Nonparametric Euler Equation Identification and Estimation

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Abstract

We consider nonparametric identification and estimation of pricing kernels, or equivalently of marginal utility functions up to scale, in consumption based asset pricing Euler equations. Ours is the first paper to prove nonparametric identification of Euler equations under low level conditions (without imposing functional restrictions or just assuming completeness). We also propose a novel nonparametric estimator based on our identification analysis, which combines standard kernel estimation with the computation of a matrix eigenvector problem. Our estimator avoids the ill-posed inverse issues associated with nonparametric instrumental variables estimators. We derive limiting distributions for our estimator and for relevant associated functionals. A Monte Carlo shows a satisfactory finite sample performance for our estimators.

JEL Codes: C14, D91, E21, G12. Keywords: Euler equations, marginal utility, pricing kernel, Fredholm equations, integral equations, nonparametric identification, asset pricing.

1 Introduction

The optimal intertemporal decision rule of an economic agent can often be characterized by first-order condition Euler equations. These equations are fundamental objects that appear in numerous branches of economics, in particular in the literatures on consumption, on savings and asset pricing, on labor supply, and on investment. Many empirical studies of dynamic optimization behaviors rely on the estimation of Euler equations. One of the original motivations of the generalized method of moments (GMM) estimator proposed by Hansen and Singleton (1982) was estimation of rational expectations based Euler equations associated with consumption based asset pricing models. In this paper we study the nonparametric identification and estimation of such Euler equations.

To fix ideas, consider a familiar consumption based asset pricing Euler equation (e.g. Cochrane 2001)

$$bE[g(C_{t+1}, V_{t+1})R_{t+1} | C_t, V_t] = g(C_t, V_t), \text{ almost surely (a.s.)} \tag{1}$$

where $b$ is the subjective discount factor, $C_t$ is consumption at time $t$, $V_t$ is a vector of other economic variables such as durables or lagged consumption (for habits) that might affect utility, $R_t$ is the gross return of an asset, and $g$ is the time homogeneous marginal utility function of consumption. Equation (1) is the first order condition that equates in real terms the marginal cost of an extra unit of the asset, purchased today, to the expected marginal benefit of the extra payoff received tomorrow.$^1$

Our work is the first to establish nonparametric point identification of the marginal utility function $g$, and by implication of the pricing kernel function $M$ (see below), under low level assumptions.$^2$ We also provide a novel nonparametric estimator based on this identification analysis, which combines standard kernel estimation with the computation of a matrix eigenvector problem. Our estimator overcomes the ill-posed inverse problem that affects existing nonparametric instrumental variables based estimators.

We take the primitives of the Euler equation to be the marginal utility function $g$, defined up to an arbitrary sign and scale normalization, and the discount factor $b$. The (nonparametric) identified set for the Euler equation is defined to be the set of all $(g, b) \in \Theta = \mathcal{G} \times (0, 1)$, for a suitable parameter space $\mathcal{G}$, that satisfy equation (1), given the true joint distribution of the data (see Tamer 2010 for a review of set identification definitions). A model is defined to be globally identified if the identified set only consists of one element.

$^1$For a formal derivation of this Euler equation, with internal or external habits, see the Appendix.

$^2$This paper is a merged revision of two earlier working papers: Lewbel, Linton and Srisuma (2011) and Escanciano and Hoderlein (2010, 2012). In particular, Escanciano and Hoderlein (2010) first used Perron-Frobenious theory for identification of marginal utilities. Some recent papers by others that establish related identification results cite these earlier versions of our paper as prior knowledge. See the next section for details.
In this paper we first show that the Euler equation is partially identified, with a finite identified set for the discount factor and an identified set for marginal utilities that is the union of finite dimensional spaces. This implies that the discount factor is also locally identified (in the sense of Fisher 1966, Rothenberg 1971 and Sargan 1983), meaning that \( b \) is nonparametrically identified within a parameter space that equals a neighborhood of the true value. We then show that if the class of utility functions is restricted to be monotone, which is a natural economic restriction, then the Euler equation model is, nonparametrically, globally point identified.

Having established identification, we next propose a novel nonparametric kernel estimator for the marginal utility function and discount factor based on our identification arguments. We provide asymptotic distribution theory for the discount factor, the marginal utility function, and for semi-parametric functionals of the marginal utility function such as the Average Relative Risk Aversion (ARRA) parameter defined below.

In the empirical asset pricing literature, the Euler equation (1) is traditionally written as

\[
E [M_{t+1} R_{t+1} \mid C_t, V_t] = E \left[ \frac{g(C_{t+1}, V_{t+1}) R_{t+1}}{g(C_t, V_t)} \mid C_t, V_t \right] = 1,
\]

where \( M_{t+1} = \frac{b g(C_{t+1}, V_{t+1})}{g(C_t, V_t)} \) is the time \( t+1 \) pricing kernel or Stochastic Discount Factor (SDF). Then, the pricing equation for asset \( R \) can be cast in the form of excess returns

\[
E [M_{t+1} (R_{t+1} - R_{0t+1}) \mid C_t, V_t] = E \left[ \frac{b g(C_{t+1}, V_{t+1}) (R_{t+1} - R_{0t+1})}{g(C_t, V_t)} \mid C_t, V_t \right] = 0,
\]

(2)

where \( R_{0t} \) denotes the return from the risk-free asset. Equation (2) is a conditional moment restriction that forms the basis of moments based estimation. In a parametric model, \( g \) (and hence \( M_t \)) is assumed known up to finite-dimensional parameters; prominent examples include Hall (1978), Hansen and Singleton (1982), Dunn and Singleton (1986), and Campbell and Cochrane (1999), among many others. Euler equations have also been specified semiparametrically, e.g., Chen and Ludvigson (2009) and Chen, Chernozhukov, Lee and Newey (2014).

Nonparametric estimators of equation (2) and similar models (taking the form of nonparametric instrumental variables models) have been proposed, by, e.g., Gallant and Tauchen (1989), Chapman (1997), Newey and Powell (2003), Ai and Chen (2003) and Darolles, Fan, Florens, and Renault (2011). However, in these applications identification is assumed rather than proved, by way of high level completeness assumptions. These models have the structure of Fredholm equations of the first kind (also called Type I equations). Solving these types of equations involves ill-posed inverse problems that can be severe, and as a result, fully nonparametric estimators of \( M_{t+1} = M(C_{t+1}, V_{t+1}, C_t, V_t) \) based on (2) can have very slow convergence rates and possibly unstable inference.

In contrast, we start by writing the pricing kernel problem in the form of equation (1) instead of equation (2), thereby estimating \( g \) instead of \( M \). The advantage is that equation (1) takes the form of
a Fredholm linear equation of the second kind (or Type II equation). As a result, unlike equation (2),
the solution of equation (1) has a well-posed generalized inverse, leading to much better asymptotic
properties for inference. In particular, in solving equation (1), a candidate discount factor $b$ and
associated marginal utility function $g$ is characterized as an eigenvalue-eigenfunction pair of a certain
conditional mean operator. Under the mild assumption that this operator is compact, a classical
result (see e.g. Kress (1999)) ensures that the number of eigenvalues is countable. The behavioral
restriction that $b < 1$ reduces this set to a finite number, leading to our finite set identification result
and hence to local identification for the discount factor. To obtain global point identification of
$b$ and $g$, we impose the additional behavioral restriction that utility is increasing in consumption,
which implies that the function $g$ is positive. Applying an infinite-dimensional extension of the
Perron-Frobenius theorem (see Krein and Rutman 1950) yields uniqueness of a positive eigenvalue-
eigenfunction pair, which then provides nonparametric point identification.

Following this identification argument, we propose a new nonparametric estimator for the mar-
ginal utility function $g$ and discount factor $b$. The estimator is based on standard kernel estimation of
a sample analogue of (1), which with finite data replaces the problem of solving for an eigenfunction
with the simpler problem of solving for a standard finite-dimensional matrix eigenvector. No nu-
merical integration or optimization is required, making the estimator straightforward to implement
(and numerically practical to bootstrap). We establish our estimator’s limiting distribution under
standard conditions, which are simpler than those associated with estimators that solve Type-I ill-
posed inverse problems, such as nonparametric instrumental variables. Our expansions show that, in
contrast to nonparametric problems leading to Type-I equations, nonparametric inference on $g$ in our
Type-II equation is to a large extent mathematically equivalent to inference on a standard conditional
mean function, and in particular has comparable rates of convergence to ordinary nonparametric re-
gression. Although our assumptions are standard, both our identification and asymptotic theory
entail machinery that is novel in the econometrics literature, applying an infinite-dimensional ex-
tension of Perron-Frobenius theory to a type II Fredholm equation (see the next section for details
comparing our results to the literature).

In addition to the pricing kernel $M_{t+1}$, another functional of the marginal utility function $g$ that
is of interest to estimate is the Arrow-Pratt coefficient of Relative Risk Aversion, and its average
value, $RRA$ and $ARRA$, given respectively by

$$RRA(c, v) = \frac{-c \partial g(c, v) / \partial c}{g(c, v)} \quad \text{and} \quad ARRA = E[RRA(C_t, V_t)].$$

We establish asymptotic normality of a nonparametric estimator of the $ARRA$. Given our estimates
of $g(c, v)$, we also provide tests of whether $g$ is independent of $v$, thereby testing whether lagged
consumption (or any other potential covariates $v$ such as durables consumption) affects the pricing
kernel. These tests are based on semiparametric functionals of $g$, which are asymptotically normal under the same type of regularity conditions we use to establish asymptotics for the ARRA.

The asymptotic theory we present in this paper is based on a random sampling framework of household-level consumption and asset data. Our theory can be easily extended to accommodate dependent and heterogeneous data as long as we can establish appropriate uniform consistency of certain nonparametric kernel density estimators and apply CLT to some of their functionals. Therefore our existing results would not be suitable to analyze aggregated consumption data that appear to be non-stationary (particularly non-recurrent time series). Otherwise, our estimates allow applied researchers to engage in a nonparametric consumption-based asset pricing theory. For example, Abbott and Gallipoli (2018) illustrate the applicability of our methods by using them to estimate nonparametrically human wealth and permanent income with PSID data.

The rest of the paper is organized as follows. After a literature review in Section 2, we provide sufficient conditions for partial identification and point identification in Section 3. We propose our kernel-type estimator in Section 4, and we investigate its asymptotic properties in Section 5. In Section 6 we describe how our asymptotic theory applies to functionals of $g$, and give some examples. We report the results of a Monte Carlo experiment in Section 7. Section 8 concludes. An Appendix contains the derivation of the Euler equation and the mathematical proofs of the main results.

2 Literature Review

The forerunners of our research are the papers by Gallant and Tauchen (1989) and Chapman (1997), who estimate nonparametrically the marginal utilities and the pricing kernel, respectively, from the Euler equation by sieves, using the moment restriction (2) (i.e. using a Type I Fredholm equation). These papers did not investigate identification, nor impose the positivity of marginal utilities, and the asymptotic properties of their nonparametric estimators were not established.

Nonparametric instrumental variables is a leading example of estimation based on a Type I Fredholm equation, yielding associated ill-posed inverse problems on estimation. Newey and Powell (2003) note that assuming statistical completeness (a high level assumption) is essentially the same as just assuming identification of this type of model. Other related examples of nonparametric and semiparametric ill-posed inverse estimation problems include Carrasco and Florens (2000), Ai and Chen (2003), Hall and Horowitz (2005), Chen and Pouzo (2009), Chen and Reiss (2010), Darolles, Fan, Florens and Renault (2011) and, more recently, Cai, Ren and Sun (2015). A particularly relevant example is Chen and Ludvigson (2009), who studied identification and estimation of a semiparametric specification of the Type-I equation (2). Their model assumes $g$ has the semiparametric form $g(C_t, V_t) = C_t^\eta h(V_t)$, where $\eta$ is a constant that determines risk aversion and $h$ is an unknown func-
tion of current and lagged values of $C_{t}/C_{t-1}$ representing habits. Virtually all parametric estimators of the asset pricing model, going back to Hansen and Singleton (1982) and including Dunn and Singleton (1986), and Campbell and Cochrane (1999), use the form of equation (2) rather than equation (1).

Many parametric rational expectations models that focus on utility or production rather than asset pricing do estimation in the form of equation (1). Early examples include Hall (1978) and Mankiw (1982) (though see Lewbel 1987 for a critique). This earlier work does not appear to recognize the theoretical integral equation advantages of casting the model in the form of equation (1). Anatolyev (1999) recognizes that this form is a Type II Fredholm equation and provides a numerical method for estimating Euler equations that makes use of this structure, but he does not consider identification or inference. We believe our paper is the first to make explicit use of this Type II Fredholm structure for identification and inference. An and Hu (2012) exploit the nature of a type II Fredholm equation to identify and estimate a measurement error rather than an Euler equation model, but they cite our working paper as prior knowledge.

Our proof of global identification makes use of extensions of the classical Perron-Frobenius theorem that positive matrices have a unique positive eigenvalue that corresponds to a unique positive eigenvector. In particular, we apply a theorem of Krein and Rutman (1950), which extends Perron-Frobenius to compact operators in Banach spaces. See, e.g., Schaefer (1974) and Abramovich and Aliprantis (2002) for a comprehensive presentation of this theory.

Versions of Perron-Frobenius have been used before in Euler equation models, though we believe we are the first to use this machinery of infinite-dimensional Perron-Frobenius theory for nonparametric identification and inference of Euler equations (see Escanciano and Hoderlein 2010). There is, however, some closely related work. Hansen and Scheinkman (2009, 2012, 2013) first used Perron-Frobenius theory for a different problem of identification than ours in a continuous-time setting, using Markov theory. In our notation, they give conditions for identification of the positive eigenfunction and eigenvalue of the operator $\phi \rightarrow E[M_{t+1}\phi(C_{t+1},V_{t+1}) \mid C_{t},V_{t}]$, assuming that the SDF $M_{t+1}$ is known. In contrast, we solve the also fundamental problem of showing that $M_{t+1}$ itself is identified, by obtaining identification of $b$ and $g$. Christensen (2015, 2017) applies identification results, based in part on our earlier working papers, to a discrete version of Hansen and Scheinkman (2009). Also recently, Ross (2015) applies the classical finite-dimensional Perron-Frobenius theorem to identify the pricing kernel and the natural probability distribution from state prices.

Regarding identification, perhaps the closest work to ours is Chen, Chernozhukov, Lee and Newey (2014). Although their paper mainly concerns local nonparametric identification, in their Euler equation application they consider a semiparametric rather than a nonparametric model like ours. Specifically, their model is the same functional form as Chen and Ludvigson (2009) described above,
but allowing for a more general conditioning set. They cite the working paper versions of our paper as prior knowledge, including the Perron-Frobenius application of Escanciano and Hoderlein (2010). They first use completeness conditions to identify the parametric \( RRA \) and then use Perron-Frobenius to identify the role of habits. In contrast, we do not require a constant \( RRA \) or require completeness conditions for identification. Thus, the setting and identification approaches of this paper and those of Chen et al. (2014) are quite different.

An alternative to our kernel based estimation would be the use of sieves. Although we focus on kernel estimates, our asymptotic theory is developed in a way that can be easily adapted to other nonparametric estimation methods, including sieves (e.g. splines) and local polynomial methods. Nonparametric sieve estimation of eigenvalue-eigenvector problems for self-adjoint operators is extensively discussed in Chen, Hansen and Sheinkman (2000, 2009), Darolles, Florens and Gouriéroux (2004) and Carrasco, Florens and Renault (2007), among others.\(^3\) However, their results cannot be applied to our model, since in our case the associated operator is not self-adjoint. Christensen (2017) proposes a nonparametric sieve estimator for the discrete Markov setting of Hansen and Sheinkman (2009), establishing asymptotic normality of the eigenvalue estimate and smooth functionals of it. See also Gobet, Hoffmann and Reiss (2004) for sieve estimation of eigenelements in diffusion models. As noted earlier, sieve estimation has more directly been applied to nonparametric and semiparametric versions of equation (2) going back to Gallant and Tauchen (1989). In comparison, our kernel based estimator has several advantages as summarized in the previous section, mainly attributable to our method of exploiting the well-posedness of equation (1). In particular, with our methods we obtain novel asymptotic distribution theory for functionals of the nonparametric utility, such as the \( ARRA \) functional. This asymptotic theory is of independent interest and has wide applicability in other situations where type-II equations arise.

3 Identification

Since our goal is the study of Euler equations, we shall take as primitives the pair \((g,b) \in \Theta \equiv \mathcal{G} \times (0,1)\), where \(\mathcal{G}\) denotes the parameter space of marginal utility functions, which satisfies some conditions below. From equation (1) it is clear that, for a given \(b\), the Euler equation cannot distinguish between \(g\) and \(h\) if there exists some constant \(k_0 \in \mathbb{R}\) such that \(g = k_0 h\) a.s., so a scale and a sign normalization must be made. For the moment we shall assume there is just one asset, and we denote its rate of return by \(R_t\). We later discuss how information from multiple assets can be

\(^3\)Section 2.5 in Carrasco et al. (2007) discusses generic methods for estimation of the singular values and associated eigenfunctions. This problem is different from estimating the principal eigenvalue and eigenfunction of a non-selfadjoint operator, but our kernel estimation strategy follows closely their logic.
used to aid identification. As seen in the previous section, for each period $t$, $C_t$ is consumption and $V_t$ is (possibly a vector of) other economic variable(s).

**Definition.** Let $S, S' \subseteq \mathbb{R}^d$ denote the supports of $(C_t, V_t)$ and $(C_{t+1}, V_{t+1})$ respectively. Let $\mu$ be a probability measure, with support $S_\mu \subseteq S \cap S'$, and let $L^2$ denote the Hilbert space $L_2(S_\mu, \mu)$ of (equivalence classes of) square $\mu$-integrable functions equipped with the inner product $\langle g, f \rangle = \int g f d\mu$ and the corresponding norm $\|g\|^2 = \langle g, g \rangle$ (we drop the domain of integration for simplicity of exposition).

Let $\mathcal{M} \subseteq L^2$ be a linear subspace, and define the linear operator $A : (\mathcal{M}, \|\cdot\|) \rightarrow (\mathcal{M}, \|\cdot\|)$ by

$$Ag(c, v) = E[g(C_{t+1}, V_{t+1}) R_{t+1} | C_t = c, V_t = v]. \quad (3)$$

We assume that $Ag$ is well-defined and $Ag \in \mathcal{M}$ for $g \in \mathcal{M}$. Examples of $\mu$ and $\mathcal{M}$ are given below. With our notation, (1) can be written in a compact form as $bAg = g$. The parameter space for $g, \mathcal{G}$, will be a subset of $\mathcal{M}$ incorporating normalization restrictions. Marginal utilities may not have finite moments around zero (where they may diverge). To overcome this problem, by suitable redefinition of $g$ we can rewrite equation (1) in the form

$$bE[C'g(C', V') (C/C') R' | C, V] = Cg(C, V). \quad (4)$$

This reparameterizes the problem in terms of $Cg(C, V)$, which under natural economic assumptions is bounded; see Lucas (1978). This identity also gives an alternative way to estimate the marginal utility function and other objects of interest, which we shall discuss further below. More generally, we could apply parametrizations with powers of consumption if necessary (i.e. with $(C')^\delta g(C', V') (C/C')^\delta$ and $C^\delta g(C, V)$ in the left and right hand side of (4), respectively). The parameter $\delta$ is chosen by the researcher (i.e. it is not estimated). Allowing for $\delta > 1$ is useful to accommodate CRRA utilities with risk aversion larger than one.

We introduce the assumption of correct specification and a formal definition of identification.

**Assumption S.** There exists $(g, b) \in \Theta \equiv \mathcal{G} \times (0, 1)$, $g \neq 0$, satisfying equation (1).

**Definition 1.** Given the joint distribution of $(R_{t+1}, C_{t+1}, V_{t+1}, C_t, V_t)$, the Euler equation is non-parametrically identified if there is a unique $(g, b) \in \Theta$ that satisfies equation (1). When the solution is unique we denote it by $\theta_0 \equiv (g_0, b_0)$.

**Definition 2.** Given the joint distribution of $(R_{t+1}, C_{t+1}, V_{t+1}, C_t, V_t)$, the identified set, denoted by $\Theta_0$, consists of elements in $\Theta$ where each $(g, b) \in \Theta_0$ satisfies equation (1) with $g \neq 0$. The sets $B_0 = \{b \in (0, 1) : \text{there is } g \in \mathcal{G} \text{ such that } (g, b) \in \Theta_0\}$ and $\mathcal{G}_0 = \{g \in \mathcal{G} : \text{there is } b \in (0, 1) \text{ such that } (g, b) \in \Theta_0\}$ are, respectively, the identified sets for $b$ and $g$. 
Therefore, the Euler equation is point identified if $\Theta_0$ is a singleton. To provide some insights on our identification and estimation strategies we consider first the case where $A$ in (3) has a finite-dimensional range. In this case, we can write

$$Ag(\cdot) = \sum_{i=1}^{I} L_i(g) \phi_i(\cdot),$$

(5)

for a set of functions $\{\phi_i\}$ that span the range of $A$, $\mathcal{R}(A) = \{f \in \mathcal{M} : \exists g \in \mathcal{M}, Ag = f\}$, and linear operators $L_i(g), i = 1, \ldots, I$. This situation arises, for example, when the support $S$ is discrete and finite. Under (5), any potential solution of (1) has to necessarily the form $g(\cdot) = \sum_{i=1}^{I} \beta_i \phi_i(\cdot)$ for a vector $\beta = (\beta_1, \ldots, \beta_I)$ satisfying the Euler equation

$$\sum_{i=1}^{I} \sum_{j=1}^{I} L_i(\phi_j) \beta_j \phi_i(c, v) = b^{-1} \sum_{i=1}^{I} \beta_i \phi_i(c, v).$$

In turn, this is the case for the solution, provided it exists, of

$$\sum_{j=1}^{I} \beta_j L_i(\phi_j) = b^{-1} \beta_i \quad 1 \leq i \leq I.$$

Therefore, $\beta$, i.e. $g$, and $b^{-1}$ are identified as any eigenelement of the $I \times I$ matrix $(L_i(\phi_j))_{i,j}$, with $b \in (0,1)$. In general, we may have more than one such eigenelement, i.e., we may have partial identification. In any case, the number of eigenvectors $\beta$ and eigenvalues is bounded by $I$, so we have a finite identified set.

As we shall show, the previous arguments extend to the general infinite-dimensional case replacing the finite-dimensionality of $\mathcal{R}(A)$ by the compactness of $A$. A linear operator $A$ is compact if it transforms bounded sets into relatively compact sets (relatively compact sets in $\mathcal{M}$ are those whose closure is compact). The compactness assumption is standard in the literature and is useful for both identification and for obtaining asymptotics of continuous functionals of $g$. Note, however, that if there are overlapping elements in $(C_{t+1}, V_{t+1})$ and $(C_{t}, V_{t})$ compactness rules out the case $\mathcal{M} = L^2$; see Carrasco, Florens and Renault (2007, Example 2.5, pg. 22). We could deal with the lack of compactness of $A$ on the whole $L^2$ by conditioning on (i.e. fixing) the overlapping components, as is common in the literature (see e.g. Blundell, Chen and Kristensen, 2007, pg. 1629). From the identification point of view there is little loss of generality by following this “conditioning” approach, however, for deriving asymptotics compactness is convenient, as it guarantees that inference will be based on well-posed generalized inverses (see the discussion at the end of this section).

**Assumption C.** $A : (\mathcal{M}, \| \cdot \|) \to (\mathcal{M}, \| \cdot \|)$ is a compact operator.

Let $\mathcal{G} \subseteq \{g \in \mathcal{M} : \|g\| = 1, \ g(c_0, v_0) > 0, \ (c_0, v_0) \in S\}$ be the parameter space for $g$. 

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Theorem 1. Suppose that Assumptions S and C hold. Then, \(B_0\) is a finite set and \(G_0\) is a finite union of finite-dimensional subsets intersected with \(G\).

Theorem 1 shows that without further assumptions the Euler equation is partially identified, with \(b\) identified up to a finite set corresponding to eigenvalues larger than one, and \(g\) is identified up to a corresponding set of eigenfunctions. The discount factor \(b\) is also locally identified, meaning that for any \(b \in B_0\) there is an open neighborhood of \(b\) that does not contain any other element in \(B_0\). Essentially, compactness of \(A\) ensures that \(B_0\) is at most countable, and the economic restriction that discount factors lie in \((0,1)\) ensures that \(B_0\) is finite.

The identified set without additional economic restrictions can be further reduced if there are multiple assets. If there are \(J\) assets, then there are \(J\) Euler equations. Applying Theorem 1 to each asset, gives an identified set for each, and the true \((g, b)\) must lie in the intersection of these identified sets. One might further shrink the identified set by imposing the restriction that \(bg(C_{t+1}, V_{t+1})R_{t+1} - g(C_t, V_t)\) is uncorrelated with all variables in the information set at time \(t\), not just measurable functions of \((C_t, V_t)\).

Assumptions S and C do not suffice for point identification in general. We consider now a shape restriction on marginal utilities, which is a common behavioral assumption that is satisfied for common parametric specifications of utility. Specifically, we impose the assumption that marginal utilities are positive. Let

\[ \mathcal{P} \equiv \{g \in \mathcal{M} : g \geq 0 \text{ } \mu - a.s.\} \tag{6} \]

denote the subset of nonnegative functions in \(\mathcal{M}\), and let \(\mathcal{P}^+ \equiv \{g \in \mathcal{M} : g > 0 \text{ } \mu - a.s.\}\) denote the subset of strictly positive functions, which is assumed to be non-empty. The assumption is then:

Assumption I. \(Ag \in \mathcal{P}^+ \text{ when } g \in \mathcal{P} \text{ and } g \neq 0\).

Assumption I is a mild condition that extends the classical assumption of a positive matrix in the Perron-Frobenius theorem to an infinite-dimensional setting, see Abramovich and Aliprantis (2002, Chapter 9) and Schaefer (1974). With our shape and normalization restrictions the parameter space is \(G \subseteq \{g \in \mathcal{P} : \|g\| = 1\}\). Note that Assumption S and I imply that \(g > 0 \text{ } \mu - a.s.\) for all \(g \in G_0\) because \(g = bAg > 0\).

Theorem 2. Let Assumptions S, C and I hold. Then, \((g, b) \in G \times (0,1)\) is point identified.

Identification can be established under weaker conditions than those of Theorem 2, however, we do not pursue these conditions here because the stronger conditions of Theorem 2 will facilitate our later asymptotic inference results. See Escanciano and Hoderlein (2012) for these weaker conditions. Our proof of Theorem 2 also shows that \(b = 1/\rho(A)\), where \(\rho(A)\) is the spectral radius of \(A\) (see the
Appendix for a definition of the spectral radius of a linear bounded operator). Following Escanciano and Hoderlein (2012) a key sufficient condition for identification of $g$ is that $A$ is irreducible; see Abramovich and Aliprantis (2002, Chapter 9) for a definition of irreducibility in a general setting.

We could consider other sufficient conditions that replace conditions on $A$ by conditions on a power of $A$, i.e. we could require that Assumptions C and I hold for $A^n$, for some $n \geq 1$ (see Escanciano and Hoderlein (2012)). It is hard to interpret these conditions, however, in a possibly non-Markovian environment, so we do not pursue them here. It is also likely that the Euler Equation is overidentified under the conditions of Theorem 2, since as noted earlier we could exploit additional information coming from multiple assets, or from uncorrelatedness with other data in the information set at time $t$.

For illustration, we consider the following examples of $\mu$ and $M$, which lead to lower level conditions for identification by Theorem 2. Assume for simplicity that $V_{t+1}$ and $V_t$ are empty, and denote by $f(c',c)$, $f'(c')$ and $f(c)$ the joint and marginal densities of $(C_{t+1}, C_t)$, respectively. Then, taking $\mu$ the distribution of $C_t$ and $M$ equals to $\mathcal{L}^2$, the operator equation $bAg = g$ can be written as

$$b \int k(c',c)g(c')d\mu(c') = g(c),$$

where $k(c',c) = r(c',c)f(c',c)/[f(c')f(c)]$ and $r(c',c) = E [R_{t+1}|C_{t+1} = c',C_t = c]$ a.s. Then, it is well known that Assumption C holds if $\int \int k^2(c',c)d\mu(c')d\mu(c) < \infty$, while Assumption I holds if $k(c',c) > 0 \mu \times \mu$ a.s. (or equivalently $r(c',c) > 0$ and $f(c',c) > 0$ almost everywhere). Note that these assumptions do not require consumption to be stationary (i.e. $f \neq f'$ is allowed). Weaker sufficient conditions for Assumption I for this choice of $\mu$ and $M$ are given in Escanciano and Hoderlein (2012).

We close our study of identification with a discussion on the degree of ill-posedness of our nonparametric problem. Assumption S implies that the operator $L = bA - I$ is not one-to-one, as $Lg = 0$ and $g \neq 0$. Therefore, solving the Euler equation (1) is an ill-posed problem (see e.g. Carrasco, Florens and Renault 2007, Section 7). However, unlike in ill-posed Type-I equations, the ill-posedness in our Type-II equation is moderate, with stable solutions. Formally, the operator $L$, although not invertible, has a continuous (i.e. bounded) Moore-Penrose pseudoinverse, which is denoted by $L^\dagger$; (see Engl, Hanke and Neubauer 1996, p. 33). To see this, note that the compactness of $A$ and the Second Riesz Theorem, see e.g. Theorem 3.2 in Kress (1999, p. 29), imply that the range of $L$, $\mathcal{R}(L) = \{f \in \mathcal{L}^2 : \exists s \in \mathcal{L}^2, Ls = f\}$, is closed. This in turn implies that $L^\dagger$ is a continuous operator by Proposition 2.4 in Engl et al. (1996). It is in this precise sense that our problem leads to well-posed rather than ill-posed generalized inverses. This property of our nonparametric problem, which results from considering Type-II equations rather than Type-I equations, has important implications for inference. For example, in the next sections we obtain rates of convergence for estimation of $g$ that are the same as those of ordinary nonparametric regression.

11
4 Estimation from Individual level-data

Our estimation strategy follows the identification strategy described above. For estimation we assume that we have a random sample of household-level data \( \{(R_{t_i+1}, C_{t_i+1}, V_{t_i+1}, C_{t_i}, V_{t_i})\}_{i=1}^n \) for \( n \) households, with possibly overlapping time periods \( t_1 \leq t_2 \leq \cdots \leq t_n \). To simplify notation denote \( W_i = (R_i', C_i', V_i', C_i, V_i) \equiv (R_{t_i+1}, C_{t_i+1}, V_{t_i+1}, C_{t_i}, V_{t_i}) \), where \( V_i = (V_{1i}, \ldots, V_{ti,i}) \) and \( V_i' = (V_{1i}', \ldots, V_{ti,i}') \) with \( \ell = \ell_1 + 1 \). We assume the data, \( \{W_i\}_{i=1}^n \), are independent and identically distributed (iid), generated with respect to an underlying parameter \( \theta_0 \equiv (g_0, b_0) \in \Theta \). Our asymptotic results can all be straightforwardly extended to non-iid data using the uniform rate results of Andrews (1995); see the discussion after Theorem 4 for further details. We shall henceforth assume that Assumptions S, C and I hold, so that \( \theta_0 \) is point-identified. Particularly, we consider \( g_0 \in \mathcal{G} \subseteq \{ g \in \mathcal{P} : \|g\| = 1 \} \).

Let the vector \( W = (R', C', V', C, V) \) have the same distribution as \( (R_i', C_i', V_i', C_i, V_i) \). We assume that the vector \( W \) is continuously distributed (the discrete case is simpler). As in the example above, we denote the Lebesgue density of \( (C, V) \) by \( f \) and that of \( (C', V') \) by \( f' \). We consider the setting described in the identification section where \( \mu \) is a probability measure with Lebesgue density \( f_\mu \) and support \( S_\mu \subseteq S \cap S' \). Henceforth, \( g \) and \( b \) denote generic elements in \( \mathcal{G} \) and \((0, 1)\), respectively.

Define the Nadaraya-Watson (NW) kernel estimator of the operator \( A \) at \( g \) as follows,

\[
\widehat{A}g(c, v) = \frac{1}{n} \sum_{i=1}^n g' R_i' \phi_i(c, v),
\]

where, for \( i = 1, \ldots, n \), \( g'_i \equiv g (C_i', V_i') \), \( \phi_i(c, v) = K_{hi}(c, v) / \hat{f}(c, v) \), while for \( v = (v_1, \ldots, v_{\ell_i}) \),

\[
\hat{f}(c, v) = \frac{1}{n} \sum_{i=1}^n K_{hi}(c, v),
\]

and

\[
K_{hi}(c, v) = h^{-\ell} K \left( \frac{c - C_i}{h} \right) \prod_{j=1}^{\ell_i} K \left( \frac{v_j - V_{ij}}{h} \right).
\]

Here, \( K \) is a univariate kernel function and \( h \equiv h_n \) is a possibly stochastic bandwidth. Note that contrary to \( A \), the operator \( \widehat{A} \) has a finite-dimensional closed range (that is spanned by the functions \( \phi_i(c, v), i = 1, \ldots, n \)). Therefore, similar to our discussion of identification in Section 3, the number of eigenvalues and eigenfunctions of \( \widehat{A} \) is finite and bounded by \( n \), and they can be computed by solving a linear system. Indeed, any eigenfunction \( \hat{g}(c, v) \) of \( \widehat{A} \) necessarily has the form \( n^{-1} \sum_{i=1}^n \hat{\beta}_i \phi_i(c, v) \), for some coefficients \( \hat{\beta}_i, i = 1, \ldots, n \), satisfying for its corresponding eigenvalue \( \hat{\lambda} \) the equation

\[
\frac{1}{n^2} \sum_{i=1}^n \sum_{j=1}^n \hat{\beta}_j \phi_j(C_i', V_i') R_i' \phi_i(c, v) = \hat{\lambda} \frac{1}{n} \sum_{i=1}^n \hat{\beta}_i \phi_i(c, v).
\]
A solution to this eigenvalue problem exists if, for all \( i = 1, \ldots, n \),

\[
\frac{1}{n} \sum_{j=1}^{n} \beta_j \phi_j(C'_i, V'_i) R'_i = \hat{\lambda} \hat{\beta},
\]

which in matrix notation can be simply written as

\[
\hat{A}_n \hat{\beta} = \hat{\lambda} \hat{\beta},
\]

where \( \hat{A}_n \) is an \( n \times n \) matrix with \( ij \)-th element \( a_{ij} = \phi_j(C'_i, V'_i) R'_i \), and \( \hat{\beta} = (\hat{\beta}_1, \ldots, \hat{\beta}_n)^\top \) (henceforth, \( v^\top \) denotes the transpose of \( v \)). Thus, let \( \hat{\lambda} \) denote the largest eigenvalue in modulus of \( \hat{A}_n \) and \( \hat{\beta} = (\hat{\beta}_1, \ldots, \hat{\beta}_n)^\top \) its corresponding eigenvector. Our estimators for \( b_0 \) and \( g_0 \) are, respectively,

\[
\hat{b} = 1/\hat{\lambda} \quad \text{and} \quad \hat{g}(c, v) = n^{-1} \sum_{i=1}^{n} \hat{\beta}_i \phi_i(c, v). \tag{7}
\]

Marginal utilities are identified up to scale and we consider a normalization that is convenient for the asymptotic analysis. Define \( \|g\|_n^2 = \int g^2 d\mu_n \), where \( \mu_n \) is an empirical measure that converges weakly to \( \mu \), and is independent of the data. For example, if \( \mu \) is the distribution of \( (C_i, V_i) \), then we can take as \( \mu_n \) the empirical measure of an independent bootstrap sample of \( \{(C_i, V_i)\}_{i=1}^{n} \), say \( \{(C_i^*, V_i^*)\}_{i=1}^{n} \), so \( \|g\|_n^2 = \sum_{i=1}^{n} g^2(C_i^*, V_i^*) / n \). Then, \( \hat{\beta} \) is normalized so that \( \hat{g} \) satisfies \( \|\hat{g}\|_n = 1 \), or equivalently \( \hat{\beta}^\top \hat{\Omega} \hat{\beta} = 1 \), where \( \hat{\Omega} \) is the \( n \times n \) matrix with entries

\[
\omega_{ij} = \frac{1}{n^3} \sum_{l=1}^{n} \phi_l(C_l^*, V_l^*) \phi_j(C_l^*, V_l^*).
\]

This normalization is useful to show under weak conditions that \( \|\hat{g}\| \rightarrow 1 \) a.s. (since \( 1 = \|\hat{g}\|_n = \|\hat{g}\| + o(1) \) a.s. by the Law of Large Numbers). Alternative normalizations and methods of proof with stronger conditions could be used. For example, we could normalize both \( \hat{g} \) and \( g_0 \) to have unit norm with respect to a known distribution (e.g. standard Gaussian). Also, we impose the sign normalization \( \langle \hat{g}, 1 \rangle_n \equiv n^{-1} \sum_{i=1}^{n} \hat{g}(C_i^*, V_i^*) > 0 \). The estimator \( \langle \hat{g}, \hat{b} \rangle \) can be easily obtained with any statistical package that computes eigenvalues and eigenvectors of matrices. There are also efficient algorithms for the computation of the so-called Perron-Frobenius root \( \hat{\lambda} \), see e.g. Chanachana (2007).

Notice that under very mild conditions the matrix \( \hat{A}_n \) itself satisfies the classic conditions of the Perron-Frobenius theorem, which guarantees that \( \hat{b} = \rho^{-1}(\hat{A}_n) \) and \( \hat{\beta} \) is the only eigenvector of \( \hat{A}_n \) with positive entries. That is, in this case we also have identification in finite samples. For example, for strictly positive kernels and strictly positive gross returns, \( \hat{A}_n \) has strictly positive entries, which then implies a positive estimator \( \hat{g}(c, v) > 0 \) and a positive discount factor \( \hat{b} \) with probability one for a fixed \( n \geq 1 \). For higher-order kernels we can take the positive part of the estimator as usual, and apply Perron-Frobenius for a sufficiently large \( n \).
The easiest way to consider simultaneously different assets in our estimation strategy is to consider individual estimates of the marginal utility for each asset by the method above and then combine the resulting estimators to reduce the variance; see e.g. Chen, Jacho-Chavez and Linton (2016).

5 Asymptotic Theory

In this section we provide conditions for the consistency and limiting distribution theory of our estimators as defined in the previous section, under a random sampling framework. We give high-level conditions that allow for general estimators of \( A \), and provide in the Appendix low-level conditions for our leading kernel-type estimator. For a generic subset \( G \) of \( \mathcal{M} \) define the restricted operator norm

\[
\|A\|_G := \sup_{g \in G \subseteq \mathcal{M} : \|g\| \leq 1} \|Ag\|.
\]

Denote by \( G_0 \) the eigenspace associated to \( b_0^{-1} \) and \( A \).

**Assumption E:**

1. The estimator \( \hat{g} \) satisfies \( \|\hat{g}\|_n = 1 \) and \( \langle \hat{g}, 1 \rangle_n > 0 \). Moreover, \( \|\hat{g}\| < \infty \) a.s.

2. \( \|\hat{A} - A\|_{G_0} \to 0 \) a.s.

Condition E.1 is just a convenient normalization for our setting. Assumption E.2 is a mild consistency condition. Note that by our identification results \( G_0 \) consists of the linear span of \( g_0 \). More generally, under Assumption C, \( G_0 \) is finite dimensional, which makes E.2 easy to check; see the Appendix for primitive conditions for kernel estimators. Our next result shows the strong \( L^2 \)-consistency of our estimators.

**Theorem 3.** Let Assumptions S, C, I and E hold. Then, \( \hat{b} \to b_0 \) a.s. and \( \|\hat{g} - g_0\| \to 0 \) a.s.

We remark that Theorem 3 also holds in the partially identified case where Assumption I is dropped and the \( L^2 \)-distance between \( \hat{g} \) and \( g_0 \) is replaced by the gaps between the eigenspaces of \( \hat{A} \) and \( A \) associated to the eigenvalues \( \hat{b}^{-1} = \rho(\hat{A}) \) and \( b_0^{-1} = \rho(A) \), respectively; see Osborn (1975).

To obtain asymptotic distribution theory for our estimators, we impose the following additional assumptions and notation. Let \( A^* \) denote the adjoint operator of \( A \), that is, the linear compact operator such that \( \langle Ag_1, g_2 \rangle = \langle g_1, A^* g_2 \rangle \) for all \( g_1, g_2 \in \mathcal{M} \). Note that \( b_0^{-1} \) is also an eigenvalue for \( A^* \); eigenvalues of \( A^* \) are complex conjugates of those of \( A \). Similarly as we did for \( g_0 \), it can be shown that under our assumptions there exists a unique (up to scale) strictly positive eigenfunction of \( A^* \) associated to \( b_0^{-1} \) (see Theorem 7.C in Zeidler (1986, vol. 1, p. 290)).
Definition 3. Let \( s \) be the unique strictly positive eigenfunction of \( A^* \) with eigenvalue \( b_0^{-1} \) and satisfying the normalization \( \langle g_0, s \rangle = 1 \).

The function \( s \) plays an important role in the asymptotics for \( \hat{b} \) and \( \hat{g} \), as does the error term
\[
\varepsilon_i = g_0(C_i', V_i') R_i' - b_0^{-1} g_0(C_i, V_i), \quad i = 1, \ldots, n.
\]
(8)

Henceforth, to simplify notation, define \( \varphi_i = \varphi(C_i, V_i) \) for any \( \varphi \in L^2 \). For asymptotic normality of our estimators we require the following standard assumption. Primitive conditions for our kernel estimator are provided in the Appendix.

Assumption N.

1. \(||\hat{A} - A||_G = o_P(n^{-1/4}) \) and \( P(\hat{g} \in G) \to 1 \) as \( n \to \infty \).
2. \( \sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, s \right\rangle = \frac{1}{\sqrt{n}} \sum_{i=1}^n s_i \varepsilon_i + o_P(1). \)
3. \( \Sigma_s \equiv E[s_i^2 \varepsilon_i^2] < \infty. \)

Theorem 4. Let Assumptions S, C, I, E and N hold. Then, as \( n \to \infty \),
\[
\sqrt{n} \left( \hat{b} - b_0 \right) \overset{d}{\to} N(0, b_0^4 \Sigma_s).
\]

The proof of Theorem 4 can be found in the Appendix. For non-iid data, we obtain the same result by replacing Lemma B1 in the Appendix by Theorem 3(b) in Andrews (1995), and then applying a CLT for weakly dependent data (see e.g. Theorem 5.20 in White (2001)). For non-overlapping time periods, \( \{s_i \varepsilon_i\}_{i=1}^n \) is a martingale difference sequence, so the same asymptotic variance is obtained.

We can estimate the asymptotic variance of \( \hat{b} \) by using the sample variance of the sequence \( \{\hat{s}_i \varepsilon_i\}_{i=1}^n \) where \( \hat{\varepsilon}_i = \hat{g}(C_i', V_i') R_i' - \hat{b}_0^{-1} \hat{g}(C_i, V_i) \), and \( \hat{s} \) is computed as our estimator \( \hat{g} \), with the normalization \( \langle \hat{g}, \hat{s} \rangle_n = 1 \). An alternative to plug-in asymptotic methods is to use bootstrap. Since the eigenvalue \( b_0^{-1} \) is simple, and isolated from other eigenvalues, standard bootstrap sampling with replacement should consistently estimate the asymptotic distribution of \( \sqrt{n}(\hat{b} - b_0) \) under iid data.


Our next result establishes an asymptotic expansion for \( \hat{g} - g_0 \). This expansion can be used to obtain rates for \( \hat{g} - g_0 \) and to establish asymptotic normality of (semiparametric) functionals of \( \hat{g} \). Define the process \( \Delta_n(c, v) \equiv n^{-1} \sum_{i=1}^n \varepsilon_i \phi_i(c, v), \) where recall that \( \phi_i(c, v) = K_{hi}(c, v) / \hat{f}(c, v) \).

Note that a standard result in kernel estimation is that for all \( (c, v) \) in the interior of \( S \), under suitable conditions,
\[
\sqrt{nh_n^\alpha} \Delta_n(c, v) \overset{d}{\to} N(0, \Sigma_{\Delta}(c, v)),
\]
with \( \Sigma_\Delta (c,v) = f^{-1}(c,v)\sigma^2 (c,v) \kappa_2, \kappa_2 = \int K^2(u)du \) and \( \sigma^2 (c,v) = E [\varepsilon_i^2 | C_i = c, V_i = v] \).

Recall \( L^\dagger \) denotes the Moore-Penrose pseudoinverse of \( L = b_0A - I \), which under our conditions is linear and continuous (cf. Section 3.1).

**Theorem 5.** Let Assumptions S, C, I, E and N hold. Then, we have

\[
(bgL_n - g_0)(c,v) = b_0 L^\dagger \Delta_n (c,v) + R_n (c,v),
\]

where as \( n \to \infty, \|R_n\| = o_P(n^{-\rho/2}) \) for any \( 0 < \rho < 1 \).

This result implies that the rates of convergence of \( \hat{g} - g_0 \) in \( L^2 \) are the same as those of the NW kernel estimator of \( E [\varepsilon_i | C_i = c, V_i = v] \). We could use the expansion of Theorem 5 to test parametric hypotheses about \( g \), i.e., \( H_0 : g_0(c,v) = g_{\eta_0}(c,v) \), against nonparametric alternatives, where the function \( g_{\eta_0}(c,v) \) is known up to a finite-dimensional unknown parameter \( \eta_0 \) (e.g. power utility). A test can be based on the discrepancy

\[
D_n = \left\| \sqrt{nh^\epsilon L} (\hat{g} - \hat{g}) \right\|^2,
\]

where \( \hat{L} = \hat{b}A - I \) and \( \tilde{g} = g_{\hat{\eta}}(c,v) \) is a parametric fit, with \( \hat{\eta} \) denoting a consistent estimator for \( \eta_0 \) under the null (e.g. a GMM estimator). Noting that \( \hat{L}\hat{g} = 0 \), \( D_n \) further simplifies to \( D_n = \|\sqrt{nh^\epsilon \hat{L}\hat{g}}\|^2 \). Similar test statistics have been suggested by Härdle and Mammen (1993) in a different context. More generally, we could test nonparametric hypotheses such as the significance of certain variables, for example \( H_0 : g_0(c,v) = g_0(c,v') \) for all \( v,v' \), against nonparametric alternatives. The same \( D_n \) can be used, where now \( \tilde{g} \) denotes a restricted estimator of \( g_0 \) under the null (e.g. our marginal utility estimator depending only on \( c \)). In each case, the expansion in Theorem 5 is instrumental in analyzing the asymptotic limiting distribution of \( D_n \), which can be readily obtained combining Theorem 5 here with the results of Härdle and Mammen (1993).

6 Summary Measures

We now consider some summary measures of the model, specifically, functionals of \( \hat{g} \). These are either behavioral parameters of interest such as the average value of relative risk aversion (ARRA), or parameters having values that are relevant for testing. We first apply the results of the previous section to establish asymptotic normality of the estimated ARRA. We then list some other functionals of interest that can, in the same way, be shown to be asymptotically normal.

Define the ARRA functional by

\[
\gamma (g) \equiv E \left[-C \partial g(C,V)/\partial c \right].
\]

(9)
The natural estimator of $\gamma (g_0)$ is the sample analog based on our estimator $\hat{g}$, i.e.

$$
\gamma_n (\hat{g}) = \frac{1}{n} \sum_{i=1}^{n} \frac{-C_i \partial \hat{g}(C_i, V_i) / \partial c}{\hat{g}(C_i, V_i)}.
$$

Under the assumptions for Theorem 6 below, $\hat{g}$ is differentiable and bounded away from zero with probability tending to one, so $\gamma_n (\hat{g})$ is well-defined for large $n$. Define the class of functions

$$
D = \left\{ (c, v) \rightarrow -c \frac{\partial \log (g(c, v))}{\partial c} : g \in \mathcal{G} \right\},
$$

and the functions

$$
d(c, v) \equiv \frac{\partial (c \times f(c, v))}{\partial c} \frac{1}{f(c, v)} \quad \text{and} \quad \chi(c, v) \equiv \frac{d(c, v)}{g_0(c, v)}. \tag{11}
$$

Also, we need to introduce some notation to be used in the asymptotic normality of $\gamma_n (\hat{g})$. Assuming $\chi \in L^2$, define

$$
\chi_s = \chi - \langle g_0, \chi \rangle \langle g_0, s \rangle^{-1} s. \tag{12}
$$

The function $\chi_s$ has a geometrical interpretation as the value of $\chi$ projected parallel to $s$ on a subspace of functions orthogonal to $g_0$. Let $L^*$ denote the adjoint operator of $L$, and let $\chi_s^*$ denote the minimum norm solution of $\chi_s = L^* r$ in $r$, i.e. $\chi_s^* = \arg \min \{ \| r \| : \chi_s = L^* r \}$, which is well defined because $\chi_s \in \mathcal{N}^\perp (L) = \mathcal{R}(L^*)$; see Luenberger (1997, Theorem 3, p. 157) for the latter equality. Here $\mathcal{N}^\perp (L)$ denotes the orthogonal complement of the null space of $L$, see Luenberger (1997, p. 52) for a definition.

We also introduce a class of smooth function $C^q(T)$ for a generic closed and convex set $T$. For any vector $a$ of $\ell$ integers define the differential operator $\partial^a_x \equiv \partial^{|a_1|}/\partial x_1^{a_1} \ldots \partial x_\ell^{a_\ell}$, where $|a_1| \equiv \sum_{i=1}^{\ell} a_i$. For any smooth function $h : T \subset \mathbb{R}^\ell \rightarrow \mathbb{R}$ and some $\eta > 0$, let $\eta$ be the largest integer smaller or equal than $\eta$, and

$$
\| h \|_{\infty, \eta} \equiv \max_{|a_1| \leq \eta} \sup_{x \in T} |\partial^a_x h(x)| + \max_{|a_1| = \eta} \sup_{x \neq x'} \frac{|\partial^a_x h(x) - \partial^a_{x'} h(x')|}{|x - x'|^{\eta - 2}}.
$$

Further, let $C^q_m (T)$ be the set of all continuous functions $h : T \subset \mathbb{R}^\ell \rightarrow \mathbb{R}$ with $\| h \|_{\infty, \eta} \leq M$ (for an integer $\eta$, the $\eta$-th derivative is assumed to be continuous). Since the constant $M$ is irrelevant for our results, we drop the dependence on $M$ and denote $C^q (T)$.

The $ARRA$ estimator behaves asymptotically as a sample average, with an influence function given by

$$
\xi_i = (\zeta_i - E [\zeta_i]) - b_0 \chi_s^*(C_i, V_i) \varepsilon_i, \tag{13}
$$

where $\zeta_i = -C_i (\partial g_0 (C_i, V_i) / \partial c) / g_0 (C_i, V_i)$. The second term in $\xi_i$ accounts for the estimation effect due to estimating $g_0$.

**ASSUMPTION CE.**
1. The class $\mathcal{D}$ defined in (10) is $P$-Donsker.\footnote{Let $P_n$ be the empirical measure with respect to $P$. Using a standard empirical process notation, define $\mathcal{G}_n = \sqrt{n} (P_n - P)$. Then $\mathcal{D}$ is $P$-Donsker if $\mathcal{G}_n$ converges weakly to $\mathcal{G}$ in the space of uniformly bounded functions on $\mathcal{D}$, $l^\infty (\mathcal{D})$, where $\mathcal{G}$ is a mean-zero $P$-Brownian bridge process with uniformly continuous sample paths with respect to the semi-metric $\rho (d,d')$ defined by $\rho (d,d') = \sqrt{\text{Var} (d(C,V) - d'(C,V))}$. For further details we refer the reader to van der Vaart and Wellner (1996).}

2. The measure $\mu$ is the probability measure of $(C,V)$ with a support $S$ that can be written as $S = [l_c,u_c] \times S_V$, for some $l_c,u_c$ with $l_c < u_c$. Furthermore, $\lim_{c \to -l_c} cf(c,v) = \lim_{c \to u_c} cf(c,v)$ for all $v \in S_V$ and $P (\min \{g_0, \hat{g}\} > \varepsilon ) \to 1$ for some $\varepsilon > 0$.

3. The function $d$ defined in (11) satisfies $d \in \mathcal{L}^2$, $E[|\xi_i|^2] < \infty$ and $\chi_0 \in \mathcal{C}^r(S)$.

Assumption CE.1 is standard in the semiparametric literature, see, e.g. Chen, Linton and Van Keilegom (2003). The following Lemma provides sufficient conditions for an example of $\mathcal{D}$ satisfying the $P$-Donsker property. Its proof is a standard exercise in empirical processes theory, and hence it is omitted.

**Lemma 2.** Suppose that $\mathcal{G}$ is a subset of $\mathcal{C}^\eta (S)$ of functions bounded away from zero, where $\eta > (2 + \ell)/2$, and that $E[C_1^2] < \infty$. Then, $\mathcal{D}$ is $P$-Donsker.

Assumption CE.2 is similar to other assumptions required in estimation of average derivatives, see Powell, Stock and Stoker (1989). This assumption guarantees that $\gamma_n (\hat{g})$ is well defined. Assumption CE.3 implies that the asymptotic variance of $\gamma_n (\hat{g})$ is finite.

**Theorem 6.** Let Assumptions S, C, I, E, N and CE hold. Then,

$$
\sqrt{n} (\gamma_n (\hat{g}) - \gamma (g_0)) \xrightarrow{d} N (0, E [\xi_i^2]) ,
$$

where $\xi_i$ is defined in (13).

Estimating the asymptotic variance of $\gamma_n (\hat{g})$ by plug-in methods would be possible but complicated. An alternative is to use the bootstrap, which can be justified along the lines of Chen, Linton and Van Keilegom (2003). Extensions of Theorem 6 to non-iid settings can be readily obtained, as discussed above in Theorem 4.

Now consider some other functionals of interest. The asymptotic normality of each can be established using the same methods as Theorem 6. As with $\gamma_n (\hat{g})$, we can use the bootstrap to estimate their limiting distributions. In our empirical work we consider a model allowing for habits, where $C_t^i = C_{t+1,i}$ and $V_i = C_{t,i}$ for two time periods $t$ and $t + 1$ (these time periods may vary across
individuals). For the remainder of this section we drop the $i$ subscript for clarity. Closely related to the ARRA are local averages defined by

$$
\rho (q, j) = E \left[ -C_{t+1} \frac{\partial g_0(C_{t+1}, C_t)}{g_0(C_{t+1}, C_t)} \bigg| C_{t+1} \in Q_q, C_t \in S_j \right],
$$

(14)

where $Q_q$ denotes the interval between the $q - 1$ and $q$ quartile of $C_{t+1}$, and $S_j$ denotes the interval between the $j - 1$ and $j$ quartile of $C_t$ for $q, j = 1, 2, 3, 4$. We refer to each of these local averages of the RRA between different quartiles as a QRRA (quartile relative risk aversion).

We can use our results to construct tests of heterogeneity in risk aversion measures as follows. The sample analogs of the QRRA parameters $\rho (q, j)$ can be shown to be asymptotically normal under the same conditions above used for the ARRA. That is, with the simplified notation $\rho (q) \equiv \rho (q, q)$ for the parameter and $\rho_n (q) \equiv \rho_n (q, q)$ for the plug-in estimator, it can be shown

$$
\sqrt{n} (\rho_n (q) - \rho (q)) \xrightarrow{d} N \left( 0, \sigma^2 (q) \right),
$$

for a suitable asymptotic variance $\sigma^2 (q), q = 1, 2, 3$ and 4. Moreover, by definition, $\sqrt{n} (\rho_n (q) - \rho (q))$ and $\sqrt{n} (\rho_n (j) - \rho (j))$ are asymptotically independent for $q \neq j$. This suggests a simple strategy for testing heterogeneity in risk aversion by means of simple pairwise t-tests for the hypotheses, for $q \neq j,$

$$
H_{0qj} : \rho (q) = \rho (j) \quad vs \quad H_{1qj} : \rho (q) \neq \rho (j).
$$

The t-statistics are constructed as

$$
t_{qj} = \frac{\sqrt{n} (\rho_n (q) - \rho_n (j))}{\sqrt{\hat{\sigma}_n^2 (q) + \hat{\sigma}_n^2 (j)}},
$$

for suitable consistent estimates $\sigma^2_n (q)$ of the asymptotic variances $\sigma^2 (q),$ for $q = 1, 2, 3$ and 4. We then reject $H_{0qj}$ when $t_{qj}$ is large in absolute value, using that $t_{qj}$ converges to a standard normal under $H_{0qj}$.

We also construct some tests for the absence of habits, i.e.

$$
\frac{\partial g_0(C_{t+1}, C_t)}{\partial C_t} = 0.
$$

Our tests are based on the functional

$$
\delta (g) = E \left[ \frac{\partial g(C_{t+1}, C_t)}{\partial C_t} \tau (C_{t+1}, C_t) \right],
$$

for various positive functions $\tau (\cdot)$. When there is no habit effect $\delta (g_0) = 0$ for any choice of $\tau$. As with $\gamma (g_0),$ for each choice of function $\tau$ we estimate $\delta (g_0)$ by plugging in $\hat{g}$ for $g_0$ and replacing the expectation with a sample average. The asymptotic normality of this estimator and its bootstrap approximation is then used for inference, analogous to our analysis of $\gamma (g_0)$. 

19
7 Monte Carlo Experiment

In this section we illustrate the finite-sample performance of our estimator described in the previous sections based on a CRRA utility function so that \( g_0(c, v) = c^{-\eta_0} \), where \( \eta_0 \) in this case equals the ARRRA. The model is then given by the Euler equation

\[
0 = b_0 E \left[ C_{t+1}^{1-\eta_0} \frac{R_{t+1}}{C_t} \right] = C_t^{-\eta_0}.
\]

We set \( b_0 = 0.95 \) and \( \eta_0 = 0.5 \). We draw a random sample of \((C_t, C_{t+1})\) from the distribution

\[
\begin{pmatrix}
\log C_t \\
\log C_{t+1}
\end{pmatrix} \sim N \left( 0, \begin{pmatrix} 0.25 & 0.1 \\
0.1 & 0.25 \end{pmatrix} \right),
\]

and construct \( R_{t+1} = b_0^{-1}(1 + \epsilon_t) (C_{t+1}/C_t)^{\eta_0} \), where \( \epsilon_t \) is distributed uniformly on \([-0.5, 0.5]\) and drawn independently of \((C_t, C_{t+1})\). This design was chosen to generate data that satisfies the Euler equation model, has realistic parameter values and consumption distribution, and avoids the approximation and other numerical errors that would result from solving each individual’s dynamic optimization problem numerically.

To save space we only report simulation results for two experiments, each with sample sizes \( n = 500 \) and \( n = 2000 \). The number of bootstrap replications used in each simulation is 200, and we repeat each simulation 1000 times. We compute our proposed nonparametric estimators and compare them to the method of moments estimator defined using the correctly specified CRRA utility function with a constant and \( C_t \) as instruments. So while our estimator attempts to recover the constant \( b_0 \) and the entire function \( g_0 \), this alternative just estimates the two constants \( b_0 \) and \( \eta_0 \), using two moments of the data. In our tables estimates from this correctly specified parametric functional form are labeled CRRA.

We consider two nonparametric estimators. The first one, which we label \( NP - 1 \), correctly conditions on just \( C_t \) (since our choice of \( g_0(c, v) \) does not depend on \( v \)), and so only entails estimation of a one-dimensional marginal utility function. The second nonparametric estimator, denoted \( NP - 2 \), uses both \( C_t \) and \( V_t \) as conditioning variables, where \( V_t = C_{t-1} \) is in this case an irrelevant habit variable. We simulate \( C_{t-1} \) by drawing from a \( N(1, 1) \) distribution that is independent of \((C_t, C_{t+1})\).

We compute our estimates using the procedure described in Section 4 that incorporates the transformation suggested in equation (4). While not necessary in theory, we find that estimates of \( g_0 \) fit better in the tails using this transformation than not, though the differences in overall integrated mean square errors and other measures of fit are small. In order to apply the transformation, note that equation (4) can be re-written as

\[
bE[g^*(C_{t+1}, V_{t+1})R_{t+1}^* \mid C_t, V_t] = g^*(C_t, V_t),
\]
where \( g^*(C_{t+1}, V_{t+1}) \equiv C_{t+1}g(C_{t+1}, V_{t+1}), g^*(C_t, V_t) \equiv C_t g(C_t, V_t) \) and \( R_{t+1}^* \equiv (C_t/C_{t+1}) R_{t+1} \). With these definitions the procedure remains as described in Section 4 after redefining the return variable, from \( R_{t+1} \) to \( R_{t+1}^* \). The procedure then yields an estimate of \( g^* \), from which the marginal utility function \( g \) is then recovered using the relation \( g(c, v) = g^*(c, v)/c \). Throughout we set the bandwidth to be \( 1.06\sigma n^{-1/3.5} \), where \( \sigma \) is the sample standard deviation of \( C_t \). This is essentially Silverman’s rule applied to the rate \( n^{-1/3.5} \). All of our estimators for \( g_0 \) are normalized to have a unit norm with respect to the empirical \( L^2 \)-norm.

For each finite-dimensional parameter and summary measure we consider, we report the mean, standard deviation, 2.5th percentile, 97.5th percentile, 95% coverage probability based on normal distribution, their bootstrap counterparts and the root mean square error.\(^5\) Table 1 reports estimates of the discount factor from our three estimators, \( CRRA \), \( NP - 1 \), and \( NP - 2 \). Table 2 reports estimates of the \( ARRA \), which for the \( CRRA \) model is just the estimated constant \( \eta_0 \), while for the nonparametric estimators the \( ARRA \) is \( \gamma (g_0) \) defined by equation (9). Table 1 shows that all of the estimators succeed in estimating the discount factor \( b \) very accurately. This is in contrast to many macro models, which often calibrate the discount factor due to the difficulty in estimating it accurately. Table 2 shows somewhat more difficulty in estimating the \( ARRA \), but the relative accuracy of our nonparametric estimates to the parametric alternative is similar. In both tables the root mean squared errors of our nonparametric estimates are seen to shrink with sample size and increase with dimensionality at rates that are generally consistent with asymptotic theory.

Figures 1 and 2 show plots of the one-dimensional nonparametric (i.e., \( NP - 1 \)) estimated marginal utility function \( g_0 \) as a function of \( C_t \). Figure 1 is \( n = 500 \) while Figure 2 is \( n = 2000 \). For each figure, the solid line denotes the mean, the dotted line denotes the 95% confidence interval, and the dashed line is the true. One can see from these figures that \( NP - 1 \) quite accurately tracks the true function. The precision of these fits can also be summarized by their integrated mean square error (weighted with respect to the true density), which is 0.0014 for \( n = 500 \) and 0.0005 for \( n = 2000 \).

Not surprisingly, estimates of the two-dimensional \( NP - 2 \) are noisier, since by design the second conditioning variable \( V_t \) is irrelevant. The results for \( NP - 2 \) can be summarized by their implied quartile averages \( QRRA \). Table 3 reports estimates of each \( QRRA \), \( \rho (q, j) \) for all quartiles \( q \) and \( j \) having \( |q - j| \leq 1 \).\(^6\) Table 3 shows that estimates of \( QRRA \) have generally about an order of magnitude larger root mean squared error than \( ARRA \), which is not surprising since each \( \rho (q, j) \) is obtained by averaging over 1/16 as much data (one quartile of current consumption and one quartile

\(^5\)The normal coverage probability is constructed ex-post using the true (simulated) standard deviation.

\(^6\)We only report pairs of quartiles \( q \) and \( s \) where \( |q - s| \leq 1 \), because a value that violates this inequality, like \( \rho (4, 1) \), corresponds to individuals who’s consumption jumps from the fourth to the first quartile, and in real data the number of such individuals who make this jump would be too small to reliably estimate their \( QRRA \).
of lagged consumption observations) as \textit{ARRA}.

One unexpected finding is that estimates of $\rho(q,j)$ display substantially larger biases and root mean squared errors for larger values of $q$ and $j$ than for smaller values, suggesting that our $NP - 2$ estimates of the marginal utility function tend to be less accurate at higher consumption levels. This can also be seen for $NP - 1$ in Figure 1, where the standard error bands widen at higher consumption levels.

In Table 4 we report estimates of $\delta(g_0)$ that can be used to test for the presence of habits in $g_0$. In our experiments estimates of $\delta(g_0)$ do not differ significantly from zero as expected, since our specification of $g_0$ does not have any habit effect. Generally, all of our parameter estimates and test statistics appear to have distributions across simulations that are reasonably well approximated by the bootstrap, e.g., biases are relatively small, bootstrap standard errors are generally close to the standard deviations across simulations, and bootstrap confidence intervals are generally close to the true. Both coverage probabilities based on the normal approximation and the bootstrap generally are relatively close to the nominal.
### Table 1: Summary statistics of Monte Carlo estimates of the discount factor $b_0$. The true is $b_0 = 0.95$. $CRR A$, $NP - 1$ and $NP - 2$ refer respectively to the parametric, one-dimensional nonparametric, and two-dimensional nonparametric estimators.

<table>
<thead>
<tr>
<th></th>
<th>$b_0$</th>
<th>Bias</th>
<th>Std</th>
<th>Lpc</th>
<th>Upc</th>
<th>Cov</th>
<th>B-Std</th>
<th>B-Lpc</th>
<th>B-Upc</th>
<th>B-Cov</th>
<th>Rmse</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n = 500$</td>
<td>$CRR A$</td>
<td>0.000</td>
<td>0.012</td>
<td>0.926</td>
<td>0.975</td>
<td>0.946</td>
<td>0.012</td>
<td>0.926</td>
<td>0.974</td>
<td>0.940</td>
<td>0.012</td>
</tr>
<tr>
<td></td>
<td>$NP - 1$</td>
<td>0.006</td>
<td>0.027</td>
<td>0.917</td>
<td>0.971</td>
<td>0.984</td>
<td>0.018</td>
<td>0.915</td>
<td>0.980</td>
<td>0.929</td>
<td>0.028</td>
</tr>
<tr>
<td></td>
<td>$NP - 2$</td>
<td>0.009</td>
<td>0.041</td>
<td>0.808</td>
<td>0.983</td>
<td>0.963</td>
<td>0.031</td>
<td>0.895</td>
<td>1.012</td>
<td>0.932</td>
<td>0.042</td>
</tr>
<tr>
<td>$n = 2000$</td>
<td>$CRR A$</td>
<td>0.000</td>
<td>0.006</td>
<td>0.938</td>
<td>0.961</td>
<td>0.960</td>
<td>0.006</td>
<td>0.938</td>
<td>0.962</td>
<td>0.950</td>
<td>0.006</td>
</tr>
<tr>
<td></td>
<td>$NP - 1$</td>
<td>0.004</td>
<td>0.020</td>
<td>0.936</td>
<td>0.960</td>
<td>0.992</td>
<td>0.009</td>
<td>0.932</td>
<td>0.965</td>
<td>0.924</td>
<td>0.020</td>
</tr>
<tr>
<td></td>
<td>$NP - 2$</td>
<td>0.005</td>
<td>0.028</td>
<td>0.862</td>
<td>0.965</td>
<td>0.974</td>
<td>0.021</td>
<td>0.922</td>
<td>0.994</td>
<td>0.946</td>
<td>0.028</td>
</tr>
</tbody>
</table>

### Table 2: Summary statistics of Monte Carlo estimates of the $ARR A$, which is $\eta_0$ for the parametric and $\gamma (g_0)$ for the nonparametric estimators. The true is $ARR A = 0.5$. $CRR A$, $NP - 1$ and $NP - 2$ refer respectively to the parametric, one-dimensional nonparametric, and two-dimensional nonparametric estimators.

<table>
<thead>
<tr>
<th></th>
<th>$ARR A$</th>
<th>Bias</th>
<th>Std</th>
<th>Lpc</th>
<th>Upc</th>
<th>Cov</th>
<th>B-Std</th>
<th>B-Lpc</th>
<th>B-Upc</th>
<th>B-Cov</th>
<th>Rmse</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n = 500$</td>
<td>$CRR A$</td>
<td>0.000</td>
<td>0.046</td>
<td>0.420</td>
<td>0.590</td>
<td>0.956</td>
<td>0.046</td>
<td>0.411</td>
<td>0.592</td>
<td>0.944</td>
<td>0.046</td>
</tr>
<tr>
<td></td>
<td>$NP - 1$</td>
<td>-0.058</td>
<td>0.107</td>
<td>0.431</td>
<td>0.714</td>
<td>0.961</td>
<td>0.101</td>
<td>0.359</td>
<td>0.751</td>
<td>0.906</td>
<td>0.122</td>
</tr>
<tr>
<td></td>
<td>$NP - 2$</td>
<td>-0.096</td>
<td>0.194</td>
<td>0.277</td>
<td>0.888</td>
<td>0.952</td>
<td>0.194</td>
<td>0.209</td>
<td>0.986</td>
<td>0.930</td>
<td>0.217</td>
</tr>
<tr>
<td>$n = 2000$</td>
<td>$CRR A$</td>
<td>0.001</td>
<td>0.023</td>
<td>0.456</td>
<td>0.545</td>
<td>0.950</td>
<td>0.023</td>
<td>0.454</td>
<td>0.544</td>
<td>0.952</td>
<td>0.023</td>
</tr>
<tr>
<td></td>
<td>$NP - 1$</td>
<td>-0.032</td>
<td>0.077</td>
<td>0.470</td>
<td>0.610</td>
<td>0.988</td>
<td>0.052</td>
<td>0.430</td>
<td>0.628</td>
<td>0.914</td>
<td>0.083</td>
</tr>
<tr>
<td></td>
<td>$NP - 2$</td>
<td>-0.067</td>
<td>0.092</td>
<td>0.412</td>
<td>0.716</td>
<td>0.934</td>
<td>0.109</td>
<td>0.355</td>
<td>0.782</td>
<td>0.906</td>
<td>0.114</td>
</tr>
</tbody>
</table>
Table 3: Summary statistics of Monte Carlo estimates of $QRRA$, which is $\rho(q,j)$ from $NP-2$.

The true is $\rho(q,j) = 0.5$ for all $q$ and $j$.  

<table>
<thead>
<tr>
<th>$QRRA$</th>
<th>Bias</th>
<th>Std</th>
<th>Lpc</th>
<th>Upc</th>
<th>Cov</th>
<th>B-Std</th>
<th>B-Lpc</th>
<th>B-Upc</th>
<th>B-Cov</th>
<th>Rmse</th>
</tr>
</thead>
<tbody>
<tr>
<td>$n = 500$</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>$\rho(1, 1)$</td>
<td>-0.158</td>
<td>0.205</td>
<td>0.273</td>
<td>1.068</td>
<td>0.910</td>
<td>0.242</td>
<td>0.115</td>
<td>1.068</td>
<td>0.878</td>
<td>0.259</td>
</tr>
<tr>
<td>$\rho(1, 2)$</td>
<td>-0.068</td>
<td>0.366</td>
<td>-0.049</td>
<td>1.167</td>
<td>0.969</td>
<td>0.358</td>
<td>-0.137</td>
<td>1.287</td>
<td>0.969</td>
<td>0.372</td>
</tr>
<tr>
<td>$\rho(2, 1)$</td>
<td>-0.149</td>
<td>0.222</td>
<td>0.242</td>
<td>1.060</td>
<td>0.932</td>
<td>0.246</td>
<td>0.145</td>
<td>1.118</td>
<td>0.904</td>
<td>0.267</td>
</tr>
<tr>
<td>$\rho(2, 2)$</td>
<td>-0.055</td>
<td>0.327</td>
<td>0.000</td>
<td>1.151</td>
<td>0.961</td>
<td>0.355</td>
<td>-0.137</td>
<td>1.274</td>
<td>0.965</td>
<td>0.331</td>
</tr>
<tr>
<td>$\rho(2, 3)$</td>
<td>-0.010</td>
<td>0.450</td>
<td>-0.240</td>
<td>1.187</td>
<td>0.973</td>
<td>0.480</td>
<td>-0.433</td>
<td>1.477</td>
<td>0.973</td>
<td>0.450</td>
</tr>
<tr>
<td>$\rho(3, 2)$</td>
<td>-0.053</td>
<td>0.326</td>
<td>-0.014</td>
<td>1.081</td>
<td>0.969</td>
<td>0.351</td>
<td>-0.121</td>
<td>1.275</td>
<td>0.966</td>
<td>0.330</td>
</tr>
<tr>
<td>$\rho(3, 3)$</td>
<td>0.009</td>
<td>0.457</td>
<td>-0.279</td>
<td>1.180</td>
<td>0.972</td>
<td>0.460</td>
<td>-0.408</td>
<td>1.428</td>
<td>0.966</td>
<td>0.457</td>
</tr>
<tr>
<td>$\rho(3, 4)$</td>
<td>-0.102</td>
<td>0.785</td>
<td>-0.850</td>
<td>1.972</td>
<td>0.963</td>
<td>0.933</td>
<td>-1.320</td>
<td>2.452</td>
<td>0.972</td>
<td>0.792</td>
</tr>
<tr>
<td>$\rho(4, 3)$</td>
<td>-0.029</td>
<td>0.400</td>
<td>-0.137</td>
<td>1.181</td>
<td>0.969</td>
<td>0.470</td>
<td>-0.345</td>
<td>1.515</td>
<td>0.978</td>
<td>0.401</td>
</tr>
<tr>
<td>$\rho(4, 4)$</td>
<td>-0.281</td>
<td>0.980</td>
<td>-0.957</td>
<td>2.378</td>
<td>0.954</td>
<td>1.079</td>
<td>-1.486</td>
<td>2.876</td>
<td>0.955</td>
<td>1.019</td>
</tr>
<tr>
<td>$n = 2000$</td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>$\rho(1, 1)$</td>
<td>-0.104</td>
<td>0.179</td>
<td>0.350</td>
<td>0.825</td>
<td>0.978</td>
<td>0.158</td>
<td>0.280</td>
<td>0.889</td>
<td>0.888</td>
<td>0.206</td>
</tr>
<tr>
<td>$\rho(1, 2)$</td>
<td>-0.023</td>
<td>0.272</td>
<td>0.125</td>
<td>0.903</td>
<td>0.984</td>
<td>0.249</td>
<td>0.048</td>
<td>1.027</td>
<td>0.954</td>
<td>0.273</td>
</tr>
<tr>
<td>$\rho(2, 1)$</td>
<td>-0.087</td>
<td>0.146</td>
<td>0.330</td>
<td>0.859</td>
<td>0.938</td>
<td>0.171</td>
<td>0.245</td>
<td>0.910</td>
<td>0.912</td>
<td>0.170</td>
</tr>
<tr>
<td>$\rho(2, 2)$</td>
<td>-0.018</td>
<td>0.214</td>
<td>0.151</td>
<td>0.882</td>
<td>0.964</td>
<td>0.251</td>
<td>0.031</td>
<td>1.030</td>
<td>0.968</td>
<td>0.214</td>
</tr>
<tr>
<td>$\rho(2, 3)$</td>
<td>-0.007</td>
<td>0.319</td>
<td>0.004</td>
<td>1.019</td>
<td>0.988</td>
<td>0.314</td>
<td>-0.104</td>
<td>1.133</td>
<td>0.956</td>
<td>0.319</td>
</tr>
<tr>
<td>$\rho(3, 2)$</td>
<td>-0.009</td>
<td>0.274</td>
<td>0.078</td>
<td>0.871</td>
<td>0.980</td>
<td>0.254</td>
<td>0.024</td>
<td>1.013</td>
<td>0.954</td>
<td>0.274</td>
</tr>
<tr>
<td>$\rho(3, 3)$</td>
<td>-0.016</td>
<td>0.376</td>
<td>0.095</td>
<td>0.956</td>
<td>0.986</td>
<td>0.310</td>
<td>-0.067</td>
<td>1.153</td>
<td>0.962</td>
<td>0.377</td>
</tr>
<tr>
<td>$\rho(3, 4)$</td>
<td>-0.078</td>
<td>0.388</td>
<td>-0.136</td>
<td>1.322</td>
<td>0.952</td>
<td>0.573</td>
<td>-0.583</td>
<td>1.722</td>
<td>0.970</td>
<td>0.396</td>
</tr>
<tr>
<td>$\rho(4, 3)$</td>
<td>-0.002</td>
<td>0.385</td>
<td>0.129</td>
<td>0.913</td>
<td>0.980</td>
<td>0.302</td>
<td>-0.054</td>
<td>1.123</td>
<td>0.964</td>
<td>0.385</td>
</tr>
<tr>
<td>$\rho(4, 4)$</td>
<td>-0.244</td>
<td>0.476</td>
<td>0.053</td>
<td>1.641</td>
<td>0.940</td>
<td>0.624</td>
<td>-0.571</td>
<td>1.948</td>
<td>0.958</td>
<td>0.535</td>
</tr>
<tr>
<td></td>
<td>$\tau (C_{t+1}, C_t)$</td>
<td>Bias</td>
<td>Std</td>
<td>Lpc</td>
<td>Upc</td>
<td>Cov</td>
<td>B-Std</td>
<td>B-Lpc</td>
<td>B-Upc</td>
<td>B-Cov</td>
</tr>
<tr>
<td>-------</td>
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<td>-------</td>
<td>-------</td>
<td>-------</td>
</tr>
<tr>
<td>$n = 500$</td>
<td>$C_{t+1}$</td>
<td>-0.002</td>
<td>0.111</td>
<td>-0.111</td>
<td>0.132</td>
<td>0.975</td>
<td>0.118</td>
<td>-0.255</td>
<td>0.200</td>
<td>0.975</td>
</tr>
<tr>
<td></td>
<td>$C_t$</td>
<td>-0.006</td>
<td>0.097</td>
<td>-0.128</td>
<td>0.125</td>
<td>0.975</td>
<td>0.118</td>
<td>-0.245</td>
<td>0.209</td>
<td>0.980</td>
</tr>
<tr>
<td></td>
<td>$C_{t+1}^2$</td>
<td>-0.010</td>
<td>0.289</td>
<td>-0.249</td>
<td>0.252</td>
<td>0.977</td>
<td>0.262</td>
<td>-0.567</td>
<td>0.438</td>
<td>0.965</td>
</tr>
<tr>
<td></td>
<td>$C_t^2$</td>
<td>-0.030</td>
<td>0.237</td>
<td>-0.331</td>
<td>0.270</td>
<td>0.967</td>
<td>0.269</td>
<td>-0.531</td>
<td>0.502</td>
<td>0.977</td>
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<tr>
<td></td>
<td>$C_{t+1}C_t$</td>
<td>-0.015</td>
<td>0.229</td>
<td>-0.209</td>
<td>0.190</td>
<td>0.972</td>
<td>0.220</td>
<td>-0.463</td>
<td>0.370</td>
<td>0.973</td>
</tr>
<tr>
<td>$n = 2000$</td>
<td>$C_{t+1}$</td>
<td>-0.005</td>
<td>0.078</td>
<td>-0.070</td>
<td>0.072</td>
<td>0.978</td>
<td>0.077</td>
<td>-0.154</td>
<td>0.131</td>
<td>0.978</td>
</tr>
<tr>
<td></td>
<td>$C_t$</td>
<td>-0.009</td>
<td>0.080</td>
<td>-0.084</td>
<td>0.072</td>
<td>0.982</td>
<td>0.077</td>
<td>-0.154</td>
<td>0.132</td>
<td>0.978</td>
</tr>
<tr>
<td></td>
<td>$C_{t+1}^2$</td>
<td>-0.013</td>
<td>0.229</td>
<td>-0.176</td>
<td>0.149</td>
<td>0.986</td>
<td>0.188</td>
<td>-0.374</td>
<td>0.319</td>
<td>0.968</td>
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<tr>
<td></td>
<td>$C_t^2$</td>
<td>-0.036</td>
<td>0.244</td>
<td>-0.270</td>
<td>0.150</td>
<td>0.986</td>
<td>0.195</td>
<td>-0.382</td>
<td>0.344</td>
<td>0.966</td>
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<tr>
<td></td>
<td>$C_{t+1}C_t$</td>
<td>-0.016</td>
<td>0.222</td>
<td>-0.146</td>
<td>0.107</td>
<td>0.984</td>
<td>0.160</td>
<td>-0.313</td>
<td>0.268</td>
<td>0.970</td>
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Table 4: Summary statistics of Monte Carlo estimates of $\delta (g_0)$, used to test for the presence of habit effects. The true value of each $\delta (g_0)$ is zero. The $\tau (C_{t+1}, C_t)$ column lists the functions that are used to define $\delta (g_0)$. 
Figure 1: Estimates of the marginal utility function $g_0$ using simulated data with $n = 500$. $Est$, $CI$, and $True$ represent respectively the one-dimensional nonparametric estimator, its 95% confidence interval, and the true.

Figure 2: Estimates of the marginal utility function $g_0$ using simulated data with $n = 2000$. $Est$, $CI$, and $True$ represent respectively the one-dimensional nonparametric estimator, its 95% confidence interval, and the true.
8 Conclusions

We investigate nonparametric identification and estimation of marginal utilities and discount factors in consumption-based asset pricing Euler equations. The main features of our nonparametric identification results are: (i) the decomposition of the pricing kernel into its marginal utility and discount factor components, cast in the form equation (1), and (ii) the use of shape restrictions (positive marginal utilities). Together, these allow us to establish nonparametric global point identification of the model. Based on our identification arguments, we propose a new nonparametric estimator for marginal utilities and the discount factor that combines standard kernel estimation with the computation of a (finite-dimensional) matrix eigenvalue-eigenvector problem. No numerical integration or optimization is involved. The estimator is based on a sample analogue of (1) and is easy to implement, since no numerical searches are required. We establish a useful expansion for the marginal utility (suitably normalized), and limiting distribution theory for the discount factor and associated functionals of the marginal utility like the average level of relative risk aversion. Due to the well posedness of equation (1), our estimator converges at comparable rates to ordinary nonparametric regression and does not suffer from issues associated with nonparametric instrumental variables estimation.

In the older version of our paper, we apply our nonparametric methods to household-level CEX data and find evidence against the common assumption of constant relative risk aversion across consumers. Our estimates are fairly insensitive to the choice of asset used (risk-free vs risky), which supports our nonparametric model. We find empirical evidence for the presence of habits, and evidence that risk aversion varies across current and lagged consumption levels in ways that are not fully captured by standard parametric or even semiparametric specifications of habits in asset pricing models. However, there are two caveats that underly our findings. One is we used common asset returns with repeated cross-section of household consumptions since we do not have data on household-level assets. Our existing theory requires suitable variation in the returns and extending it to clustering data is beyond the scope of this paper. Another concern with the dataset we used is the presence of measurement error in consumption, see e.g. Alan, Attanasio and Browning (2009) and references therein. Escanciano (2018) has recently shown that the identification and estimation of the discount factor in this paper is robust to the presence of measurement error under rather general conditions. Measures of risk aversion, however, are more sensitive to measurement error and may require different estimation strategies to achieve robustness. These strategies are well beyond the scope of this paper and will be investigated in future research.
9 Appendix

9.1 Euler Equation Derivation

To encompass a large class of existing Euler equation and asset pricing models, consider utility functions that in addition to ordinary consumption, may include both durables and habit effects. Let \( U \) be a time homogeneous period utility function, \( b \) is the one period subjective discount factor, \( C_t \) is expenditures on consumption, \( D_t \) is a stock of durables, and \( Z_t \) is a vector of other variables that affect utility and are known at time \( t \). Let \( V_t \) denote the vector of all variables other than \( C_t \) that affect utility in time \( t \). In particular, \( V_t \) contains \( Z_t \), \( V_t \) contains \( D_t \) if durables matter, and \( V_t \) contains lagged consumption \( C_{t-1}, C_{t-2} \) and so on if habits matter.

The consumer's time separable utility function is

\[
\max_{\{C_t, D_t\}_{t=0}^{\infty}} E \left[ \sum_{t=0}^{\infty} b^t U(C_t, V_t) \right].
\]

The consumer saves by owning durables and by owning quantities of risky assets \( A_{jt}, j = 1, \ldots, J \). Letting \( C_t \) be the numeraire, let \( P_t \) be the price of durables \( D_t \) at time \( t \) and let \( R_{jt} \) be the gross return in time period \( t \) of owning one unit of asset \( j \) in period \( t - 1 \). Assume the depreciation rate of durables is \( \delta \). Then without frictions the consumer's budget constraint can be written as, for each period \( t \),

\[
C_t + (D_t - \delta D_{t-1}) P_t + \sum_{j=1}^{J} A_{jt} \leq \sum_{j=1}^{J} (A_{jt} - A_{jt-1} R_{jt})
\]

We may interpret this model either as a representative consumer model, or a model of individual agents which may vary by their initial endowments of durables and assets and by \( \{Z_t\}_{t=0}^{\infty} \). The Lagrangean is

\[
E \left[ \sum_{t=0}^{T} b^t U(C_t, V_t) - \left( C_t + (D_t - \delta D_{t-1}) P_t + \sum_{j=1}^{J} (A_{jt} - A_{jt-1} R_{jt}) \right) \lambda_t \right]
\]

with Lagrange multipliers \( \{\lambda_t\}_{t=0}^{\infty} \).

Consider the roles of durables and habits. For durables, define

\[
g_d(C_t, V_t) = \frac{\partial U(C_t, V_t)}{\partial D_t}
\]

which will be nonzero only if \( V_t \) contains \( D_t \). For habits, we must handle the possibility of both internal or external habits. Habits are defined to be internal (or internalized) if the consumer considers both the direct effects of current consumption on future utility through habit as well as through the budget constraint. In the above notation, habits are internal if the consumer takes into
account the fact that, due to habits, changing $C_t$ will directly change $V_{t+1}, V_{t+2}$ etc. Otherwise, if the consumer ignores this effect when maximizing, then habits called external.

If habits are external or if there are no habit effects at all, then define the marginal utility function $g$ by

$$g(C_t, V_t) = \frac{\partial U(C_t, V_t)}{\partial C_t}$$

If habits exist and are internal then define the function $\tilde{g}$ by

$$\tilde{g}(I_t) = \sum_{\ell=0}^L b^\ell E \left[ \frac{\partial U(C_{t+\ell}, V_{t+\ell})}{\partial C_t} \mid I_t \right].$$

where $L$ is such that $V_t$ contains $C_{t-1}, C_{t-2}, \ldots, C_{t-L}$, and $I_t$ is all information known or determined by the consumer at time $t$ (including $C_t$ and $V_t$). For external habits, we can write $\tilde{g}(I_t) = g(C_t, V_t)$, while for internal habits define

$$g(C_t, V_t) = E [\tilde{g}(I_t) \mid C_t, V_t].$$

With this notation, regardless of whether habits are internal or external, we may write the first order conditions associated with the Lagrangean (15) as

$$\lambda_t = b^t \tilde{g}(I_t)$$

$$\lambda_t = E [\lambda_{t+1} R_{j_{t+1}} \mid I_t] \quad j = 1, \ldots, J$$

$$\lambda_t P_t = b^t g_d(C_t, V_t) - \delta E [\lambda_{t+1} P_{t+1} \mid I_t]$$

Using the consumption equation $\lambda_t = b^t \tilde{g}(I_t)$ to remove the Lagrangeans in the assets and durables first order conditions gives

$$b^t \tilde{g}(I_t) = E [b^{t+1} \tilde{g}(I_{t+1}) R_{j_{t+1}} \mid I_t] \quad j = 1, \ldots, J$$

$$b^t \tilde{g}(I_t) P_t = b^t g_d(C_t, V_t) - \delta E [b^{t+1} \tilde{g}(I_{t+1}) P_{t+1} \mid I_t].$$

Taking the conditional expectation of the asset equations, conditioning on $C_t, V_t$, yields the Euler equations for asset $j$

$$g(C_t, V_t) = b E [g(C_{t+1}, V_{t+1}) R_{j_{t+1}} \mid C_t, V_t] \quad j = 1, \ldots, J,$$

for all $t$. Therefore, given the pair $(U, b)$ of utility function and discounting factor the optimal decision satisfies the Euler equations for all asset $j$.

### 9.2 Results on kernel estimators

#### 9.2.1 Assumptions

This section collects some results on kernel estimates, providing primitive assumptions for the general conditions of the main text. We need to introduce some notation from empirical processes theory.
To measure the complexity of the class $\mathcal{F}$, we can employ covering or bracketing numbers. Here, for simplicity, we focus on bracketing numbers. Given two functions $l, u$, a bracket $[l, u]$ is the set of functions $f \in \mathcal{F}$ such that $l \leq f \leq u$. An $\varepsilon$-bracket with respect to $\|\cdot\|$ is a bracket $[l, u]$ with $\|l - u\| \leq \varepsilon$, $\|l\| < \infty$ and $\|u\| < \infty$ (note that $u$ and $l$ not need to be in $\mathcal{F}$). The covering number with bracketing $N_h(\varepsilon, \mathcal{F}, \|\cdot\|)$ is the minimal number of $\varepsilon$-brackets with respect to $\|\cdot\|$ needed to cover $\mathcal{F}$. An envelope for $\mathcal{F}$ is a function $F$, such that $F(c, v) \geq \sup_{g \in \mathcal{F}} |g(c, v)|$ for all $(c, v)$. Denote by $\mathcal{K}(r)$ the class of $r$-order kernels $K$ that are Lipschitz continuous on the support $[-1, 1]$, symmetric, integrate to one, and such that for some $r \geq 2$: $\int u^r K(u) du = \delta_{l0}$ for $l = 0, \ldots, r - 1$, where $\delta_{l0}$ denotes Kronecker’s delta, and $\int u^r K(u) du \neq 0$.

**Assumption A1:**

1. For each $\varepsilon > 0$, $\log N_{\varepsilon}(\varepsilon, \mathcal{F}, \|\cdot\|) \leq C\varepsilon^{-v}$ for some $v < 2$. The class $\mathcal{F}$ is such that $g_0 \in \mathcal{F}$ and has an envelope $F$ such that $\sup_{(c, v) \in \mathcal{S}} E[|F(C', V') R'| |C = c, V = v] < \infty$ for some $\delta > 2$. Moreover,

$$
\lim_{\delta \to 0} \sup_{\|c_1 - c_2\| < \delta} \sup_{g \in \mathcal{F}} \|E[g(C', V') R' | C = c_1, V = v_1] - E[g(C', V') R' | C = c_2, V = v_2]\| = 0.
$$

2. The density function $f(\cdot)$ is bounded away from zero on $S_\mu$ and is continuous on $S$.

3. The kernel $K$ satisfies $K \in \mathcal{K}(2)$.

4. As $n \to \infty$, the possibly stochastic bandwidth $h \equiv h_n$ satisfies $P(l_n \leq h_n \leq u_n) \to 1$ for deterministic sequences of positive numbers $l_n$ and $u_n$ such that: $u_n \downarrow 0$ and $l_n^{(s)/2} n / \log n \to \infty$.

Examples of classes $\mathcal{F}$ satisfying A1.1 abound in the literature; see van der Vaart and Wellner (1996). The uniform equicontinuity in Assumption A1.1 follows under smoothness conditions on the joint density of $W$. The remaining conditions in Assumption A1 are self-explanatory. For A1.3 we could also use kernels with unbounded support that satisfy some smoothness and integrability conditions. Finally, note that A1.4 allows for data-driven bandwidth choices, which are common in applied work.

For asymptotic normality of our estimators we require the following assumption.

**Assumption A2:**

1. The density function $f$ satisfies $f \in C^r(S_\mu)$, where $r$ as in A2.4 below.

2. Functions in $R(A)$ are in $C^r(S_\mu)$ with uniformly equicontinuous $r$-th derivative on $S_\mu$. 

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3. The function \( s \) given in Definition 3 above satisfies \( s \in \mathcal{C}^r(S_\mu) \).

4. The kernel \( K \) satisfies \( K \in \mathcal{K}(r) \), for \( r \geq 2 \).

5. For \( l_n \) and \( u_n \) defined in A1.5, it also holds that \( l_n^2 n / \log n \to \infty \) and \( n u_n^{2r} \to 0 \) as \( n \to \infty \).

### 9.2.2 Some generic results

Let \( Y \in \mathbb{R}^{d_Y} \) denote the elements of \( (C',V') \) that do not overlap with \( (C,V) \), let \( X \in \mathbb{R}^{d_X} \) denote the overlapping elements of \( (C',V') \) and \( (C,V) \) and let \( Z \in \mathbb{R}^{d_Z} \) denote the elements of \( (C,V) \) that do not overlap with \( (C',V') \). We denote by \( \psi \equiv (\varphi,c,v) \) a generic element of the set \( \Psi \equiv \mathcal{F} \times S_\mu \). Let \( f(c,v) \) denote the density of \( (C,V) \) evaluated at \( (c,v) \). Define the regression function \( m(\psi) \equiv E[\varphi(C',V')|C = c, V = v] \). Then, an estimator for \( m(\psi) \) is given by

\[
\hat{m}_h(\psi) = \frac{1}{nh} \sum_{i=1}^{n} \varphi(C'_i, V'_i) R_i K \left( \frac{c - C_i}{h} \right) \prod_{j=1}^{d_Y} K \left( \frac{ v_j - V_{ji} }{h} \right) = \frac{\hat{T}_h(\psi)}{\hat{f}(c,v)}.
\]

Henceforth, we abstract from measurability issues that may arise (see van der Vaart and Wellner (1996) for ways to deal with lack of measurability). The following lemma draws heavily on Einmahl and Mason (2005), who in their Theorem 3 and Corollary 3 proved the univariate version of it (i.e. the case \( d_Y = 1 \)).

**Lemma B1.** Suppose that Assumption A1 holds. Then, a.s.

\[
\sup_{l_n \leq h \leq u_n, \psi \in \Psi} |\hat{m}_h(\psi) - m(\psi)| = o(1). \tag{17}
\]

If, in addition, A2 holds, then a.s.

\[
\sup_{l_n \leq h \leq u_n, \psi \in \Psi} |\hat{m}_h(\psi) - m(\psi)| = O \left( \sqrt{\frac{\ln n}{m_n^r} + u_n^r} \right). \tag{18}
\]

**Proof.** By the Triangle inequality

\[
\left| \hat{m}_h(\psi) - m(\psi) \right| \leq \left| \hat{m}_h(\psi) - E[\hat{T}_h(\psi)] \right| + \left| E[\hat{T}_h(\psi)] - m(\psi) \right| \leq \frac{1}{\hat{f}(c,v)} \left| \hat{T}_h(\psi) - E[\hat{T}_h(\psi)] \right| + \frac{\left| E[\hat{T}_h(\psi)] \right|}{\hat{f}(c,v) E[\hat{f}(c,v)]} \left| \hat{f}(c,v) - E[\hat{f}(c,v)] \right|
\]

\[
+ \frac{1}{E[\hat{f}(c,v)]} \left| E[\hat{T}_h(\psi)] - T(\psi) \right| + \frac{|T(\psi)|}{E[\hat{f}(c,v)] f(c,v)} \left| E[\hat{f}(c,v)] - f(c,v) \right|,
\]

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where $T(\psi) \equiv m(\psi)f(c,v)$. We shall apply a variation of Theorem 4 in Einmahl and Mason (2005) to obtain uniform rates for $\hat{T}_h(\psi) - E[\hat{T}_h(\psi)]$; the rates for $\hat{f}(c,v) - E[\hat{f}(c,v)]$ follow analogously and are simpler to obtain (see their Theorem 1, p. 1.3). Our conditions imply the assumptions needed for their Theorem 4, where the bracketing conditions replace their covering conditions (see their Remark 3 and Lemma B.4 in Escanciano, Jacho-Chávez and Lewbel (2014)). Then, we conclude a.s.

$$\sup_{l_n \leq h \leq u_n} \sup_{\psi \in \Psi} |\hat{T}_h(\psi) - E[\hat{T}_h(\psi)]| = O\left(\frac{\ln n}{n^{1/2}}\right).$$

On the other hand, Lemma 2 in Einmahl and Mason (2005) and the uniform equicontinuity of Assumption A2.2 yield

$$\sup_{l_n \leq h \leq u_n} \sup_{\psi \in \Psi} |E[\hat{T}_h(\psi)] - T(\psi)| = o(1),$$

and likewise for the density bias term. This together with the above expansion for $\hat{m}_h - m$ completes the proof of (17).

To obtain rates for the bias terms we need the smoothness conditions of Assumption A2. A standard Taylor expansion argument, the higher-order property of the kernel and the uniform equicontinuity of the $r$–th derivative imply that

$$\sup_{l_n \leq h \leq u_n} \sup_{\psi \in \Psi} |E[\hat{T}_h(\psi)] - T(\psi)| = O\left(u_n^r\right),$$

and similarly for the density bias term. The proof is completed by standard arguments using the boundedness away from zero of $f(c,v)$ over the domain.

**Lemma B2.** Suppose that Assumption A1 holds. Then, as $n \to \infty$, a.s.

$$\left\|\hat{A} - A\right\|_F = \sup_{g \in \mathcal{F} : \|g\| \leq 1} \left\|\hat{A}g - Ag\right\| = o(1).$$

If, in addition, A2 holds, then a.s.

$$\sup_{l_n \leq h \leq u_n} \left\|\hat{A} - A\right\|_F = O\left(\frac{\ln n}{n^{1/2}} + u_n^r\right).$$

**Proof.** The result follows directly from Lemma B1.

We introduce a useful class of functions:

**Definition 4.** Let $\mathcal{L}^2(r)$ be the class of functions $\varphi \in \mathcal{L}^2$ such that $\Sigma_\varphi \equiv E[\varphi_i^2 c_i^2] < \infty$ and $\varphi$ is $r$–times continuously differentiable.
\textbf{Lemma B3.} Suppose that Assumptions A1 and A2 hold. Then, for any \( \varphi \in L^2(r) \), it holds that

\[
\sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, \varphi \right\rangle = \frac{1}{\sqrt{n}} \sum_{i=1}^{n} \varphi_i \varepsilon_i + o_P(1),
\]

and then

\[
\sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, \varphi \right\rangle \xrightarrow{d} N(0, \Sigma_\varphi).
\]

\textbf{Proof.} Define

\[
\hat{T} g_0(c,v) = \frac{1}{n} \sum_{i=1}^{n} g_{0i} R_i' K_{hi} (c,v),
\]

with \( g_{0i} \equiv g_0 (C'_i, V'_i) \) and note that \( \hat{A} g_0(c,v) = \hat{T} g_0(c,v) / \hat{f}(c,v) \). Using standard arguments, we write

\[
\left( \hat{A} - A \right) g_0(c,v) = a_n(c,v) + r_n(c,v),
\]

where

\[
a_n(c,v) = f^{-1}(c,v) \left( \hat{T} g_0(c,v) - T g_0(c,v) - A g_0(c,v) \left( \hat{f}(c,v) - f(c,v) \right) \right),
\]

\[
T g_0(c,v) \equiv f(c,v) A g_0(c,v), \quad \hat{T} g_0(c,v) \equiv \hat{f}(c,v) \hat{A} g_0(c,v) \quad \text{and}
\]

\[
r_n(c,v) \equiv - \frac{\hat{f}(c,v) - f(c,v)}{\hat{f}(c,v)} a_n(c,v).
\]

Lemma B1 and our conditions on the bandwidth imply \( \|r_n\| = o_P(n^{-1/2}) \). It then follows that \( \left\langle \left( \hat{A} - A \right) g_0, \varphi \right\rangle \) has the following expansion

\[
\int \varphi(c,v)[\hat{T} g_0(c,v) - T g_0(c,v)]dcdv \tag{19}
\]

\[
- \int \varphi(c,v) A g_0(c,v) [\hat{f}(c,v) - f(c,v)]dcdv \tag{20}
\]

\[
+ o_P(n^{-1/2}).
\]

We now look at terms (19)-(20). Firstly, it follows from standard arguments and A2.5 that the difference between \( T g_0(c,v) \) and \( E[\hat{T} g_0(c,v)] \) is \( O_P(n^{-1/2}) \) by the condition \( nu_n^2 \to 0 \). Hence,

\[
\int \varphi(c,v)[\hat{T} g_0(c,v) - T g_0(c,v)]dcdv = \int \varphi(c,v)[\hat{T} g_0(c,v) - E(\hat{T} g_0(c,v))]dcdv + o_P(n^{-1/2})
\]

\[
= \frac{1}{n} \sum_{i=1}^{n} g_{0i} R_i' \int \varphi(c,v) K_{hi} (c,v) dcdv - \int \varphi(c,v) E(g_{0i}' R_i K_{hi} (c,v))dcdv + o_P(n^{-1/2}),
\]

\[
= \frac{1}{n} \sum_{i=1}^{n} \varphi(C_i, V_i) g_{0i}' R_i - E[\varphi(C_i, V_i) A g_0(C_i, V_i)] + o_P(n^{-1/2}),
\]

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where the last equality follows from the standard change of variables argument and our Assumption A2. Likewise, the term (20) becomes $n^{-1/2} \sum_{i=1}^{n} \varphi(C_i, V_i)A g_0 (C_i, V_i) - E[\varphi (C_i, V_i) A g_0 (C_i, V_i)] + o_p(n^{-1/2})$. In conclusion, we have

$$\sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, \varphi \right\rangle = \frac{1}{\sqrt{n}} \sum_{i=1}^{n} \varphi(C_i, V_i) \varepsilon_i + o_p(1).$$

Then, the result follows from a standard central limit theorem, since $\{\varphi(C_i, V_i) \varepsilon_i\}_{i=1}^{n}$ is iid with zero mean and finite variance. ■

For a generic function $r \in L^2$, define

$$r_s = r - \langle g_0, r \rangle (g_0, s)^{-1} s.$$

Also for $r \in \mathcal{N}^\perp(L) = \mathcal{R}(L^*)$ denote by $r^*$ the unique minimum norm solution of $r = L^* r^*$. Note that for $r \in \mathcal{R}(L^*)$, $r^*_s$ does not depend on the solution $r^*$ considered of $r = L^* r^*$ (whether or not is minimum norm). This follows because under our conditions $\mathcal{N}(L^*)$ is the linear span generated by $s$.

**Lemma B4.** Let Assumptions S, C, I, E, N and A1-A2 hold. If $\varphi \in \mathcal{N}^\perp(L)$, so $\varphi = L^* \varphi^*$ for some $\varphi^*$, and if $\varphi_s^* \in L^2(r)$, then

$$\sqrt{n} \left\langle \hat{g} - g_0, \varphi \right\rangle \xrightarrow{d} N \left( 0, b_0^2 \Sigma_{\varphi_s^*} \right).$$

**Proof.** Note that by (23) below and the adjoint property

$$\sqrt{n} \left\langle \hat{g} - g_0, \varphi \right\rangle = \sqrt{n} \left\langle \hat{g} - g_0, L^* \varphi^* \right\rangle$$

$$= \sqrt{n} \langle L(\hat{g} - g_0), \varphi^* \rangle$$

$$= -\sqrt{n} \left( \hat{b} - b_0 \right) b_0^{-1} \langle g_0, \varphi^* \rangle - b_0 \sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, \varphi^* \right\rangle + o_p(1).$$

Then, by the proof of Theorem 4, this can be further simplified to

$$b_0 \sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, s \langle g_0, \varphi^* \rangle - \varphi^* \right\rangle = -b_0 \sqrt{n} \left\langle \left( \hat{A} - A \right) g_0, \varphi_s^* \right\rangle + o_p(1).$$

Then, the result follows from the last display and Lemma B3. ■

### 9.3 Main Proofs

The spectral radius $\rho(A)$ of a linear continuous operator $A$ on a Banach space $X$ is defined as $\sup_{\lambda \in \sigma(A)} |\lambda|$, where $\sigma(A) \subset \mathbb{C}$ denotes the spectrum of $A$. Any compact operator $A$ has a discrete
spectrum, so that \( \sigma(A) \) is simply the set of eigenvalues of \( A \). For more definitions and further details see Kress (1999, Chapter 3.2). The operator \( B \) is called positive if \( Bg \in \mathcal{P} \) when \( g \in \mathcal{P} \).

**Proof of Theorem 1.** By Assumption C the set of countable eigenvalues of \( A \) has zero as a limit point, and thus, the set of eigenvalues \( \lambda \) with \( \lambda^{-1} \in (0,1) \) is a finite set. By Theorem 3.1 in Kress (1999) for each such eigenvalue there is a finite-dimensional eigenvector space. ■

**Proof of Theorem 2.** Let \( A^* \) denote the adjoint of \( A \), which is also compact and positive by well known results in functional analysis. Assumption S implies that \( \rho(A) > 0 \). Also notice that the eigenvalues of \( A^* \) are complex conjugates of those of \( A \) (in particular, \( \rho(A) = \rho(A^*) \)). Then, by the Krein-Rutman’s theorem (see Theorem 7.C in Zeidler (1986, vol. 1, p. 290)) there is exactly one solution to \( bA g = g \) with \( g \in \mathcal{P} \) and \( \|g\| = 1 \) and a solution to \( bA^* s = s \) with \( s > 0 \). Note \( \langle g, s \rangle = b \langle Ag, s \rangle = b \langle g, A^* s \rangle = b\rho(A) \langle g, s \rangle \). Hence, since \( g \) and \( s \) are strictly positive, \( \langle g, s \rangle \neq 0 \), and then \( b = \rho^{-1}(A) \). ■

**Proof of Theorem 3.** By Theorems 1 and 2 in Osborn (1975), there is a constant \( M \) such that

\[
\left| \widehat{b}^{-1} - b_0^{-1} \right| \leq M \| \widehat{A} - A \|_{G_0} \tag{21}
\]

and

\[
\| \widehat{g} - \tilde{g} \| \leq M \| \widehat{A} - A \|_{G_0}, \tag{22}
\]

where \( \tilde{g} = \langle \tilde{g}, g_0 \rangle g_0 \) is the projection of \( \tilde{g} \) on \( g_0 \). Thus, by \( 0 < b_0, \hat{b} < 1 \), a.s,

\[
\left| \hat{b} - b_0 \right| \leq M \left| \hat{b} \times b_0 \right| \| \widehat{A} - A \|_{G_0} \leq M \| \widehat{A} - A \|_{G_0},
\]

and by Assumption E.2 \( \left| \hat{b} - b_0 \right| = o(1) \) a.s.

To conclude that \( \| \tilde{g} - g_0 \| = o(1) \) a.s. we need to show that \( \| \tilde{g} - g_0 \| = o(1) \) a.s. First, we show that \( \langle \tilde{g}, g_0 \rangle \) is positive for sufficiently large \( n \). To see this, note by the normalization, (22) and the Law of Large Numbers (LLN), a.s.

\[
\langle \tilde{g}, 1 \rangle_n = \langle \tilde{g}, 1 \rangle_n + o(1) = \langle \tilde{g}, g_0 \rangle \langle g_0, 1 \rangle_n + o(1) > 0,
\]

so \( \langle \tilde{g}, g_0 \rangle > 0 \) for large enough \( n \).
Next,

\[ 1 = \|\hat{g}\| + o(1) \text{ (by normalization and LLN)} \]
\[ = \|\hat{g}\| + o(1) \text{ (by } \|\hat{g} - g\| \leq M\|\hat{A} - A\|_{G_0}) \]
\[ = |\langle \hat{g}, g_0 \rangle| + o(1), \text{ (by definition of } \hat{g}) \]

which then implies \( \|\hat{g} - g_0\| = |\langle \hat{g}, g_0 \rangle - 1| = o(1) \) \text{ a.s. Hence, by the triangle inequality, } \|\hat{g} - g_0\| = o(1) \text{ a.s.} \]

**Proof of Theorem 4.** By definition

\[ \hat{b}\hat{A}\hat{g} - b_0Ag_0 = \hat{g} - g_0. \]

Write the left hand side of the last display as

\[ \left( \hat{b} - b_0 \right) A\hat{g} + b_0 \left( \hat{A} - A \right) g_0 + b_0A(\hat{g} - g_0) + \hat{R}, \]

where \( \hat{R} = \left( \hat{b} - b_0 \right) \left( \hat{A} - A_0 \right) \hat{g} + b_0 \left( \hat{A} - A \right) (\hat{g} - g_0). \) Then, after noticing that (by definition of \( s \)),

\[ \langle b_0A(\hat{g} - g_0), s \rangle = \langle \hat{g} - g_0, s \rangle, \]

we obtain

\[ \left( \hat{b} - b_0 \right) b_0^{-1} \langle \hat{g}, s \rangle + b_0 \left( \left( \hat{A} - A \right) g_0, s \right) + \langle \hat{R}, s \rangle = 0. \]

By the proof of Theorem 3, it is straightforward to show that, for a \( C > 0, \)

\[ \|\hat{R}\| \leq C \left\{ \|\hat{A} - A\|_{G_0}^2 + \|\hat{A} - A\|_{G_0 - \{g_0\}} \|\hat{g} - g_0\| \right\} \]

and

\[ \|\hat{g} - g_0\| \leq \|\hat{g} - \bar{g}\| + \|\bar{g} - g_0\| \]
\[ \leq M\|\hat{A} - A\|_{G_0} + \|\langle \hat{g}, g_0 \rangle - 1| \]
\[ \leq 2M\|\hat{A} - A\|_{G_0} + \|\hat{g}\| - \|\bar{g}\|, \]

which implies by Assumption N.1 and \( \|\hat{g}\| - \|\hat{g}\|_{n} = O_P(n^{-1/2}) \)

\[ \|\hat{R}\| = o_P(n^{-1/2}). \]

Then, Cauchy-Schwarz inequality yields

\[ \left| \langle \hat{R}, s \rangle \right| \leq \|\hat{R}\| \|s\| \]
\[ = o_P(n^{-1/2}). \]
Then, by continuity of the inner product, \( \langle \hat{g}, s \rangle \to_p \langle g_0, s \rangle \equiv 1 \), and by Slutzky Theorem
\[
\sqrt{n} \left( \hat{b} - b_0 \right) = -\sqrt{n} b_0^2 \left( \left( \hat{A} - A \right) g_0, s \right) + o_P(1).
\]
Hence, the result follows from Assumption N.2.

**Proof of Theorem 5.** Define the operators \( L = b_0 A - I \), and its estimator \( \hat{L} = \hat{b} \hat{A} - I \). Then, by definition
\[
0 = \hat{L} \hat{g} - L g_0 = L(\hat{g} - g_0) + (\hat{L} - L)(\hat{g} - g_0).
\]
First, from previous results it is straightforward to show as in Theorem 4
\[
\mathbb{P} \left( \frac{1}{n} \sum_{i=1}^{n} \left( \hat{g}_i - g_0 \right) \right) = o_P(n^{-1/2})
\]
and
\[
\mathbb{P} \left( \frac{1}{n} \sum_{i=1}^{n} \left( \hat{L} - L \right) g_0 - b_0(\hat{A} - A) g_0 \right) = O_P(n^{-1/2}).
\]
Hence, in \( L^2 \),
\[
L(\hat{g} - g_0) = -b_0(\hat{A} - A) g_0 + R_n,
\]where \( R_n \) satisfies the conditions of the Theorem.

**Proof of Theorem 6.** Set \( \hat{\zeta}(C_i, V_i) = -C_i \partial \hat{g}(C_i, V_i) \partial c / \hat{g}(C_i, V_i) \), which estimates consistently \( \zeta(C_i, V_i) = -C_i (\partial g_0(C_i, V_i) / \partial c) / g_0(C_i, V_i) \). Then, using standard empirical processes notation, write
\[
\sqrt{n} \left( \gamma_n(\hat{g}) - \gamma(g_0) \right) = \sqrt{n} \left( P_n \hat{\zeta} - P \hat{\zeta} \right) + \sqrt{n} \left( P \hat{\zeta} - P \zeta \right).
\]
By the \( P \)-Donsker property of \( D \), \( P (\hat{g} \in G) \to 1 \) and the consistency of \( \hat{g} \),
\[
\sqrt{n} \left( P_n \hat{\zeta} - P \hat{\zeta} \right) = \sqrt{n} \left( P_n \zeta - P \zeta \right) + o_P(1).
\]
Since \( \hat{g} - g_0 \) is bounded with probability tending to one, we can apply integration by parts and use Assumption CE to write
\[
\sqrt{n} \left( P \hat{\zeta} - P \zeta \right) = \sqrt{n} \left( \log(\hat{g}) - \log(g_0), d \right) + o_P(1)
\]where the last equality follows from the Mean Value Theorem and the lower bounds on \( g \) and \( \hat{g} \). Note that \( \chi \in N^\perp(L) \), since \( \langle g_0, \chi \rangle = E[d(C, V)] = 0 \). Then, by Lemma B4
\[
\sqrt{n} \left( P \hat{\zeta} - P \zeta \right) = -\frac{b_0}{\sqrt{n}} \sum_{i=1}^{n} \chi_{s}^*(C_i, V_i) \xi_i + o_P(1),
\]
and therefore

\[
\sqrt{n} (\gamma_n (\hat{g}) - \gamma (g_0)) = \frac{1}{\sqrt{n}} \sum_{i=1}^{n} (\zeta(C_i, V_i) - P\zeta) - b_0 \chi_s^*(C_i, V_i) \epsilon_i + o_P(1).
\]

The result then follows from the Lindeberg-Levy central limit theorem and \(E[\epsilon_i | C_i, V_i] = 0\). □
References


