t-1}, X_{t-1}^*}$ under the following assumptions:

1. Assumption 1 (first-order Markov and limited feedback) holds;
2. For any given $w_t \in \mathcal{W}_t$ there exists $(\bar{w}_t, w_{t-1}, \bar{w}_{t-1})$ with $\bar{w}_t \neq w_t$ and $\bar{w}_{t-1} \neq w_{t-1}$ such that the matrices $L_{V_{t+1}, w_t | w_{t-1}, V_{t-2}}$, $L_{V_{t+1}, \bar{w}_t | w_{t-1}, V_{t-2}}$, $L_{V_{t+1}, w_t | \bar{w}_{t-1}, V_{t-2}}$, and $L_{V_{t+1}, \bar{w}_t | \bar{w}_{t-1}, V_{t-2}}$ are all invertible, and that
3. $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j) < \infty$ for all $j \in \mathcal{X}_t^*$ and $k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j_1) \neq k(w_t, \bar{w}_t, w_{t-1}, \bar{w}_{t-1}, j_2)$ for $j_1 \neq j_2$.
4. A known quantile of $f_{V_{t+1} | W_t, X_t^*}(\cdot | w_t, x_t^*)$ is monotonic in x_t^* . Without loss of generality, we normalize x_t^* to be that quantile.

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