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**ABSTRACT**

We consider a general class of nonlinear optimal policy problems involving forward-looking constraints (such as the Euler equations that are typically present as structural equations in DSGE models), and show that it is possible, under regularity conditions that are straightforward to check, to derive a problem with linear constraints and a quadratic objective that approximates the exact problem. The LQ approximate problem is computationally simple to solve, even in the case of moderately large state spaces and flexibly parameterized disturbance processes, and its solution represents a local linear approximation to the optimal policy for the exact model in the case that stochastic disturbances are small enough. We derive the second-order conditions that must be satisfied in order for the LQ problem to have a solution, and show that these are stronger, in general, than those required for LQ problems without forward-looking constraints. We also show how the same linear approximations to the model structural equations and quadratic approximation to the exact welfare measure can be used to correctly rank alternative simple policy rules, again in the case of small enough shocks.

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Linear-quadratic (LQ) optimal-control problems have been the subject of an extensive literature.<sup>1</sup> General characterizations of their solutions and useful numerical algorithms to compute them are now available, allowing models with fairly large state spaces, complicated dynamic linkages, and a range of alternative informational assumptions to be handled.<sup>2</sup> And the extension of the classic results of the engineering control literature to the case of forward-looking systems of the kind that naturally arise in economic policy problems when one allows for rational expectations on the part of the private sector has proven to be fairly straightforward.<sup>3</sup>

An important question, however, is whether optimal policy problems of economic interest should take this convenient form. It is easy enough to apply LQ methodology if one specifies an *ad hoc* quadratic loss function on the basis of informal consideration of the kinds of instability in the economy that one would like to reduce, and posits linear structural relations that capture certain features of economic time series without requiring these relations to have explicit choice-theoretic foundations, as in early applications to problems of monetary policy.<sup>4</sup> But it is highly unlikely that the analysis of optimal policy in a DSGE model will involve either an exactly quadratic utility function or exactly linear constraints.

We shall nonetheless argue that LQ problems can usefully be employed as approximations to exact optimal policy problems in a fairly broad range of cases. Since an LQ problem necessarily leads to an optimal decision rule that is linear, the most that one could hope to obtain with any generality would be for the solution to the LQ problem to represent a *local linear approximation* to the actual optimal policy — that is, a first-order Taylor approximation to the true, nonlinear optimal policy rule. In this paper we present conditions under which this will be the case, and show how to derive an LQ approximate problem corresponding to any member of a general class of optimal policy problems.

The conditions under which the solution to an LQ approximate problem will yield a correct local linear approximation to optimal policy are in fact more restrictive

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<sup>1</sup>Important references include Bertsekas (1976), Chow (1975), Hansen and Sargent (2004), Kendrick (1981), Kwakernaak and Sivan (1972), and Sargent (1987). See Kendrick (2005) for an overview of the use of LQ methods in economics.

<sup>2</sup>For numerical algorithms see, among others, Amman (1996), Anderson *et al.* (1996), Amman and Kendrick (1999), Diaz-Gimenez (1999), Gerali and Lippi (2005), Hansen and Sargent (2004), and Söderlind (1999).

<sup>3</sup>See, *e.g.*, Backus and Driffill (1986) for a useful review.

<sup>4</sup>Notable examples include Kalchbrenner and Tinsley (1975) and Leroy and Waud (1977).

than might be expected, as noted for example by Judd (1996, pp. 536-539; 1998, pp. 507-508). In particular, it does *not* suffice that the objective and constraints of the exact problem be continuously differentiable a sufficient number of times, that the solution to the LQ approximate problem imply a stationary evolution of the endogenous variables, and that the exogenous disturbances be small enough (though each of these conditions is obviously *necessary*, except in highly special cases). An approach that simply computes a second-order Taylor-series approximation to the utility function and a first-order Taylor-series approximation to the model structural relations in order to define an approximate LQ problem — the approach criticized by Judd (1996, 1998) that we have elsewhere (Benigno and Woodford, 2006a) called “naive LQ approximation” — may yield a linear policy rule with coefficients very different from those of a correct linear approximation to the optimal policy in the case of small enough disturbances.<sup>5</sup>

The discussion by Judd (1996, pp. 536-539) might seem to imply that LQ approximation is an inherently mistaken idea — that it cannot be expected, other than in cases so special as to represent an essentially fortuitous result, to yield a correct approximation to optimal policy at all. Nonetheless, it is quite generally possible to construct an alternative quadratic objective function that *will* result in a correct local LQ approximation, in the sense that the linear solution to the LQ problem is a correct linear approximation to the solution to the exact problem. The correct method was illustrated in the important paper of Magill (1977), that applied results of Fleming (1971) from the optimal-control literature to derive a local LQ approximation to a continuous-time multi-sector optimal growth model. Here we show how the method of Magill can be used in the context of discrete-time dynamic optimization problems where some of the structural relations are forward-looking, as is almost inevitably the case in optimal monetary or fiscal policy problems.<sup>6</sup>

Of course the problems that can arise as a result of “naive” LQ optimization can also be avoided through the use of alternative perturbation techniques, as explained by Judd. Approaches that are widely used in the recent literature on policy

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<sup>5</sup>For an example illustrating this possibility, see Benigno and Woodford (2006a). The same problem can also result in incorrect welfare rankings of alternative simple policies, as discussed by Kim and Kim (2003, 2006).

<sup>6</sup>See also Levine *et al.* (2007) for another application to a discrete-time problem, and additional discussion of how our method relates to that of Magill.

analysis in DSGE models include either (i) deriving the first-order conditions that characterize optimal (Ramsey) policy using the exact (nonlinear) objective and constraints, and then log-linearizing *these conditions* in order to obtain an approximate solution to them, rather than separately approximating the objective and constraints before deriving the first-order conditions;<sup>7</sup> or (ii) obtaining a higher-order (at least second-order) perturbation solution for the equilibrium implied by a given policy by solving a higher-order approximation to the constraints, and then evaluating welfare under the policy using this approximate solution.<sup>8</sup> These methods can also be used to correctly calculate a linear approximation to the optimal policy rule, and when applied to the problem considered here, provide alternative approaches to calculating the same solution.<sup>9</sup>

Despite the existence of these alternative perturbation approaches to the analysis of optimal policy, we believe that it remains useful to show how a correct form of LQ analysis is possible in the case of a fairly general class of problems. One reason is that the ability to translate a policy problem into this form allows one to use the extensive body of theoretical analysis and numerical techniques that have been developed for LQ problems. Another is that casting optimal policy analysis in DSGE models in this form can allow comparisons between welfare-based policy analysis and analyses of optimal policy based on *ad hoc* stabilization objectives (which have often been expressed as LQ problems). We also show that the LQ formulation of the approximate policy problem makes it straightforward to analyze whether a solution to first-order conditions for optimal policy also satisfies the relevant second-order conditions for optimality, and to rank suboptimal policy rules by a criterion that is consistent with the characterization given of optimal policy.

We first explain the essential problem with naive LQ optimization in section 1, in the context of a simple, finite-dimensional example, and also use this example to

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<sup>7</sup>Recent applications of this method to problems of optimal monetary and fiscal policy include King and Wolman (1999), Khan *et al.* (2003) and Schmitt-Grohé and Uribe (2004b).

<sup>8</sup>For discussions of methods for executing computations of this kind in general classes of forward-looking equation systems see, among others, Jin and Judd (2002), Kim *et al.* (2003), and Schmitt-Grohé and Uribe (2004a). These methods have been used in many recent numerical analyses of optimal policy (e.g., Schmitt-Grohé and Uribe, 2007).

<sup>9</sup>As shown in section 2.3 below, our method computes the same coefficients for a linear policy rule as are obtained by linearization of the first-order conditions for the exact policy problem. The general intuition for this result is discussed in section 1.

explain why the approach used by Magill (1977) avoids the problem. In section 2, we present a general class of dynamic optimization problems with forward-looking constraints, and derive an LQ approximate problem associated with any problem in this class. Section 3 discusses the general algebraic form of the first- and second-order conditions for optimality in the LQ approximate problem. Section 4 shows how the quadratic objective for stabilization policy derived in section 2 can also be used to compute welfare comparisons between alternative sub-optimal policies, in the case that the stochastic disturbances are small enough. Finally, section 5 discusses applications of the method described here and concludes.

## 1 Naive and Correct LQ Approximations

Here we review the reason why naive LQ approximation is generally incorrect, as noted by Judd (1996, 1998), in the context of a simple static optimization problem that allows us to explain the issues in terms of simple multivariate calculus. We then illustrate how Magill's (1977) approach solves the problem, in the context of this static example, before turning in the next section to the additional complications raised by dynamic problems.

### 1.1 A Static Example

Suppose that we wish to find the policy  $y(\xi)$  that maximizes an objective  $U(y; \xi)$ , where  $y$  is an  $n$ -vector of endogenous variables and  $\xi$  is a vector of exogenous disturbances; we assume that  $U$  is at least twice continuously differentiable with respect to the arguments  $y$ . Suppose furthermore that the possible outcomes  $y$  that can be achieved by policy in any state of the world  $\xi$  are those values consistent with the structural equations

$$F(y; \xi) = 0, \tag{1.1}$$

where  $F$  is a vector of  $m$  functions (for some  $m < n$ ), again each at least twice continuously differentiable. We assume that  $m < n$  so that there is at least one direction in which it is possible for the outcome  $y$  to be varied by policy. We might suppose that  $y$  is determined by equations (1.1) together with an additional set of  $n - m$  equations of the form

$$G(y; i, \xi) = 0, \tag{1.2}$$

where  $i$  is a vector of  $n - m$  instrument settings (or control variables); but the nature of the additional equations (1.2) does not matter for our conclusions below, as long as the derivative matrices

$$\begin{bmatrix} D_y F \\ D_y G \end{bmatrix}, \quad D_i G$$

are of full rank when the partial derivatives are evaluated at the point around which we conduct our local analysis. We shall suppose that there exists a solution  $y^{opt}(\xi)$  to this problem for all  $\xi$  in some neighborhood of 0 (the case of “zero disturbances”).

Now let  $\bar{y}$  be the outcome under an optimal policy in the case that  $\xi = 0$ ; that is, it maximizes  $U(y; 0)$  subject to the constraints  $F(y; 0) = 0$ .<sup>10</sup> We wish to obtain a local linear approximation to the function  $y^{opt}(\xi)$  for values of  $\xi$  near enough to 0. In the case that  $y^{opt}(\xi)$  is differentiable at  $\xi = 0$ , such a linear approximation exists, with coefficients of the linear rule given by the derivatives of  $y^{opt}$ , since by Taylor’s theorem,

$$y^{opt}(\xi) = \bar{y} + D y^{opt} \cdot \xi + \mathcal{O}(\|\xi\|^2), \quad (1.3)$$

where the partial derivatives are evaluated at  $\xi = 0$ .

In many problems, differentiability can be established, and the derivatives (and hence the coefficients of the linear approximation) calculated, using the implicit function theorem. We can write a Lagrangian for this problem

$$\mathcal{L}(y; \xi; \lambda) \equiv U(y; \xi) + \lambda' F(y; \xi) \quad (1.4)$$

where  $\lambda$  is an  $m$ -vector of multipliers associated with the constraints (1.1); there must exist a vector of multipliers for which the optimal policy minimizes the Lagrangian for each possible value of  $\xi$ . It follows that the optimal policy  $y^{opt}(\xi)$  must satisfy the (exact, nonlinear) first-order conditions obtained by differentiating the Lagrangian,

$$D_y U(y; \xi) + \lambda' D_y F(y; \xi) = 0, \quad (1.5)$$

in addition to the structural relations. The system consisting of (1.1) together with (1.5) is then a system of  $n + m$  nonlinear equations implicitly defining functions  $y(\xi)$  and  $\lambda(\xi)$ . A correct local approximation to the solution to these equations can be

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<sup>10</sup>Note that we must compute our local approximations to the objective and constraints around this optimal point if there is to be any hope that consideration of these local approximations alone can correctly identify the optimal policy rule even in the case that  $\xi$  is small.

obtained (under the regularity condition stated below) by linearizing equations (1.1) and (1.5) around the unperturbed solution  $y(0) = \bar{y}$ ,  $\lambda(0) = \bar{\lambda}$ , and solving these linear equations for  $y$  and  $\lambda$  as linear functions of  $\xi$ .

In this method, we replace the exact constraints (1.1) by their linearized form,

$$D_y F \cdot \tilde{y} + D_\xi F \cdot \xi = 0, \quad (1.6)$$

where we use the notation  $\tilde{y} \equiv y - \bar{y}$ , and partial derivatives are evaluated at  $\bar{y}$ . Similarly, the linearization of the first-order conditions (1.5) is given by

$$\tilde{y}' D_{yy}^2 U + \xi' D_{\xi y}^2 U + \tilde{\lambda}' D_y F + \sum_k \bar{\lambda}_k [\tilde{y}' D_{yy}^2 F^k + \xi' D_{\xi y}^2 F^k] = 0, \quad (1.7)$$

where  $\tilde{\lambda} \equiv \lambda - \bar{\lambda}$ , and  $k$  indexes the  $m$  individual constraints  $F^k$ . The linear system consisting of (1.6)–(1.7) has a unique solution if and only if

$$\det \begin{bmatrix} D_{yy}^2 U + \sum_k \bar{\lambda}_k D_{yy}^2 F^k & (D_y F)' \\ D_y F & 0 \end{bmatrix} \neq 0. \quad (1.8)$$

This is also precisely the regularity condition under which the implicit function theorem guarantees that there exists a differentiable solution  $y(\xi)$  to the system consisting of (1.1) and (1.5), for values of  $\xi$  in a neighborhood of 0. Moreover, a local linear approximation of the form (1.3) exists, equal precisely to the solution to the linear system (1.6)–(1.7).

We wish to compare this correct linear approximation with the linear solution to an LQ optimization problem obtained by approximating the objective  $U$  and the constraints (1.1). In the case of any policy  $y(\xi)$  such that  $\tilde{y} = \mathcal{O}(\|\xi\|)$ ,<sup>11</sup> a second-order Taylor series expansion of  $U$  yields

$$\begin{aligned} U(y; \xi) &= \bar{U} + D_y U \cdot \tilde{y} + D_\xi U \cdot \xi + \frac{1}{2} \tilde{y}' D_{yy}^2 U \cdot \tilde{y} + \\ &\quad \frac{1}{2} \xi' D_{\xi\xi}^2 U \cdot \xi + \tilde{y}' D_{y\xi}^2 U \cdot \xi + \mathcal{O}(\|\xi\|^3) \\ &= D_y U \cdot \tilde{y} + \frac{1}{2} \tilde{y}' D_{yy}^2 U \cdot \tilde{y} + \tilde{y}' D_{y\xi}^2 U \cdot \xi + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3), \end{aligned} \quad (1.9)$$

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<sup>11</sup>Note that in the case that (1.8) holds and a local characterization of optimal policy can be given using the implicit function theorem, as discussed in the previous paragraph, the optimal policy  $y^{opt}(\xi)$  has this property. More generally, in our discussion below of the use of local approximations to rank alternative policies, we shall restrict attention to policies with this property.



where the various matrices of partial derivatives are each evaluated at  $(\bar{y}; 0)$ . The expression “t.i.p.” refers to terms that are independent of the policy chosen (such as the constant term and terms that depend only on the exogenous disturbances); the form of these terms is irrelevant in obtaining a correct ranking of alternative policies.

A naive LQ approximation of this problem can then be obtained by replacing the exact objective  $U(y; \xi)$  by the quadratic objective

$$U^Q(y; \xi) \equiv D_y U \cdot \tilde{y} + \frac{1}{2} \tilde{y}' D_{yy}^2 U \cdot \tilde{y} + \tilde{y}' D_{y\xi}^2 U \cdot \xi, \quad (1.10)$$

and replacing the exact constraints (1.1) by their linearized form (1.6). We wish to consider whether the policy that maximizes  $U^Q(y; \xi)$  subject to the constraints (1.6) represents a correct local linear approximation to the true optimal policy of the form (1.3).

In general, it does not. The policy that maximizes the naive quadratic objective (1.10) subject to the linearized constraints (1.6) satisfies linear first-order conditions

$$D_y U + \tilde{y}' D_{yy}^2 U + \xi' D_{\xi y}^2 U + \lambda' D_y F = 0. \quad (1.11)$$

The naive LQ-optimal policy is then obtained by solving the system of equations consisting of (1.6) and (1.11) for  $y$  and  $\lambda$  as linear functions of  $\xi$ . Because the two final terms on the left-hand side of (1.7) are missing in (1.11), the naive method will generally yield incorrect coefficients for the linear policy rule.

The fact that LQ analysis using the quadratic objective (1.10) yields an incorrect linear approximation to optimal policy is related to the fact that a linear approximation to the equilibrium outcome under a given policy rule does not suffice for a correct welfare ranking of alternative policies, even to second order. In the case of any (sufficiently differentiable) policy  $y(\xi)$  that is optimal in the absence of disturbances (*i.e.*, when  $\xi = 0$ ), a local linear approximation is given by<sup>12</sup>

$$y^L(\xi) \equiv \bar{y} + D_\xi y \cdot \xi.$$

But substituting  $y^L(\xi)$  into  $U^Q$  to obtain a quadratic function of  $\xi$  but does not in general result in an approximation to  $U$  that is accurate to second order; instead,

$$U(y(\xi); \xi) = U^Q(y^L(\xi); \xi) + \sum_j D_j U [\xi' D_{\xi\xi}^2 y^j \cdot \xi] + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \quad (1.12)$$

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<sup>12</sup>This linear approximation is the one that is given by solution of the linearized structural relations (1.6) using a similar linear approximation to the policy rule, in the case in which these linear relations have a determinate solution, again as a consequence of the implicit function theorem.

Here the second term on the right-hand side indicates omitted (policy-dependent) second-order terms in a correct approximation to  $U$  that result from second-order dependence of the equilibrium outcome  $y$  on the state  $\xi$ , omitted as a result of the linearization of  $y(\xi)$ . Such terms generally exist if the gradient of the welfare criterion  $D_y U$  is non-zero when evaluated at the unperturbed optimal policy.<sup>13</sup>

Both this problem and the incorrect outcome of LQ optimization can be solved, however, by using an *alternative* quadratic approximation to  $U$ . In order to obtain correct welfare rankings of alternative policies, it suffices that a quadratic function  $U^*(y; \xi)$  be such that

$$U(y; \xi) = U^*(y; \xi) + \mathcal{O}(\|\xi\|^3) \quad (1.13)$$

in the case of any  $y(\xi)$  satisfying (1.1) and such that  $\tilde{y} = \mathcal{O}(\|\xi\|)$ . If we find an alternative objective  $U^*$  that is also *purely quadratic*, in the sense of containing no linear terms ( $D_y U^* = 0$ ), then welfare can be evaluated to second order by  $\hat{U}(y^L(\xi); \xi)$ ;<sup>14</sup> and the error in LQ approximation of the optimal policy rule is eliminated as well. This is why the approach of Magill (1977) yields a correct LQ approximation.

In fact, an objective  $U^*$  of this form can quite generally be found. The key is to use a second-order Taylor series approximation to the constraints (1.1) to replace the linear terms in (1.9) with purely quadratic terms.<sup>15</sup> A second-order approximation to the structural relations (1.1), of the same form as the approximation (1.9), implies that

$$D_y F^k \cdot \tilde{y} = -\frac{1}{2} \tilde{y}' D_{yy}^2 F^k \cdot \tilde{y} - \tilde{y}' D_{y\xi}^2 F^k \cdot \xi + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3)$$

in the case of any  $(y; \xi)$  satisfying (1.1). The fact that  $\bar{y}$  is an optimal policy when the disturbances are zero implies that

$$D_y U = -\bar{\lambda}' D_y F, \quad (1.14)$$

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<sup>13</sup>See Woodford (2002; 2003, sec. 6.1) and Sutherland (2002) for further discussion.

<sup>14</sup>While this method (like the approach of simply computing a second-order approximation to  $y(\xi)$  and substituting this into objective) relies upon computing a second-order approximation to the model structural relations, the second-order approximation need be used only *once*, in determining the coefficients of the quadratic objective  $U^*$ , rather than having to be used again each time one seeks to evaluate the welfare associated with yet another candidate policy.

<sup>15</sup>A similar method is used by Sutherland (2002) to compute correct second-order approximations to welfare under alternative policies. However, his second-order approximation is computed for a particular parametric class of policies, while we derive a quadratic loss function that yields a correct welfare measure for *any* feasible policy.

where  $\bar{\lambda}$  is a vector of Lagrange multipliers associated with the constraints (1.1) in the case of zero disturbances. It then follows that

$$\begin{aligned} D_y U \cdot \tilde{y} &= - \sum_k \bar{\lambda}_k D_y F^k \cdot \tilde{y} \\ &= \frac{1}{2} \sum_k \bar{\lambda}_k \tilde{y}' D_{yy}^2 F^k \cdot \tilde{y} + \sum_k \bar{\lambda}_k \tilde{y}' D_{y\xi}^2 F^k \cdot \xi + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \end{aligned}$$

We can then use this expression to substitute for the term  $D_y U \cdot \tilde{y}$  in (1.9), yielding

$$\begin{aligned} U(y; \xi) &= \frac{1}{2} \tilde{y}' [D_{yy}^2 U \\ &\quad + \sum_k \bar{\lambda}_k D_{yy}^2 F^k] \cdot \tilde{y} + \tilde{y}' [D_{y\xi}^2 U + \sum_k \bar{\lambda}_k D_{y\xi}^2 F^k] \cdot \xi + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \end{aligned}$$

This is an approximation of the form (1.13), where

$$U^*(y; \xi) \equiv \frac{1}{2} \tilde{y}' [D_{yy}^2 U + \sum_k \bar{\lambda}_k D_{yy}^2 F^k] \cdot \tilde{y} + \tilde{y}' [D_{y\xi}^2 U + \sum_k \bar{\lambda}_k D_{y\xi}^2 F^k] \cdot \xi. \quad (1.15)$$

Use of the corrected quadratic objective (1.15) solves the problems associated with the use of  $U^Q$  discussed above. In particular, the LQ problem of maximizing (1.15) subject to the linearized constraints (1.6) satisfies linear first-order conditions of precisely the form (1.7). Hence this linear policy represents a correct linear approximation to the optimal policy  $y^{opt}(\xi)$ . There is a simple reason for this; in the case of any functions  $y(\xi), \lambda(\xi)$  such that  $\tilde{y}, \tilde{\lambda}$  are both of order  $\mathcal{O}(\|\xi\|)$ ,<sup>16</sup> a second-order Taylor expansion of the Lagrangian (1.4) takes the form

$$\mathcal{L}(y; \xi; \lambda) = U^*(y; \xi) + \tilde{\lambda}' [D_y F \cdot \tilde{y} + D_\xi F \cdot \xi] + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \quad (1.16)$$

But this is just the Lagrangian for the correct LQ problem, and the first-order conditions obtained by differentiating this approximate Lagrangian (which is the Lagrangian of the proposed approximate policy problem) agree, to first order, with those obtained by differentiating the exact Lagrangian.

The objective (1.15) can also be used to correctly rank alternative policies (none of which need be fully optimal), as long as these policies imply that  $y(0) = \bar{y}$ .<sup>17</sup>

<sup>16</sup>Again, the implicit function theorem implies that in the case that (1.8) is satisfied, the functions  $y(\xi), \lambda(\xi)$  that solve the exact Lagrangian problem have this property.

<sup>17</sup>Kim and Kim (2006) illustrate how the method expounded here can be used, for example, to correctly rank alternative policies with regard to international risk-sharing, in an example where naive LQ analysis sometimes gives an incorrect ranking.

One can easily verify that in the case of any feasible differentiable policy with this property,

$$U(y; \xi) = U^*(y^L(\xi); \xi) + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3),$$

where  $y^L(\xi)$  is the linear approximation to the policy in question. Hence using this criterion, welfare can be correctly evaluated to second order, using only a linear approximation to equilibrium outcomes under the policy in question; the problem resulting from “naive” linearization discussed by Kim and Kim (2003, 2006) is thus avoided.

## 1.2 Special Cases

While “naive” LQ optimization yields an incorrect linear approximation to optimal policy in general, as discussed above, it is an adequate approach under certain more restrictive conditions. This means that it is possible to use the simpler approach when it is used with sufficient care, as has often been the case in the literature. Our exposition above also makes clear in which cases a “naive” LQ approximation is possible. These are cases in which the additional terms present in (1.7) but not in (1.11) necessarily vanish.

One such case is when the constraints (1.1) are all exactly *linear*, in which case the second derivatives of the functions  $F^k$  vanish.<sup>18</sup> Sometimes it is possible to arrange for a problem to have constraints of this form, through some combination of restrictive specification of one’s model and careful choice of the variables in terms of which the problem is written, as in Kydland and Prescott (1982).<sup>19</sup> But while ingenuity can extend the range of applicability of naive LQ optimization, the class of models that can be put in this form is likely to be fairly restrictive.

Another such case is when the unperturbed optimum  $\bar{y}$  is also an *unconstrained* optimum in the case of zero disturbances, so that the multipliers  $\bar{\lambda}$  vanish, even though the constraints bind in general (and the associated Lagrange multipliers are

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<sup>18</sup>The importance of this condition for application of LQ approximation is stressed by Diaz-Jimenez (1999).

<sup>19</sup>Kydland and Prescott eliminate one nonlinear constraint by combining the production function with the utility function of the representative household, to obtain an objective written as a function of the paths of hours, capital and investment spending. The only remaining constraint is then a linear relation between investment spending and the dynamics of the capital stock; the linearity of this relation depends on their omission of an convex adjustment costs for the capital stock.

non-zero) in the presence of small disturbances  $\xi$ .<sup>20</sup> Again, sometimes it is possible to arrange for a problem to have this form, through some combination of restrictive model specification and an appropriate change of variables, as in Rotemberg and Woodford (1997).<sup>21</sup> But once again, the class of cases to which this result can be applied are likely to be quite restrictive. And in any event, as Judd (1996) stresses, it is undesirable for one's computational approach to yield correct answers only when the problem is expressed in terms of one set of variables rather than another. The approach described in section 1.1 eliminates the need for restrictions of the kind discussed in this section.

### 1.3 LQ Approximation in Models with Uncertainty

In the simple static example presented above, we have supposed that the value of the complete vector of disturbances  $\xi$  is known before any policy decisions must be made, and before any of the endogenous variables are determined; there is therefore no issue of policy choice under uncertainty. However, similar reasoning applies in the case of choice under uncertainty. Here we illustrate this through a reinterpretation of the analysis presented above.

Consider a problem in which there are two periods, and at least one dimension of policy must be decided in period 1, while at least one dimension of uncertainty about the exogenous disturbances is resolved only in period 2. Suppose that there are  $k$  possible states ( $s = 1, 2, \dots, k$ ) in period 2, each with some probability  $\pi(s)$  of occurring. The aim of policy is to maximize expected welfare

$$E[\hat{U}(y_1, y_2; \xi_1, \xi_2)] \equiv \sum_s \pi(s) \hat{U}(y_1, y_2(s); \xi_1, \xi_2(s)), \quad (1.17)$$

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<sup>20</sup>The importance of this condition for application of LQ approximation is stressed by Woodford (2002).

<sup>21</sup>Rotemberg and Woodford write their welfare objective in terms of the paths of sectoral output levels, rather than the consumption or hours worked by households, by substituting the production function and market-clearing condition into the objective. They then consider a model in which the steady-state path for output represents an optimal allocation of resources, because the stickiness of prices does not affect the equilibrium allocation of resources in the steady state. But the optimality of the steady-state allocation depends both on restriction of attention to policy rules consistent with zero steady-state inflation, and an assumption that subsidies offset the steady-state distortions that would otherwise result from firms' market power. The present method allows LQ approximation to be applied without these restrictive assumptions, as shown in Benigno and Woodford (2005a)

where  $\hat{U}$  is a smooth function of the  $n_1$  endogenous variables  $y_1$  that are determined in period 1, the  $n_2$  endogenous variables  $y_2$  that are determined in period 2, as well as the exogenous disturbances  $\xi_1$  that are realized in period 1 and the exogenous disturbances  $\xi_2$  that are realized in period 2. The possible equilibrium outcomes that can be achieved by policy are those consistent with the  $m_1$  structural equations

$$\mathbb{E}[\hat{f}_1(y_1, y_2; \xi_1, \xi_2)] = 0, \quad (1.18)$$

where each of the elements of  $\hat{f}_1$  is a smooth function of the same arguments as  $\hat{U}$ , and with the  $m_2$  structural equations

$$\hat{f}_2(y_1, y_2; \xi_1, \xi_2) = 0, \quad (1.19)$$

each of which must hold exactly, regardless of the state  $s$  that occurs in period 2.<sup>22</sup> Equations (1.18), together with policy decisions in period 1, determine the endogenous variables  $y_1$ ;<sup>23</sup> we suppose that  $m_1 < n_1$ , so that there is at least one dimension along which policy can vary in period 1. Equations (1.19), together with policy decisions in period 2 (if any) and the variables determined in period 1, then determine the endogenous variables  $y_2$  in whichever state  $s$  happens to be realized; we suppose that  $m_2 \leq n_2$ .

Under regularity conditions of the same kind as are needed in the static case, the method presented above can be used to derive a valid LQ approximation to this kind of policy problem as well. In fact, the calculations presented in section 1.1 are directly applicable. The right-hand side of (1.17) defines an objective  $U(y; \xi)$ , where the vectors  $y$  and  $\xi$  now specify all possible realizations of the random variables:

$$y \equiv \begin{bmatrix} y_1 \\ y_2(1) \\ \vdots \\ y_2(k) \end{bmatrix}, \quad \xi \equiv \begin{bmatrix} \xi_1 \\ \xi_2(1) \\ \vdots \\ \xi_2(k) \end{bmatrix}.$$

The left-hand side of (1.18) can similarly be written as a vector of  $m_1$  functions of the vectors  $y$  and  $\xi$ , while for each possible state  $s$ , the left-hand side of (1.19) is a vector

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<sup>22</sup>On the left-hand side of (1.17) and in equation (1.18),  $y_2$  and  $\xi_2$  are random variables, whereas in (1.19) these symbols refer to the values of those variables that are realized in period 2.

<sup>23</sup>Note that these structural relations need not all involve expectations; we allow for the case in which some elements of  $\hat{f}_1$  may not depend on either  $y_2$  or  $\xi_2$ . The important feature of the relations (1.18) is that they involve only information available in period 1.

of  $m_2$  functions of the vectors  $y$  and  $\xi$ .<sup>24</sup> Hence the complete system of structural relations, for both periods and for all possible states in period 2, can be written as a system of the form (1.1), where now  $F$  is a vector of  $m = m_1 + k \cdot m_2$  functions of  $y$  and  $\xi$ . (Under the assumptions made in the previous paragraph,  $m < n$ , as assumed in section 1.1, where  $n = n_1 + k \cdot n_2$  is the length of the vector  $y$ .)

If  $\hat{U}$ ,  $\hat{f}_1$ , and  $\hat{f}_2$  are continuously differentiable functions (of whatever order) of  $\hat{y}' \equiv (y_1, y_2)'$  and  $\hat{\xi}' \equiv (\xi_1, \xi_2)'$ , then it follows that  $U$  and  $F$  will be correspondingly differentiable functions of  $y$  and  $\xi$ . The results of section 1.1 are then directly applicable. In the case that the implicit function theorem can be applied to derive a local linear approximation to optimal policy,<sup>25</sup> that correct linear approximation corresponds to the linear policy that solves an LQ optimization problem.

In the case that the uncertainty is small — *i.e.*, in each state  $s$ ,  $\xi_2(s)$  is close to the *same* value  $\bar{\xi}_2$  (which we may denote as zero, without loss of generality) — there is additional structure that we can exploit in writing the approximate LQ problem. Let us suppose not only that this is true (*i.e.*, that the unperturbed problem corresponds to  $\xi = 0$  in all elements), but also that the solution to the unperturbed problem is

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<sup>24</sup>Of course, these latter functions depend on only a subset of the elements of  $y$  and  $\xi$ , namely, those corresponding to  $y_1, y_2(s), \xi_1$ , and  $\xi_2(s)$ , where  $s$  is the particular state for which the structural relations are written.

<sup>25</sup>Here by “linear approximation” we mean that  $y^{opt}(\xi)$  is approximated by a linear function of  $\xi$ , which differs from the exact function by a residual that is at most of order  $\mathcal{O}(\|\xi\|^2)$ , as in (1.3). Note that this is a different sense than the one proposed by Judd (1996) for stochastic models. Judd considers a perturbation of a dynamic optimization problem under certainty by varying a factor that scales the amplitude of a mean-zero random disturbance to fundamentals; the disturbance is multiplied by  $\sqrt{\epsilon}$ , so that the variance of the disturbance is proportional to  $\epsilon$ . Judd considers a “linear approximation” to be a perturbation of the (deterministic) solution for the  $\epsilon = 0$  case that includes all terms linear in  $\epsilon$ , including any non-zero derivatives of the average values of endogenous variables with respect to the variance of the disturbance; as a consequence, Judd refers to the approximation obtained by linearizing the first-order conditions for optimality as “at most a half-linear approximation” (1996, p. 538). In the present approach, if we write  $\xi_2(s) = \sqrt{\epsilon}\bar{\xi}_2(s)$ , where the random variable  $\bar{\xi}_2$  remains fixed as we vary  $\sqrt{\epsilon}$ , then only terms of order  $\mathcal{O}(\sqrt{\epsilon})$  are considered to be of “first order” (*i.e.*, of order  $\mathcal{O}(\|\xi\|)$ ); terms linear in  $\epsilon$ , such as the terms indicating how the average values of variables vary linearly with the variance of the disturbances, are treated as part of the residual of order  $\mathcal{O}(\|\xi\|^2)$ . Of course, this is in no way intended to deny that it may be of interest to calculate such effects; however, in the case of small enough random disturbances (*i.e.*, a small enough value of  $\sqrt{\epsilon}$ ), these effects should be small relative to the ones taken account of in the linear approximation derived here.

deterministic (*i.e.*,  $y_2(s) = \bar{y}_2$  for all  $s$  as well).<sup>26</sup> In this case, the solution  $\bar{y}$  to the unperturbed problem must correspond to the solution  $(\bar{y}_1, \bar{y}_2)$  to the deterministic problem

$$\max_{\hat{y}} \hat{U}(\hat{y}; 0) \quad \text{s.t.} \quad \hat{F}(\hat{y}; 0) = 0,$$

where  $\hat{F}$  is the vector of  $n_1 + n_2$  functions  $\hat{f}_1$  and  $\hat{f}_2$ . This latter solution must satisfy first-order conditions of the form

$$D_{\hat{y}} \hat{U} = -\hat{\lambda}' D_{\hat{y}} \hat{F}, \quad (1.20)$$

where all derivatives are evaluated at  $(\bar{y}_1, \bar{y}_2; 0, 0)$ . It follows that the solution to the unperturbed problem (in which  $y_2$  is however allowed to be state-dependent) satisfies first-order conditions of the form (1.14), in which the vector of Lagrange multipliers is given by

$$\bar{\lambda}' = [\hat{\lambda}'_1, \pi(1)\hat{\lambda}'_2, \dots, \pi(k)\hat{\lambda}'_2].$$

Because the steady-state vector of Lagrange multipliers takes this form (and the functions  $U$  and  $F$  are additively separable across states), the Lagrangean for the stochastic policy problem can be written in the form

$$\mathcal{L} = \text{E}[\hat{U}(\hat{y}; \hat{\xi}) + \hat{\lambda}' \hat{F}(\hat{y}; \hat{\xi})],$$

in the case in which the multipliers  $\lambda$  are set equal to  $\bar{\lambda}$ . It follows that the correct quadratic objective (1.15) for the LQ approximation is of the form

$$U^*(y; \xi) = \text{E}[\hat{U}^*(\hat{y}; \hat{\xi})],$$

where

$$\hat{U}^*(\hat{y}; \hat{\xi}) \equiv \frac{1}{2} \tilde{y}' [D_{\hat{y}\hat{y}}^2 \hat{U} + \sum_k \hat{\lambda}_k D_{\hat{y}\hat{y}}^2 \hat{F}^k] \cdot \tilde{y} + \tilde{y}' [D_{\hat{y}\hat{\xi}}^2 \hat{U} + \sum_k \hat{\lambda}_k D_{\hat{y}\hat{\xi}}^2 \hat{F}^k] \cdot \hat{\xi}, \quad (1.21)$$

and  $\tilde{y}$  denotes the difference between  $\hat{y}$  and the unperturbed optimal values. Note that (1.21) is just the quadratic objective for the LQ approximation to the deterministic

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<sup>26</sup>This latter property necessarily follows from the assumption that  $\xi$  is deterministic in the case of a strictly convex problem, but the implication need not follow for all problems of the more general sort that we consider here, so an additional assumption is required. In section 3, we show how to check whether the deterministic optimal steady state is indeed at least a local optimum of the stochastic problem as well.



problem of maximizing  $\hat{U}(\hat{y}; \hat{\xi})$  subject to the constraints  $\hat{F}(\hat{y}; \hat{\xi}) = 0$ , derived using the method of section 1.1.

So the correct quadratic objective for the LQ approximation of the stochastic problem is just the expected value of the quadratic objective for the corresponding deterministic problem. Similarly, the local linear approximation to the constraints (1.18) is of the form

$$\Phi_1 \tilde{y}_1 + \Phi_2 \mathbb{E}[\tilde{y}_2] + \Psi_1 \xi_1 + \Psi_2 \mathbb{E}[\xi_2] = 0, \quad (1.22)$$

where the matrices of coefficients are the same as in the deterministic model, while the local linear approximation to constraints (1.19) is of exactly the same form as in the deterministic model. Thus all coefficients of both the quadratic objective and the linear constraints are the same as in the LQ approximation to the deterministic problem (and represent partial derivatives of the functions  $\hat{U}$  and  $\hat{F}$ , evaluated at the optimum of the deterministic problem when  $\hat{\xi} = 0$ ); the only difference in the form of the two LQ problems is the fact that expected values are taken in (1.21) and (1.22).

As is well known, the solution to a stochastic LQ problem of this kind exhibits the property of *certainty equivalence*. In particular, the linear approximation to  $y^{opt}(\xi)$  is of the form

$$\begin{aligned} y_1 &= M\xi_1 + N\mathbb{E}[\xi_2], \\ y_2 &= Py_1 + Q\xi_1 + R\xi_2, \end{aligned}$$

where the matrices of coefficients are the same as in the linear approximation to  $\hat{y}^{opt}(\hat{\xi})$  in the corresponding deterministic (perfect foresight) problem,

$$\begin{aligned} y_1 &= M\xi_1 + N\xi_2, \\ y_2 &= Py_1 + Q\xi_1 + R\xi_2. \end{aligned}$$

Of course, this does not mean that the *exact* solution for optimal policy in the presence of uncertainty generally possesses the property of certainty equivalence. However, departures from certainty equivalence (for example, effects of a mean-preserving change in the variance of  $\xi_2$  on the optimal choice of  $y_1$ ) represent contributions to  $y^{opt}(\xi)$  that are at most of order  $\mathcal{O}(\|\xi\|^2)$ .

## 1.4 Qualifications

While the conditions under which a valid LQ approximation is possible are fairly general, several qualifications are in order. First of all, the LQ approximation, when valid, is purely *local* in character; it can only provide an approximate characterization of optimal policy to the extent that disturbances are sufficiently small. Whether the disturbances are small enough for this to be a useful approximation will depend upon the application; and a judgment about how accurate the approximation is likely to be is not possible on the basis of the coefficients of the LQ approximate problem alone. And like all perturbation approaches, it depends on sufficient differentiability of the problem;<sup>27</sup> it cannot be applied, for example, to problems in which there are inequality constraints that sometimes bind but at other times do not. Moreover, the LQ approximation provides at best a linear approximation to optimal policy. More general perturbation methods<sup>28</sup> can instead be used to compute approximations of any desired order, assuming sufficient differentiability of the objective and constraints. In this respect, LQ approximation is hardly a substitute for an understanding of general perturbation calculations, as stressed by Judd (1996).

Second, a correct LQ approximation yields a correct linear approximation to the optimal policy in the case that linearization of the first-order conditions (1.5) would also yield a system of linear equations that can be solved to obtain a linear approximation to optimal policy. If the regularity condition (1.8) fails, the implicit function theorem cannot be applied to obtain a linear approximation in this way, and the LQ approach similarly fails to provide a correct linear approximation to optimal policy. This is a problem that can certainly arise in cases of economic interest, such as the portfolio problem treated by Judd and Guu (2001), and more complex perturbation methods (used to deal with singular perturbations) can still be employed in

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<sup>27</sup>Whether the objective and constraints are sufficiently differentiable in a given application may depend on the choice of variables in terms of which one writes these functions. Kim *et al.* (2007) provide an example of an optimal stabilization policy problem in which the first-order conditions describing optimal policy cannot be linearized — so that the LQ methodology expounded here would also not be applicable — when the state variables are assumed to include the square root of a measure of price dispersion, rather than the measure of price dispersion itself. Linearization is instead possible under the alternative choice of variables.

<sup>28</sup>A generally useful approach to obtaining a higher-order Taylor series approximation to  $y^{opt}(\xi)$  is to solve a higher-order Taylor series approximation to the first-order conditions (1.5), using the approach explained by Judd (1996, 1998).

such cases, as Judd and Guu show. But it is a problem the existence of which can be diagnosed within the LQ analysis itself: for when the condition (1.8) fails, the first-order conditions of the LQ problem fail to determine a unique solution. Thus it remains true that *if* the LQ analysis determines a unique linear policy, this will be a correct linear approximation to optimal policy. But when the LQ analysis implies that optimal policy is indeterminate, this need not be a correct conclusion; there may instead be a unique optimal policy, a correct linear approximation to which depends on higher-order derivatives than those considered in the LQ approximation. An identical caveat applies to the method of linearization of the first-order conditions characterizing optimal policy.

Third, a correct LQ approximation yields a correct linear approximation to the optimal policy only in the case that the perturbed solution to the first-order conditions (1.5) characterized by the implicit function theorem is in fact an optimum. It cannot be taken as obvious that the first-order conditions suffice for optimality, since in applications of interest, the structural relations (1.1) often define a non-convex set. The question of convexity can be addressed at least locally by evaluating the relevant second-order conditions corresponding to a given solution to the first-order conditions. This is straightforward within the LQ analysis itself: one simply needs to verify the concavity of the quadratic objective  $U^*$  in  $\tilde{y}$ , for vectors  $\tilde{y}$  in the linear subspace such that  $D_y F \cdot \tilde{y} = 0$ . This is an algebraic property of the coefficients of the LQ problem, involving the signs of certain principal minors of the matrix appearing in (1.8), as shown by Debreu (1952).<sup>29</sup> But of course, verification of the second-order conditions for optimality still only guarantees that the solution to the LQ problem approximates a *local* welfare optimum. The question of global optimality of the solution cannot be treated using purely local methods, and is often quite difficult in dynamic stochastic models.

We turn now to the additional complications that arise in applying this method to dynamic, stochastic policy problems. Foremost among these complications are ones that result from the presence of forward-looking constraints, indicating the way in which equilibrium determination is affected by forward-looking optimizing decisions on the part of the public.

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<sup>29</sup>We discuss the generalization of this characterization to the dynamic case in section 3.

## 2 LQ Approximation of a Problem with Forward-Looking Constraints

We wish to consider an abstract discrete-time dynamic optimal policy problem of the following sort.<sup>30</sup> Suppose that the policy authority wishes to determine the evolution of an (endogenous) state vector  $\{y_t\}$  for  $t \geq t_0$  to maximize an objective of the form

$$V_{t_0} \equiv E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \pi(y_t, \xi_t), \quad (2.1)$$

where  $0 < \beta < 1$  is a discount factor, the period objective  $\pi(y, \xi)$  is a concave function of  $y$ , and  $\xi_t$  is a vector of exogenous disturbances. The evolution of the endogenous states must satisfy a system of backward-looking structural relations

$$F(y_t, \xi_t; y_{t-1}) = 0 \quad (2.2)$$

and a system of forward-looking structural relations

$$E_t g(y_t, \xi_t; y_{t+1}) = 0, \quad (2.3)$$

that both must hold for each  $t \geq t_0$ , given the vector of initial conditions  $y_{t_0-1}$ .

Conditions of the form (2.2) allow current endogenous variables to depend on lagged states; for example, these relations could include a technological relation between the capital stock carried into the next period, current investment expenditure, and the capital stock carried into the current period.<sup>31</sup> Conditions of the form (2.3) instead allow current endogenous variables to depend on current expectations regarding future states; for example, these relations could include an Euler equation for the optimal timing of consumer expenditure, relating current consumption to expected consumption in the next period and the expected rate of return on saving.<sup>32</sup> While the most general notation would allow both leads and lags in all of the structural

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<sup>30</sup>Applications of the framework proposed here are discussed in section 5.

<sup>31</sup>The next period's capital stock and the current investment expenditure would both be elements of  $y_t$ ; the vector  $\xi_t$  could include a random disturbance to investment adjustment costs.

<sup>32</sup>Current consumption and the current period ex-post return on saving in the previous period would both be elements of  $y_t$ ; the vector  $\xi_t$  could include a random disturbance to the impatience to consume. Note that without loss of generality we may suppose that the vector  $\xi_t$  includes all information available in period  $t$  regarding future exogenous disturbances.

equations, supposing that there are equations of these two types will make clearer the different types of complications arising from the two distinct types of intertemporal linkages. We shall suppose that the number  $n_F$  of constraints of the first type each period plus the number  $n_g$  of constraints of the second type is less than the number  $n_y$  of endogenous state variables each period, so that there is at least one dimension along which policy can continuously vary the outcome  $y_t$  each period, given the past and expected future evolution of the endogenous variables. A  $t_0$ -optimal commitment (the standard Ramsey policy problem) is then the state-contingent evolution  $\{y_t\}$  consistent with equations (2.2)–(2.3) for all  $t \geq t_0$  that maximizes (2.1).

## 2.1 A Recursive Policy Problem

As is well-known, the presence of the forward-looking constraints (2.3) implies that a  $t_0$ -optimal commitment is not generally time-consistent. If, however, we suppose that a policy to apply from period  $t_0$  onward must be chosen subject to an additional set of constraints on the acceptable values of  $y_{t_0}$ , it is possible for the resulting policy problem to have a recursive structure.<sup>33</sup> While this is not necessary for the method of LQ approximation to be applicable,<sup>34</sup> it is necessary in order for both our approximate quadratic objective and approximate linear constraints to involve coefficients that are time-invariant, and correspondingly for our derived linear approximation to optimal policy to involve time-invariant coefficients. In the case of the unconstrained (Ramsey) optimal policy problem, the  $t_0$ -optimal policy generally does not imply constant values of the endogenous variables, even when there are no random disturbances and the functions  $\pi$ ,  $F$  and  $g$  are all time-invariant, as assumed above; correspond-

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<sup>33</sup>The fact that a recursive structure can be restored, allowing dynamic-programming methods to be employed, through a suitable modification of the assumed objective and/or constraints has been known since the seminal work of Kydland and Prescott (1980). Marcet and Marimon (1998) provide a detailed analysis of an approach that modifies the policy objective by adding additional multiplier terms; the additional terms in the objective of the modified problem of Marcet and Marimon lead to the same additional terms in the Lagrangian for the policy problem as the additional constraints that we introduce here. We prefer to introduce initial pre-commitments because of the more transparent connection of the modified problem to the original policy problem under this exposition. The first-order conditions for optimal policy in the recursive policy problem that we propose are the same as those derived by Marcet and Marimon, except in the initial period.

<sup>34</sup>This is illustrated by the treatment of a simple example with forward-looking constraints in section 1.3 above.

ingly, a local approximation to Ramsey policy in the case of small disturbances must involve derivatives evaluated along this non-constant path, so that the coefficients of the linear approximation are generally time-varying. The case considered here is clearly more convenient computationally, and it is arguable that the solution to this kind of problem represents a more appealing policy commitment as well.<sup>35</sup>

As discussed in Benigno and Woodford (2003, 2005a), in order to obtain a problem with a recursive structure (the solution to which can be described by a time-invariant policy rule), we must choose initial pre-commitments regarding  $y_{t_0}$  that are *self-consistent*, in the sense that the policy that is chosen subject to these constraints would also satisfy constraints of exactly the same form in all later periods as well. The required initial pre-commitments are of the form

$$g(y_{t_0-1}, \xi_{t_0-1}; y_{t_0}) = \bar{g}_{t_0}, \quad (2.4)$$

where  $\bar{g}_{t_0}$  may depend on the exogenous state at date  $t_0$ . Note that we assume the existence of a pre-commitment only about those aspects of  $y_{t_0}$  the anticipation of which back in period  $t_0 - 1$  should have been relevant to equilibrium determination then; there is no need for any stronger form of commitment in order to render optimal policy time-consistent.

We are thus interested in characterizing the state-contingent policy  $\{y_t\}$  for  $t \geq t_0$  that maximizes (2.1) subject to constraints (2.2) – (2.4). Such a policy is *optimal from a timeless perspective* if  $\bar{g}_{t_0}$  is chosen, as a function of predetermined or exogenous states at  $t_0$ , according to a self-consistent rule.<sup>36</sup> This means that the initial pre-commitment is determined by past conditions through a function

$$\bar{g}_{t_0} = \bar{g}(\xi_{t_0}, \mathbf{y}_{\mathbf{t}_0-1}), \quad (2.5)$$

where  $\mathbf{y}_{\mathbf{t}}$  is an extended state vector;<sup>37</sup> this function has the property that under optimal policy, given this initial pre-commitment, the state-contingent evolution of the economy will satisfy

$$g(y_{t-1}, \xi_{t-1}; y_t) = \bar{g}(\xi_t, \mathbf{y}_{\mathbf{t}-1}) \quad (2.6)$$

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<sup>35</sup>See Giannoni and Woodford (2002) and Woodford (2003, chap. 7) for further discussion.

<sup>36</sup>See Giannoni and Woodford (2002), Woodford (2003, chap. 7), or Benigno and Woodford (2005a) for further discussion.

<sup>37</sup>The extended state vector may include both endogenous and exogenous variables, the values of which are realized in period  $t$  or earlier. More specific assumptions about the nature of the extended state vector are made below; see the discussion of equation (2.8).

in each possible state of the world at each date  $t \geq t_0$  as well. Thus the initial constraint is of a form that one would optimally commit oneself to satisfy at all (subsequent) dates.

Let  $V(\bar{g}_{t_0}; y_{t_0-1}, \xi_{t_0}, \xi_{t_0-1})$  be the maximum achievable value of the objective (2.1) in this problem.<sup>38</sup> Then the infinite-horizon problem just defined is equivalent to a sequence of one-period decision problems in which, in each period  $t \geq t_0$ , a value of  $y_t$  is chosen and state-contingent one-period-ahead pre-commitments  $\bar{g}_{t+1}(\xi_{t+1})$  (for each of the possible states  $\xi_{t+1}$  in the following period) are chosen so as to maximize

$$\pi(y_t, \xi_t) + \beta E_t V(\bar{g}_{t+1}; y_t, \xi_{t+1}, \xi_t), \quad (2.7)$$

subject to the constraints

$$F(y_t, \xi_t; y_{t-1}) = 0,$$

$$g(y_{t-1}, \xi_{t-1}; y_t) = \bar{g}_t,$$

$$E_t \bar{g}_{t+1} = 0,$$

given the values of  $\bar{g}_t, y_{t-1}, \xi_{t-1}$ , and  $\xi_t$ , all of which are predetermined and/or exogenous in period  $t$ . It is this recursive policy problem that we wish to study; note that it is only when we consider this problem (as opposed to the unconstrained Ramsey problem) that it is possible, in general, to obtain a deterministic steady state as an optimum in the case of suitable initial conditions, and hence only in this case that we can hope to approximate the optimal policy problem around such a steady state.<sup>39</sup>

The solution to the recursive policy problem just defined involves values for the endogenous variables  $y_t$  given by a policy function of the form

$$y_t = y^*(\bar{g}_t, y_{t-1}, \xi_t, \xi_{t-1}),$$

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<sup>38</sup>We assume, to economize on notation, that the exogenous state vector  $\xi_t$  evolves in accordance with a Markov process. Hence  $\xi_t$  summarizes not only all of the disturbances that affect the structural relations at date  $t$ , but all information at date  $t$  about the subsequent evolution of the exogenous disturbances. This is important in order for a time-invariant value function to exist with the arguments indicated.

<sup>39</sup>In the literature on Ramsey policy, one sometimes sees approximate characterizations of optimal policy computed by log-linearizing around a steady state that Ramsey policy approaches asymptotically in the absence of random disturbances. But in such a case, there is no guarantee that the approximate characterization would be accurate even in the case of arbitrarily small disturbances, as Ramsey policy need not be near the steady state except asymptotically.

and a choice of the following period's pre-commitment  $\bar{g}_{t+1}$  of the form

$$\bar{g}_{t+1} = g^*(\xi_{t+1}; \bar{g}_t, y_{t-1}, \xi_t, \xi_{t-1}),$$

where  $y^*$  and  $g^*$  are time-invariant functions. Let us suppose furthermore that the evolution of the extended state vector depends only on the evolution of the two vectors  $\{y_t, \xi_t\}$ , through a recursion of the form

$$\mathbf{y}_t = \psi(\xi_t, y_t, \mathbf{y}_{t-1}); \quad (2.8)$$

this system of identities *defines* the extended state vector, the elements of which consist essentially of linear combinations of current and lagged elements of the vectors  $y_t$  and  $\xi_t$ . (To simplify notation, we shall suppose that the current values  $y_t$  and  $\xi_t$  are among the elements of  $\mathbf{y}_t$ .) The initial pre-commitment (2.5) is then self-consistent if

$$g^*(\xi_{t+1}; \bar{g}(\xi_t, \mathbf{y}_{t-1}), y_{t-1}, \xi_t, \xi_{t-1}) = \bar{g}(\xi_{t+1}, \psi(\xi_t, y^*(\bar{g}_t, y_{t-1}, \xi_t, \xi_{t-1}), \mathbf{y}_{t-1})) \quad (2.9)$$

for all possible values of  $\xi_{t+1}, \xi_t$ , and  $\mathbf{y}_{t-1}$ .<sup>40</sup> Note that this implies that equation (2.6) is satisfied at all times.

## 2.2 A Correct LQ Local Approximation

We now derive a corresponding LQ problem using local approximations to both the objective and the constraints of the above problem. In order for these local approximations to involve coefficients that remain the same over time, we compute them near the special case of an optimal policy that involves values of the state variables that are constant over time. This special case involves both zero disturbances and suitably chosen initial conditions; we then seek to approximately characterize optimal policy for nearby problems in which the disturbances are small and the initial conditions are *close* to satisfying the assumed special conditions. To be precise, we assume both an initial state  $y_{t_0-1}$  and initial pre-commitments  $\bar{g}_{t_0}$  such that the optimal policy in the case of zero disturbances is a steady state, *i.e.*, such that  $y_t = \bar{y}$  for all  $t$ , for some vector  $\bar{y}$ . (Our subsequent calculations then assume that both  $y_{t_0-1}$  and  $\bar{g}_{t_0-1}$  are close enough to being consistent with this steady state.) In order to define the

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<sup>40</sup>Both sides of this equation involve only the elements of  $\xi_{t+1}, \xi_t$ , and  $\mathbf{y}_{t-1}$ , on the understanding that  $y_{t-1}$  and  $\xi_{t-1}$  are both elements of  $\mathbf{y}_{t-1}$ .



steady state, we must consider the nature of optimal policy in the exact problem just defined.

The first-order conditions for the exact policy problem can be obtained by differentiating a Lagrangian of the form

$$\mathcal{L}_{t_0} = V_{t_0} + E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\lambda'_t F(y_t, \xi_t; y_{t-1}) + \beta^{-1} \varphi'_{t-1} g(y_{t-1}, \xi_{t-1}; y_t)], \quad (2.10)$$

where  $\lambda_t$  and  $\varphi_t$  are Lagrange multipliers associated with constraints (2.2) and (2.3) respectively, for any date  $t \geq t_0$ , and we use the notation  $\beta^{-1} \varphi_{t_0-1}$  for the Lagrange multiplier associated with the additional constraint (2.4). This last notational choice allows the first-order conditions to be expressed in the same way for all periods. Optimality requires that the joint evolution of the processes  $\{y_t, \xi_t, \lambda_t, \varphi_t\}$  satisfy

$$\begin{aligned} D_y \pi(y_t, \xi_t) + \lambda'_t D_y F(y_t, \xi_t; y_{t-1}) + \beta E_t \lambda_{t+1}' D_{\bar{y}} F(y_{t+1}, \xi_{t+1}; y_t) \\ + E_t \varphi'_t D_y g(y_t, \xi_t; y_{t+1}) + \beta^{-1} \varphi'_{t-1} D_{\bar{y}} g(y_{t-1}, \xi_{t-1}; y_t) = 0 \end{aligned} \quad (2.11)$$

at each date  $t \geq t_0$ , where  $D_y$  denotes the vector of partial derivatives of any of the functions with respect to the elements of  $y_t$ , while  $D_{\bar{y}}$  means the vector of partial derivatives with respect to the elements of  $y_{t+1}$  and  $D_{\bar{y}}$  means the vector of partial derivatives with respect to the elements of  $y_{t-1}$ .

An *optimal steady state* is then described by a collection of vectors  $(\bar{y}, \bar{\lambda}, \bar{\varphi})$  satisfying

$$\begin{aligned} D_y \pi(\bar{y}, 0) + \bar{\lambda}' D_y F(\bar{y}, 0; \bar{y}) + \beta \bar{\lambda}' D_{\bar{y}} F(\bar{y}, 0; \bar{y}) \\ + \bar{\varphi}' D_y g(\bar{y}, 0; \bar{y}) + \beta^{-1} \bar{\varphi}' D_{\bar{y}} g(\bar{y}, 0; \bar{y}) = 0, \end{aligned} \quad (2.12)$$

$$F(\bar{y}, 0; \bar{y}) = 0, \quad (2.13)$$

$$g(\bar{y}, 0; \bar{y}) = 0. \quad (2.14)$$

We shall suppose that such a steady state exists, and assume (in the policy problem with random disturbances) an initial state  $y_{t_0-1}$  near  $\bar{y}$  — more precisely, such that  $y_{t_0-1} - \bar{y} = \mathcal{O}(\|\xi\|)$  — and an initial pre-commitment such that  $\bar{g}_{t_0} = \mathcal{O}(\|\xi\|)$  as well.<sup>41</sup> Once the optimal steady state has been computed, we make no further use of conditions (2.11); our proposed method does not require that we directly seek to solve these equations.

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<sup>41</sup>Note that the steady-state value of  $\bar{g}$  is equal to  $g(\bar{y}, 0; \bar{y}) = 0$ .

Instead, we now consider local approximations to the objective and constraints near an optimal steady state. We can compute a second-order Taylor expansion of the period objective function  $\pi$ , obtaining an expression of exactly the form (1.9). Substituting this into (2.1), we obtain the approximate objective

$$V_{t_0} = E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ D_y \pi \cdot \tilde{y}_t + \frac{1}{2} \tilde{y}'_t D_{yy}^2 \pi \cdot \tilde{y}_t + \tilde{y}'_t D_{y\xi}^2 \pi \cdot \xi_t \right] + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \quad (2.15)$$

This would be used as the quadratic objective in what we have called the “naive” LQ approximation.

However, (2.15) is not the only valid quadratic approximation to (2.1). Taylor’s theorem implies that it is the only quadratic function that correctly approximates (2.1) in the case of *arbitrary* (small enough) variations in the state variables, but there are others that will also correctly approximate (2.1) in the case of variations that are *consistent with the structural relations*. We can obtain an infinite number of alternative quadratic welfare measures by adding to (2.15) arbitrary multiples of quadratic (Taylor series) approximations to functions that must equal zero in order for the structural relations to be satisfied. Among these, we are able to find a welfare measure that is *purely quadratic*, i.e., that contains no non-zero linear terms, as in Benigno and Woodford (2003), so that a linear approximation to the equilibrium evolution of the endogenous variables under a given policy rule suffices to allow the welfare measure to be evaluated to second order. The key to this is using a second-order approximation to the structural relations to substitute purely quadratic terms for the linear terms  $D_y \pi \cdot \tilde{y}_t$  in the sum (2.15), as in Sutherland (2002).

A similar second-order Taylor series approximation can be written for each of the

functions  $F^k$ . It follows that

$$\begin{aligned}
\sum_{t=t_0}^{\infty} \beta^{t-t_0} \bar{\lambda}' F(y_t, \xi_t; y_{t-1}) &= \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left\{ \bar{\lambda}' [D_y F \cdot \tilde{y}_t + D_{\tilde{y}} F \cdot \tilde{y}_{t-1}] \right. \\
&\quad + \bar{\lambda}_k \left[ \frac{1}{2} \tilde{y}'_t D_{yy}^2 F^k \cdot \tilde{y}_t + \tilde{y}'_t D_{y\xi}^2 F^k \cdot \xi_t + \tilde{y}'_{t-1} D_{\tilde{y}\xi}^2 F^k \cdot \xi_t \right. \\
&\quad \left. \left. + \frac{1}{2} \tilde{y}'_{t-1} D_{\tilde{y}\tilde{y}}^2 F^k \cdot \tilde{y}_{t-1} + \tilde{y}'_t D_{y\tilde{y}}^2 F^k \cdot \tilde{y}_{t-1} \right] \right\} \\
&\quad + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3) \\
&= \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left\{ \bar{\lambda}' [D_y F + \beta D_{\tilde{y}} F] \cdot \tilde{y}_t \right. \\
&\quad + \frac{1}{2} \bar{\lambda}_k \left[ \tilde{y}'_t D_{yy}^2 F^k \cdot \tilde{y}_t + 2\tilde{y}'_t D_{y\xi}^2 F^k \cdot \xi_t + 2\beta \tilde{y}'_t D_{\tilde{y}\xi}^2 F^k \cdot \xi_{t+1} \right. \\
&\quad \left. \left. + \beta \tilde{y}'_t D_{\tilde{y}\tilde{y}}^2 F^k \cdot \tilde{y}_t + 2\tilde{y}'_t D_{y\tilde{y}}^2 F^k \cdot \tilde{y}_{t-1} \right] \right\} \\
&\quad + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \tag{2.16}
\end{aligned}$$

Using a similar Taylor series approximation of each of the functions  $g^i$ , we correspondingly obtain

$$\begin{aligned}
\sum_{t=t_0}^{\infty} \beta^{t-t_0-1} \bar{\varphi}' g(y_{t-1}, \xi_{t-1}; y_t) &= \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left\{ \bar{\varphi}' [D_y g + \beta^{-1} D_{\tilde{y}} g] \cdot \tilde{y}_t \right. \\
&\quad + \frac{1}{2} \bar{\varphi}_i \left[ \tilde{y}'_t D_{yy}^2 g^i \cdot \tilde{y}_t + 2\tilde{y}'_t D_{y\xi}^2 g^i \cdot \xi_t + 2\beta^{-1} \tilde{y}'_t D_{\tilde{y}\xi}^2 g^i \cdot \xi_{t-1} \right. \\
&\quad \left. \left. + \beta^{-1} \tilde{y}'_t D_{\tilde{y}\tilde{y}}^2 g^i \cdot \tilde{y}_t + 2\beta^{-1} \tilde{y}'_t D_{y\tilde{y}}^2 g^i \cdot \tilde{y}_{t-1} \right] \right\} \\
&\quad + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \tag{2.17}
\end{aligned}$$

It then follows from constraints (2.2)–(2.4) that in the case of any admissible policy,<sup>42</sup>

$$\begin{aligned}
\beta^{-1}\bar{\varphi}'\bar{g}_{t_0} &= E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\bar{\lambda}'F(y_t, \xi_t; y_{t-1}) + \beta^{-1}\bar{\varphi}'g(y_{t-1}, \xi_{t-1}; y_t)] \\
&= E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left\{ [\bar{\lambda}'(D_y F + \beta D_{\tilde{y}} F) + \bar{\varphi}'(D_y g + \beta^{-1} D_{\tilde{y}} g)] \cdot \tilde{y}_t \right. \\
&\quad + \frac{1}{2} \bar{\lambda}_k [\tilde{y}'_t D_{yy}^2 F^k \cdot \tilde{y}_t + 2\tilde{y}'_t D_{y\xi}^2 F^k \cdot \xi_t + 2\beta \tilde{y}'_t D_{\tilde{y}\xi}^2 F^k \cdot \xi_{t+1} \\
&\quad + \beta \tilde{y}'_t D_{\tilde{y}y}^2 F^k \cdot \tilde{y}_t + 2\tilde{y}'_t D_{\tilde{y}\tilde{y}}^2 F^k \cdot \tilde{y}_{t-1}] \\
&\quad + \frac{1}{2} \bar{\varphi}_i [\tilde{y}'_t D_{yy}^2 g^i \cdot \tilde{y}_t + 2\tilde{y}'_t D_{y\xi}^2 g^i \cdot \xi_t + 2\beta^{-1} \tilde{y}'_t D_{\tilde{y}\xi}^2 g^i \cdot \xi_{t-1} \\
&\quad + \beta^{-1} \tilde{y}'_t D_{\tilde{y}y}^2 g^i \cdot \tilde{y}_t + 2\beta^{-1} \tilde{y}'_t D_{\tilde{y}\tilde{y}}^2 g^i \cdot \tilde{y}_{t-1}] \left. \right\} \\
&\quad + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3), \tag{2.18}
\end{aligned}$$

where we have used (2.16) and (2.17) to substitute for the  $F$  and  $g$  terms respectively.

We can write this more compactly in the form

$$\begin{aligned}
\beta^{-1}\bar{\varphi}'\bar{g}_{t_0} &= E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left\{ \Phi \cdot \tilde{y}_t + \frac{1}{2} [\tilde{y}'_t H \cdot \tilde{y}_t + 2\tilde{y}'_t R \tilde{y}_{t-1} + 2\tilde{y}'_t Z(L) \xi_{t+1}] \right\} \\
&\quad + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3), \tag{2.19}
\end{aligned}$$

where

$$\begin{aligned}
\Phi &\equiv \bar{\lambda}'[D_y F + \beta D_{\tilde{y}} F] + \bar{\varphi}'[D_y g + \beta^{-1} D_{\tilde{y}} g], \\
H &\equiv \bar{\lambda}_k [D_{yy}^2 F^k + \beta D_{\tilde{y}\tilde{y}}^2 F^k] + \bar{\varphi}_i [D_{yy}^2 g^i + \beta^{-1} D_{\tilde{y}\tilde{y}}^2 g^i], \\
R &\equiv \bar{\lambda}_k D_{y\tilde{y}}^2 F^k + \bar{\varphi}_i \beta^{-1} D_{\tilde{y}y}^2 g^i, \\
Z(L) &\equiv \beta \bar{\lambda}_k D_{\tilde{y}\xi}^2 F^k + (\bar{\lambda}_k D_{y\xi}^2 F^k + \bar{\varphi}_i D_{y\xi}^2 g^i) \cdot L + \beta^{-1} \bar{\varphi}_i D_{\tilde{y}\xi}^2 g^i \cdot L^2.
\end{aligned}$$

Using (2.12), we furthermore observe that<sup>43</sup>

$$\Phi = -D_y \pi.$$

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<sup>42</sup>Note that we here include (2.4) among the constraints that a policy must satisfy. We shall call any evolution that satisfies (2.2)–(2.3) a “feasible” policy. Under this weaker assumption, the left-hand sides of (2.18) and (2.19) must instead be replaced by  $\beta^{-1}\bar{\varphi}'g(y_{t_0-1}, \xi_{t_0-1}; y_{t_0})$ .

<sup>43</sup>This is the point at which our calculations rely on the assumption that the steady state around which we compute our local approximations is optimal.

With this substitution in (2.19), we obtain an expression that can be solved for

$$E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} D_y \pi \cdot \tilde{y}_t,$$

which can in turn be used to substitute for the linear terms in (2.15). We thus obtain an alternative quadratic approximation to (2.1),<sup>44</sup>

$$V_{t_0} = \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\tilde{y}'_t Q \cdot \tilde{y}_t + 2\tilde{y}'_t R \tilde{y}_{t-1} + 2\tilde{y}'_t B(L) \xi_{t+1}] + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3), \quad (2.20)$$

where now

$$\begin{aligned} Q &\equiv D_{yy}^2 \pi + H, \\ B(L) &\equiv Z(L) + D_{y\xi}^2 \pi \cdot L. \end{aligned} \quad (2.21)$$

Since (2.20) involves no linear terms, it can be evaluated (up to a residual of order  $\mathcal{O}(\|\xi\|^3)$ ) using only a linear approximation to the evolution of  $\tilde{y}_t$  under a given policy rule.

It follows that a correct LQ approximation to the original problem is given by the problem of choosing a state-contingent evolution  $\{\tilde{y}_t\}$  for  $t \geq t_0$  to maximize the objective

$$V_{t_0}^Q(\tilde{y}; \xi) \equiv \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\tilde{y}'_t A(L) \tilde{y}_t + 2\tilde{y}'_t B(L) \xi_{t+1}] \quad (2.22)$$

subject to the constraints that

$$C(L) \tilde{y}_t = f_t, \quad (2.23)$$

$$E_t D(L) \tilde{y}_{t+1} = h_t \quad (2.24)$$

for all  $t \geq t_0$ , and the additional initial constraint that

$$D(L) \tilde{y}_{t_0} = \tilde{h}_{t_0}, \quad (2.25)$$

where now

$$A(L) \equiv Q + 2R \cdot L, \quad (2.26)$$

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<sup>44</sup>Here we include  $\bar{g}_{t_0}$  among the “terms independent of policy.” If we consider also policies that are not necessarily consistent with the initial pre-commitment, the left-hand side of (2.20) is more generally equal to  $V_{t_0} + \beta^{-1} \bar{\varphi}' g(y_{t_0-1}, \xi_{t_0-1}; y_{t_0})$ . This generalization of (2.20) is used in the derivation of equation (4.3) below.

$$C(L) \equiv D_y F + D_{\bar{y}} F \cdot L, \quad (2.27)$$

$$f_t \equiv -D_\xi F \cdot \xi_t,$$

$$D(L) \equiv D_{\bar{y}} g + D_y g \cdot L, \quad (2.28)$$

$$h_t \equiv -D_\xi g \cdot \xi_t, \quad (2.29)$$

$$\tilde{h}_{t_0} \equiv h_{t_0-1} + \bar{g}_{t_0}.$$

### 2.3 An Equivalent Lagrangian Approach

In the case that the objective (2.22) is concave,<sup>45</sup> the first-order conditions associated with the LQ problem just defined characterize the solution to that problem. Here we show that these linear equations also correspond to a local linear approximation to the first-order conditions associated with the exact problem, *i.e.*, the modified Ramsey policy problem defined in section 2.1, and hence that the solution to the LQ problem represents a local linear approximation to optimal policy from a timeless perspective.<sup>46</sup>

As already noted, the first-order conditions for the exact policy problem are obtained by differentiating the Lagrangian  $\mathcal{L}_{t_0}$  defined in (2.10). This yields the system of first-order conditions (2.11). The linearization of these first-order conditions around the optimal steady state is in turn the set of linear equations that would be obtained by differentiating a quadratic approximation to  $\mathcal{L}_{t_0}$  around that same steady state. Hence we are interested in computing such a local approximation, for the case in which  $y_t - \bar{y}$ ,  $\lambda_t - \bar{\lambda}$ , and  $\varphi_t - \bar{\varphi}$  are each of order  $\mathcal{O}(\|\xi\|)$  for all  $t$ . (Here the steady-state values of the Lagrange multipliers  $\bar{\lambda}, \bar{\varphi}$  are again given by the solution to equations (2.12) – (2.14).)

We may furthermore write the Lagrangian in the form

$$\mathcal{L}_{t_0} = \bar{\mathcal{L}}_{t_0} + \tilde{\mathcal{L}}_{t_0},$$

where

$$\bar{\mathcal{L}}_{t_0} = V_{t_0} + E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ \bar{\lambda}' F(y_t, \xi_t; y_{t-1}) + \beta^{-1} \bar{\varphi}' g(y_{t-1}, \xi_{t-1}; y_t) \right],$$

<sup>45</sup>The algebraic conditions under which this is so are discussed in the next section.

<sup>46</sup>See also Levine *et al.* (2007) for a similar discussion of the equivalence between our approach and the Lagrangian approach.

$$\tilde{\mathcal{L}}_{t_0} = E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ \tilde{\lambda}'_t F(y_t, \xi_t; y_{t-1}) + \beta^{-1} \tilde{\varphi}'_{t-1} g(y_{t-1}, \xi_{t-1}; y_t) \right],$$

$$\tilde{\lambda}_t \equiv \lambda_t - \bar{\lambda}, \quad \tilde{\varphi}_t \equiv \varphi_t - \bar{\varphi}.$$

We can then use equations (2.15) and (2.18) to show that the local quadratic approximation to  $\tilde{\mathcal{L}}_{t_0}$  is given by<sup>47</sup>

$$\bar{\mathcal{L}}_{t_0} = V_{t_0}^Q + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3).$$

In addition, the fact that  $\tilde{\lambda}_t, \tilde{\varphi}_t$  are both of order  $\mathcal{O}(\|\xi\|)$  means that a local quadratic approximation to the other term is given by

$$\tilde{\mathcal{L}}_{t_0} = E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ \tilde{\lambda}'_t \tilde{F}(y_t, \xi_t; y_{t-1}) + \beta^{-1} \tilde{\varphi}'_{t-1} \tilde{g}(y_{t-1}, \xi_{t-1}; y_t) \right] + \mathcal{O}(\|\xi\|^3),$$

where  $\tilde{F}$  and  $\tilde{g}$  are local linear approximations to the functions  $F$  and  $g$  respectively.

Hence the local quadratic approximation to the complete Lagrangian is given by

$$\begin{aligned} \mathcal{L}_{t_0} &= V_{t_0}^Q + E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ \tilde{\lambda}'_t \tilde{F}(y_t, \xi_t; y_{t-1}) + \beta^{-1} \tilde{\varphi}'_{t-1} \tilde{g}(y_{t-1}, \xi_{t-1}; y_t) \right] \\ &\quad + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \end{aligned} \tag{2.30}$$

But this is identical (up to terms independent of policy) to the Lagrangian for the LQ problem of maximizing  $V_{t_0}^Q$  subject to the linearized constraints. Hence the first-order conditions obtained from this approximate Lagrangian (which coincide with the local linear approximation to the first-order conditions for the exact problem) are identical to the first-order conditions for the LQ problem, and their solutions are identical as well.

### 3 Characterizing Optimal Policy

We now study necessary and sufficient conditions for a policy to solve the LQ problem of maximizing (2.22) subject to constraints (2.23) – (2.25). Let  $\mathcal{H}$  be the Hilbert space

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<sup>47</sup>It is worth noting that this equality holds in the case of all feasible policies, whether or not the policy is consistent with the initial pre-commitment (2.4). This is important for our discussion of the welfare evaluation of suboptimal policies in section 4.

of (real-valued) stochastic processes  $\{\tilde{y}_t\}$  such that

$$E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \tilde{y}'_t \tilde{y}_t < \infty. \quad (3.1)$$

We are interested in solutions to the LQ problem that satisfy the bound (3.1) because it guarantees that the objective  $V^Q$  is well-defined (and is generically required for it to be so). Of course, our LQ approximation to the original problem is only guaranteed to be accurate in the case that  $\tilde{y}_t$  is always sufficiently small; hence a solution to the LQ problem in which  $\tilde{y}_t$  grows without bound, but at a slow enough rate for (3.1) to be satisfied, need not correspond (even approximately) to any optimum (or local optimum) of the exact problem. In this section, however, we take the LQ problem at face value, and discuss the conditions under which it has a solution, despite the fact that we should in general only be interested in bounded solutions.

### 3.1 A Lagrangian Approach

The Lagrangian for this problem is given by

$$\mathcal{L}_{t_0}^Q = \frac{1}{2} \left\{ E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ \tilde{y}'_t A(L) \tilde{y}_t + 2\tilde{y}'_t B(L) \xi_{t+1} + 2\tilde{\lambda}'_t C(L) \tilde{y}_t + 2\beta^{-1} \tilde{\varphi}'_{t-1} D(L) \tilde{y}_t \right] \right\}.$$

(Note that this is just (2.30), omitting the terms independent of policy and those of third or higher order.) Differentiation of the Lagrangian then yields a system of linear first-order conditions

$$\begin{aligned} \frac{1}{2} E_t \{ [A(L) + A'(\beta L^{-1})] \tilde{y}_t \} + E_t [B(L) \xi_{t+1}] \\ + E_t [C'(\beta L^{-1}) \tilde{\lambda}_t] + \beta^{-1} D'(\beta L^{-1}) \tilde{\varphi}_{t-1} = 0 \end{aligned} \quad (3.2)$$

that must hold for each  $t \geq t_0$  under an optimal policy. (Here we use the notation  $M'$  for the transpose of a matrix  $M$ .) These conditions, together with (2.23) – (2.25), form a linear system to be solved for the joint evolution of the processes  $\{\tilde{y}_t, \tilde{\lambda}_t, \tilde{\varphi}_t\}$  given the exogenous disturbance processes  $\{\xi_t\}$  and the initial conditions  $\tilde{y}_{t_0-1}$  and the



initial pre-commitment  $\bar{g}_{t_0}$  (or  $\hat{h}_{t_0}$ ). This type of system of linear stochastic difference equations is easy to solve using standard methods.<sup>48</sup>

The first-order conditions (3.2) are easily shown to be *necessary* for optimality, but they are not generally *sufficient* for optimality as well; one must also verify that second-order conditions for optimality are satisfied. (In the case of an LQ problem, satisfaction of the second-order conditions implies global, and not just local, optimality; so we need not check any further conditions. But because our LQ problem is only a local approximation to the original policy problem, a global optimum of the LQ problem still may only correspond to a local optimum of the exact problem.) We next consider these additional conditions.

Let us consider the subspace  $\mathcal{H}_1 \subset \mathcal{H}$  of processes  $\hat{y} \in \mathcal{H}$  that satisfy the additional constraints

$$C(L)\hat{y}_t = 0 \tag{3.3}$$

$$E_t D(L)\hat{y}_{t+1} = 0 \tag{3.4}$$

for each date  $t \geq t_0$ , along with the initial commitments

$$D(L)\hat{y}_{t_0} = 0, \tag{3.5}$$

where we define  $\hat{y}_{t_0-1} \equiv 0$  in writing (3.3) for period  $t = t_0$  and in writing (3.5). This subspace is of interest because if a process  $\tilde{y} \in \mathcal{H}$  satisfies constraints (2.23) – (2.25), another process  $y \in \mathcal{H}$  with  $y_{t_0-1} = \tilde{y}_{t_0-1}$  satisfies those constraints as well if and only if  $y - \tilde{y} \in \mathcal{H}_1$ . We may now state our first main result.

**Proposition 1** *For  $\{\tilde{y}_t\} \in \mathcal{H}$  to maximize the quadratic form (2.22), subject to the constraints (2.23) – (2.25) given initial conditions  $\tilde{y}_{t_0-1}$  and  $\bar{g}_{t_0}$ , it is necessary and sufficient that (i) there exist Lagrange multiplier processes<sup>49</sup>  $\tilde{\varphi}, \tilde{\lambda} \in \mathcal{H}$  such that the processes  $\{\tilde{y}_t, \tilde{\varphi}_t, \tilde{\lambda}_t\}$  satisfy (3.2) for each  $t \geq t_0$ ; and (ii)*

$$V^Q(\hat{y}) \equiv V_{t_0}^Q(\hat{y}; 0) = \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\hat{y}'_t A(L) \hat{y}_t] \leq 0 \tag{3.6}$$

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<sup>48</sup>See, for example, Giannoni and Woodford (2002) for discussion of the solution of an equation system of this form using an eigenvector-decomposition method.

<sup>49</sup>Note that  $\tilde{\varphi}_t$  is also assumed to be defined for  $t = t_0 - 1$ .

for all processes  $\hat{y} \in \mathcal{H}_1$ , where in evaluating (3.6) we define  $\hat{y}_{t_0-1} \equiv 0$ . A process  $\{\tilde{y}_t\}$  with these properties is furthermore uniquely optimal if and only if

$$V^Q(\hat{y}) < 0 \tag{3.7}$$

for all processes  $\hat{y} \in \mathcal{H}_1$  that are non-zero almost surely.

The proof is given in the Appendix. The case in which the stronger condition (3.7) holds — *i.e.*, the quadratic form  $V^Q(\hat{y})$  is negative definite on the subspace  $\mathcal{H}_1$  — is the one of primary interest to us, since it is in this case that we know that the process  $\{\tilde{y}_t\}$  represents at least a local welfare maximum in the exact problem. In this case we can also show that pure randomization of policy reduces the welfare objective (2.22), and hence is locally welfare-reducing in the exact problem as well, as is discussed further in Benigno and Woodford (2005a).

### 3.2 A Dynamic Programming Approach

We can furthermore establish a useful characterization of the algebraic conditions under which the second-order conditions (3.7) are satisfied. These are most easily developed by considering the recursive formulation of our optimal policy problem presented in section 2.1.<sup>50</sup> Let us suppose that the exogenous state vector  $\xi_t$  evolves according to a linear law of motion

$$\xi_{t+1} = \Gamma \xi_t + \epsilon_{t+1}, \tag{3.8}$$

where  $\Gamma$  is a matrix, all of the eigenvalues of which have modulus less than  $\beta^{-1/2}$ , and  $\{\epsilon_t\}$  is an i.i.d. vector-valued random sequence, drawn each period from a distribution with mean zero and a variance-covariance matrix  $\Sigma$ .<sup>51</sup> In this case, our LQ

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<sup>50</sup>This section has been improved by the suggestions of Paul Levine and Joe Pearlman.

<sup>51</sup>These assumptions ensure that the process  $\{\xi_t\}$  satisfies a bound of the form (3.1). If we further wish to ensure that the disturbances are bounded, so that our local approximations can be expected to be accurate in the event of small enough disturbances, we may assume further that all eigenvalues of  $\Gamma$  have a modulus less than 1, and that  $\epsilon_{t+1}$  is drawn from a distribution with bounded support. We may assume that, like the other structural relations in this section, (3.8) is merely a local linear approximation. Finally, note that the assumption of a law of motion of the form (3.8) allows for disturbances with arbitrarily complex forms of serial correlation, simply by adding elements to the vector  $\xi_t$  reflecting past exogenous states.

approximate policy problem has a recursive formulation, in which the continuation problem from any period  $t$  forward depends on the extended state vector

$$\mathbf{z}_t \equiv \begin{bmatrix} \tilde{y}_{t-1} \\ \tilde{h}_t \\ \xi_t \\ \xi_{t-1} \end{bmatrix}. \quad (3.9)$$

Let  $\bar{V}^Q(\mathbf{z}_t)$  denote the maximum attainable value of the continuation objective  $V_t^Q$ , if the process  $\{\tilde{y}_\tau\}$  from date  $t$  onward is chosen to satisfy constraints (2.23)–(2.24) for all  $\tau \geq t$ , an initial precommitment of the form

$$D(L)\tilde{y}_t = \tilde{h}_t, \quad (3.10)$$

and the bound (3.1). As usual in an LQ problem of this form, it can be shown that the value function is a quadratic function of the extended state vector,

$$\bar{V}^Q(\mathbf{z}_t) = \frac{1}{2} \mathbf{z}_t' P \mathbf{z}_t, \quad (3.11)$$

where  $P$  is a symmetric matrix to be determined. In characterizing the solution to the problem, it is useful to introduce notation for partitions of the matrix  $P$ . Let  $P_{ij}$  (for  $i, j = 1, 2, 3, 4$ ) be the 16 blocks obtained when  $P$  is partitioned in both directions conformably with the partition of  $\mathbf{z}_t$  in (3.9), and let

$$\mathbf{P}_i \equiv [P_{i1} \ P_{i2} \ P_{i3} \ P_{i4}]$$

(for  $i = 1, 2, 3, 4$ ) be the four blocks obtained when  $P$  is partitioned only vertically.

In the recursive formulation of the approximate LQ problem, in each period  $t$ ,  $\tilde{y}_t$  is chosen, and a precommitment  $\tilde{h}_{t+1}(\xi_{t+1})$  is chosen for each possible state in the period  $t + 1$  continuation, so as to maximize

$$\frac{1}{2} \tilde{y}_t' A(L) \tilde{y}_t + E_t[\tilde{y}_t' B(L) \xi_{t+1}] + \beta E_t \bar{V}^Q(\mathbf{z}_{t+1}), \quad (3.12)$$

subject to the constraints that  $\tilde{y}_t$  satisfy (2.23) and (3.10), and that the choices of  $\{\tilde{h}_{t+1}(\xi_{t+1})\}$  satisfy

$$E_t \tilde{h}_{t+1} = h_t. \quad (3.13)$$

To simplify the discussion, we shall further assume that

$$\text{rank} \begin{bmatrix} C_0 \\ D_0 \end{bmatrix} = n_F + n_g, \quad (3.14)$$

where here and below we write lag polynomials in the form  $X(L) = \sum_j X_j L^j$ . This condition implies that the constraints (2.23) and (3.10) include neither any redundant constraints nor any constraints that are inconsistent in the case of a generic state  $\mathbf{z}_t$ .

The first-order conditions for the optimal choice of  $\tilde{y}_t$  in this single-period problem are of the form

$$[A_0 + (1/2)A_1L]\tilde{y}_t + E_t[B(L)\xi_{t+1}] + \beta\mathbf{P}_1 E_t\mathbf{z}_{t+1} + C'_0\tilde{\lambda}_t + D'_0\tilde{\psi}_t = 0, \quad (3.15)$$

where  $\tilde{\lambda}_t, \tilde{\psi}_t$  are the Lagrange multipliers associated with constraints (2.23) and (3.10) respectively. Condition (3.15) together with the constraints (2.23) and (3.10) constitute a system of  $n = n_y + n_F + n_g$  linear equations to solve for  $\tilde{y}_t, \tilde{\lambda}_t$ , and  $\tilde{\psi}_t$  as functions of  $\mathbf{z}_t$ . This system can be written in the matrix form  $M y_t^\dagger = -G \mathbf{z}_t$ , where

$$M \equiv \begin{bmatrix} A_0 + \beta P_{11} & C'_0 & D'_0 \\ C_0 & 0 & 0 \\ D_0 & 0 & 0 \end{bmatrix}, \quad y_t^\dagger \equiv \begin{bmatrix} \tilde{y}_t \\ \tilde{\lambda}_t \\ \tilde{\psi}_t \end{bmatrix}, \quad (3.16)$$

and  $G$  is a matrix of coefficients, the first two columns of which (of particular interest here) are

$$G_1 \equiv \begin{bmatrix} (1/2)A_1 \\ C_1 \\ D_1 \end{bmatrix}, \quad G_2 \equiv \begin{bmatrix} 0 \\ 0 \\ -I \end{bmatrix}.$$

This has a determinate solution if and only if  $M$  is non-singular. This is evidently a necessary condition for strict concavity of the policy problem, and we shall assume that it holds in the remainder of our discussion.<sup>52</sup> Given this assumption, the unique solution is

$$y_t^\dagger = -M^{-1} G \mathbf{z}_t. \quad (3.17)$$

The first-order conditions for the optimal choice of the precommitments  $\{\tilde{h}_{t+1}(\xi_{t+1})\}$  are that

$$\beta\mathbf{P}_2 \mathbf{z}_{t+1} = -\tilde{\varphi}_t \quad (3.18)$$

in each possible state  $\xi_{t+1}$  that can succeed the given state in period  $t$ , where  $\tilde{\varphi}_t$  is the Lagrange multiplier associated with constraint (3.13); note that the value of  $\tilde{\varphi}_t$

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<sup>52</sup>We are actually only interested in whether there exists a unique solution for  $\tilde{y}_t$ . However, condition (3.14) implies that there can be no vector  $y^\dagger \neq 0$  such that  $My^\dagger = 0$ , unless it involves  $\tilde{y} \neq 0$ . Thus if  $M$  is singular, there are necessary multiple solutions for  $\tilde{y}_t$  if there are any solutions at all, and not just multiple solutions for the Lagrange multipliers.

depends only on the state in period  $t$ . The fact that the left-hand side of (3.18) must be the same in each state  $\xi_{t+1}$  implies that

$$P_{22} [\tilde{h}_{t+1} - h_t] + P_{23} \epsilon_{t+1} = 0$$

in each state. This allows a determinate solution for  $\tilde{h}_{t+1}$  if and only if  $P_{22}$  is non-singular; this too is evidently a necessary condition for concavity, and is assumed from here on.<sup>53</sup> Under this assumption, (3.18) together with (3.10) implies that

$$\tilde{h}_{t+1} = h_t - P_{22}^{-1} P_{23} \epsilon_{t+1}. \quad (3.19)$$

We can also solve uniquely for the Lagrange multiplier,

$$\begin{aligned} \tilde{\varphi}_t &= -\beta \mathbf{P}_2 E_t \mathbf{z}_{t+1} \\ &= -\beta P_{21} \tilde{y}_t - \beta P_{22} h_t - \beta [P_{23} \Gamma + P_{24}] \xi_t. \end{aligned} \quad (3.20)$$

Equations (3.17) and (3.19) completely describe the optimal dynamics of the variables  $\{\tilde{y}_t, \tilde{h}_t\}$ , starting from some initial conditions  $(\tilde{y}_{t_0-1}, \tilde{h}_{t_0})$ , given the evolution of the exogenous states  $\{\xi_t\}$ . The system consisting of these solutions for  $\tilde{y}_t$  and  $\tilde{h}_{t+1}(\xi_{t+1})$ , together with the law of motion (3.8), can be written in the form

$$\mathbf{z}_{t+1} = \Phi \mathbf{z}_t + \Psi \epsilon_{t+1}, \quad (3.21)$$

for certain matrices  $\Phi$  and  $\Psi$ . If we partition  $\Phi$  in the same way as  $P$ , it follows from the form of the solutions obtained above that  $\Phi_{ij} = 0$  for all  $i \geq 2, j \leq 2$ . From this (together with our assumption about the eigenvalues of  $\Gamma$ ) it follows that all eigenvalues of  $\Phi$  have modulus less than  $\beta^{-1/2}$  if and only if all eigenvalues of

$$\Phi_{11} \equiv [-I \ 0 \ 0] \ M^{-1} G_1 \quad (3.22)$$

have this property. Hence there exists a determinate solution to the first-order conditions for optimal policy, *i.e.*, a unique solution satisfying the bound (3.1), if and only if  $M$  and  $P_{22}$  are non-singular matrices, and all eigenvalues of  $\Phi_{11}$  have modulus less than  $\beta^{-1/2}$ .

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<sup>53</sup>If  $P_{22}$  is singular, it is obvious that there are multiple solutions for  $\tilde{h}_{t+1}(\xi_{t+1})$  consistent with the first-order conditions, but one might wonder if these correspond to multiple state-contingent evolutions  $\{\tilde{y}_t\}$ . In fact they do, for a single state-contingent evolution  $\{\tilde{y}_t\}$  is consistent with only one process  $\{\tilde{h}_t\}$ , which can be determined from (3.10).

Note that the solution (3.21) involves elements of the matrix  $P$ . We can solve for those elements of  $P$  in the following way. It follows from the assumed representation (3.11) for the value function that the vector of partial derivatives with respect to  $\tilde{y}_{t-1}$  will equal

$$\bar{V}_1^Q = \mathbf{P}_1 \mathbf{z}_t.$$

On the other hand, application of the envelope theorem to the problem (3.12) implies that

$$\bar{V}_1^Q = G'_1 y_t^\dagger = -G'_1 M^{-1} G \mathbf{z}_t. \quad (3.23)$$

Equating the corresponding coefficients in these two representations, we observe that

$$P_{1j} = -G'_1 M^{-1} G_j$$

for  $j = 1, 2, 3, 4$ . A similar argument implies that

$$P_{2j} = -G'_2 M^{-1} G_j \quad (3.24)$$

for  $j = 1, 2, 3, 4$ .

These expressions involve the matrix  $M$ , which depends on  $P_{11}$ ; but the system

$$P_{11} = -G'_1 M(P_{11})^{-1} G_1 \quad (3.25)$$

is a set of  $n_y^2$  equations to solve for the  $n_y^2$  elements of  $P_{11}$ .<sup>54</sup> Once we have solved for  $P_{11}$ , we know the matrix  $M$ , and can solve for the other elements of  $P$ . In particular, we can solve for

$$P_{22} = -G'_2 M^{-1} G_2, \quad (3.26)$$

and check whether it is non-singular, as required in (3.19). The other elements of  $P$  can be solved for using the same method.<sup>55</sup>

Thus far, we have discussed only the implications of the first-order conditions for the single-period optimization problem. Again, the question arises whether a solution to the first-order conditions corresponds to a maximum of (3.12). The second-order conditions for a finite-dimensional optimization problem are well-known. First, the objective is strictly concave in  $\tilde{y}_t$  if and only if the matrix  $A_0 + \beta P_{11}$  is such that

$$\tilde{y}' [A_0 + \beta P_{11}] \tilde{y} < 0$$

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<sup>54</sup>Actually, because  $P_{11}$  is symmetric, and the system (3.25) has the same symmetry, we need only solve a system of  $n(n+1)/2$  equations for  $n(n+1)/2$  independent quantities.

<sup>55</sup>Details of the algebra are provided in a note on computational issues available from the authors.

for all  $\tilde{y} \neq 0$  such that

$$C_0 \tilde{y} = 0, \quad D_0 \tilde{y} = 0.$$

Using a result of Debreu (1952),<sup>56</sup> we can state algebraic conditions on these matrices that are easily checked. For each  $r$  such that  $n_F + n_g + 1 \leq r \leq n_y$ , let  $M_r$  be the lower-right square block of  $M$  of size  $n_F + n_g + r$ .<sup>57</sup> Then the concavity condition stated above holds if and only if  $\det M_r$  has the same sign as  $(-1)^r$ , for each  $n_F + n_g + 1 \leq r \leq n_y$ . Note that in the case that policy is *unidimensional* — meaning that there is a single instrument to set each period, which suffices to determine the evolution of the endogenous variables, so that  $n_F + n_g = n_y - 1$  — then this requirement reduces to the single condition that the determinant of  $M$  have the same sign as  $(-1)^{n_y}$ .

Second, in each possible state  $\xi_{t+1}$  in the following period, the continuation objective  $\bar{V}^Q(\mathbf{z}_{t+1})$  is a concave function of  $\tilde{h}_{t+1}(\xi_{t+1})$  if and only if the submatrix  $P_{22}$  is *negative definite*, i.e., such that  $\tilde{h}' P_{22} \tilde{h} < 0$  for all  $\tilde{h} \neq 0$ . This condition is also straightforward to check using the Debreu theorem: the principal minors of  $P_{22}$  must have alternating signs.

These two conditions are obviously necessary for strict concavity of the single-period problem, and hence for strict concavity of the infinite-horizon optimal policy problem. In fact, they are also sufficient, yielding the following result.

**Proposition 2** *Suppose that the exogenous disturbances have a law of motion of the form (3.8), where  $\Gamma$  is a matrix the eigenvalues of which all have modulus less than  $\beta^{-1/2}$ , and that the constraints satisfy the rank condition (3.14), where  $n_F + n_g < n_y$ . Then the LQ policy problem has a determinate solution, given by (3.21), if and only if (i) there exists a solution  $P_{11}$  to equations (3.25) such that for each of the minors of the matrix  $M$  defined in (3.16),  $\det M_r$  has the same sign as  $(-1)^r$ , for each  $n_F + n_g + 1 \leq r \leq n_y$ ; (ii) the eigenvalues of the matrix  $\Phi_{11}$  defined in (3.22) all have modulus less than  $\beta^{-1/2}$ ; and (iii) the matrix  $P_{22}$  defined in (3.26) is negative definite, i.e., is such that its  $r$ th principle minor has the same sign as  $(-1)^r$ , for each  $1 \leq r \leq n_g$ .*

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<sup>56</sup>See also Theorem 1.E.17 of Takayama (1985).

<sup>57</sup>Given (3.14), we can order the elements of  $\tilde{y}_t$  so that the left  $(n_F + n_g) \times (n_F + n_g)$  block of the matrix in (3.14) is non-singular, and we assume that this has been done when forming these submatrices.

The proof of this proposition is also given in the Appendix. Note that the conditions stated in the proposition are necessary and sufficient both for the existence of a determinate solution to the first-order conditions, and for the quadratic form  $V^Q(\psi)$  to satisfy the strict concavity condition (3.7). In the case that either condition (i) or (iii) is violated, there may exist a determinate solution to the first-order conditions, but it will not represent an optimum, owing to violation of the second-order conditions.

The fact that condition (iii) is needed in addition to conditions (i)–(ii) in order to ensure that we have a concave problem indicates an important respect in which the theory of LQ optimization with forward-looking constraints is not a trivial generalization of the standard theory for backward-looking problems, since conditions (i)–(ii) are sufficient in a backward-looking problem of the kind treated by Magill (1977).<sup>58</sup> It also shows that the second-order conditions for a stochastic problem are more complex than they would be in the case of a deterministic policy problem (again, unlike what is true of purely backward-looking LQ problems). For in a deterministic version of our problem with forward-looking constraints, conditions (i)–(ii) would also be sufficient for concavity, and thus for the solution to the first-order conditions to represent an optimum.

In a deterministic version of the problem — where we not only assume that  $\xi_t = 0$  each period, but we restrict our attention to policies under which the evolution of the variables  $\{\tilde{y}_t\}$  is purely deterministic (and hence perfectly forecastable), so that we seek to characterize the optimal *perfect foresight equilibrium*, without addressing the question whether this is also optimal among the larger set of possible *rational-expectations equilibria*.<sup>59</sup> — the constraints on possible equilibria are the purely backward-looking constraints (2.23) and

$$D(L)\tilde{y}_t = \tilde{h}_t \tag{3.27}$$

for each  $t \geq t_0$ , where we specify  $\tilde{h}_t = h_{t-1} = 0$  for all  $t \geq t_0 + 1$ . This is a purely

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<sup>58</sup>See Levine *et al.* for a derivation of the second-order conditions for a backward-looking, deterministic LQ problem, using what is essentially a discrete-time version of the approach of Magill. In some cases, conditions (i)–(ii) are both necessary and sufficient for concavity, even in the presence of forward-looking constraints. The problem treated in Benigno and Woodford (2005a) is an example of this kind. Note that in that paper an alternative, frequency-domain characterization of the conditions for concavity is used, that is discussed more generally in Benigno and Woodford (2006b).

<sup>59</sup>Additional equilibria can be attained, by randomization of policy, even in the case that there are no exogenous random disturbances. This may or may not allow an increase in welfare relative to the optimal deterministic policy.



backward-looking problem, so that the standard second-order conditions apply. And it should be obvious that, as there is no longer a choice of  $\tilde{h}_{t+1}(\xi_{t+1})$  to be made each period, our argument above for the necessity of condition (iii) would not apply.

But conditions (i)–(ii) are not generally a sufficient condition to guarantee that (3.7) is satisfied, in the presence of forward-looking constraints (2.24), if policy randomization is allowed.<sup>60</sup> Because constraints (2.24) need hold only in expected value, random policy may be able to vary the paths of the endogenous variables (in some states of the world) in directions that would not be possible in the corresponding deterministic problem, and this makes the algebraic conditions required for (3.7) to hold more stringent. Specifically, the value function for the continuation problem must be a strictly concave function of the state-contingent pre-commitment  $\tilde{h}_{t+1}$  made for the following period, or it is possible to randomize  $\tilde{h}_{t+1}$  (requiring a corresponding randomization of subsequent policy) without changing the fact that constraint (2.24) is satisfied in period  $t$ . Hence condition (iii) is necessary in the stochastic case.<sup>61</sup> It can also easily be shown that condition (iii) is not implied in general by conditions (i)–(ii).

A simple example may clarify this point. Suppose that  $y_t$  has two elements, and that the only constraint on what policy can achieve is a single, forward-looking constraint

$$E_t[\delta\tilde{y}_{1,t} - \tilde{y}_{1,t+1}] = 0 \tag{3.28}$$

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<sup>60</sup>Our remarks here apply even in the case that the “fundamental” disturbances  $\{\xi_t\}$  are purely deterministic; what matters is whether policy may be contingent upon random events. As is discussed further in Benigno and Woodford (2005a), when the second-order conditions fail to hold, policy randomization can be welfare-improving, even when the random variations in policy are unrelated to any variation in fundamentals.

<sup>61</sup>Levine *et al.* (2007) provide a different argument for a condition similar to our condition (iii) as a necessary condition for optimality in a model with a forward-looking constraint, which does not require a consideration of stochastic policy. They consider Ramsey-optimal policy rather than optimality from a timeless perspective; that is, they assume no initial precommitment (2.25). In this case, the deterministic optimal policy problem is like the one considered above, except that (3.27) need hold only in periods  $t \geq t_0 + 1$ ; the optimal policy is then the same as in the backward-looking problem just discussed, except that instead of taking  $\tilde{h}_{t_0}$  as given, one is free to choose  $\tilde{h}_{t_0}$  so as to maximize (2.22). This latter problem has a solution only if the value function  $\bar{V}_{t_0}^Q$  is bounded above, for a given vector  $\tilde{y}_{t_0-1}$ , and this is true in general only if it is a strictly concave function of  $\tilde{h}_{t_0}$ . The validity of this argument, however, depends on considering an exact LQ problem, rather than an LQ local approximation to a problem that may have different global behavior.

for all  $t \geq t_0$ , where  $\delta < \beta^{-1/2}$ . (The path of  $\{\tilde{y}_{2,t}\}$  can be freely chosen, subject to the bound (3.1).) An initial pre-commitment specifies the value that  $\tilde{y}_{1,t_0}$  must have. In the corresponding deterministic problem, constraint (3.28) implies that one must have

$$\tilde{y}_{1,t+1} = \delta \tilde{y}_{1,t}$$

for each  $t \geq t_0$ , and this, together with the pre-commitment, uniquely determines the entire path of the sequence  $\{\tilde{y}_{1,t}\}$  that must be brought about by deterministic policy. Hence the second-order condition for the deterministic problem requires only that the objective be a concave function of the path of  $\{\tilde{y}_{2,t}\}$ . But if random policies are considered, it is also possible for  $\{\tilde{y}_{1,t}\}$  to evolve in accordance with any law of motion

$$\tilde{y}_{1,t+1} = \delta \tilde{y}_{1,t} + \epsilon_{t+1},$$

where  $\{\epsilon_t\}$  is any martingale difference sequence with a suitable bound on its asymptotic variance; in this simple example, the set of possible evolutions  $\{\tilde{y}_{1,t}\}$  is independent of the evolution chosen for  $\{\tilde{y}_{2,t}\}$ . Whether randomization of the path of  $\{\tilde{y}_{1,t}\}$  can increase the value of the policy objective obviously depends on terms in the objective involving the path of  $\{\tilde{y}_{1,t}\}$  (including cross terms), and not just the terms involving the path of  $\{\tilde{y}_{2,t}\}$ . Hence the conditions required for a concave optimization problem are more stringent in this case.<sup>62</sup>

## 4 Welfare Evaluation of Alternative Policy Rules

We have argued that another advantage of our approach is that it can be used not only to derive a linear approximation to a fully optimal policy commitment, but also to compute approximate welfare comparisons between alternative rules (neither of which may be fully optimal), that will correctly rank these rules in the case that random disturbances are small enough. Because empirically realistic models are inevitably fairly complex, a fully optimal policy rule is likely to be too complex to represent a realistic policy proposal; hence comparisons among alternative simple (though suboptimal) rules are of considerable practical interest. Here we discuss how this can be done.

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<sup>62</sup>In the Appendix, we illustrate the application of the conditions in Proposition 2 to this example.

We do not propose to simply evaluate (a local approximation to) expected discounted utility  $V_{t_0}$  under a candidate policy rule, because the optimal policy locally characterized above (*i.e.*, optimal policy “from a timeless perspective”) does not maximize this objective; hence ranking rules according to this criterion would lead to the embarrassing conclusion that there exist policies better than the optimal policy. (We could, of course, define “optimal policy” as the policy that maximizes  $V_{t_0}$ ; but this would result in a time-inconsistent policy recommendation, as noted earlier.) Thus we wish to use a criterion that ranks rules according to how close they come to solving the recursive policy problem defined in section 2.1, rather than how close they come to maximizing  $V_{t_0}$ .

Of course, if we restrict our attention to policies that necessarily satisfy the initial pre-commitment (2.4), there is no problem; our optimal rule will be the one that maximizes  $V_{t_0}$ , or (in the case of small enough shocks) the one that maximizes  $V_{t_0}^Q$ . But *simple* policy rules are unlikely to precisely satisfy (2.4); thus in order to be able to select the best rule from some simple class, we need an alternative criterion, one that is defined for *all* policies that are close enough to being optimal, in a sense that is to be defined. At the same time, we wish it to be a criterion the maximization of which implies that one has solved the constrained optimization problem defined in section 2.1.

## 4.1 A Lagrangian Approach

Our Lagrangian characterization of optimal policy suggests such a criterion. The timelessly optimal policy from date  $t_0$  onward — that is, the policy that maximizes  $V_{t_0}$  subject to the initial constraint (2.4) in addition to the feasibility constraints (2.2)–(2.3) — is also the policy that maximizes the Lagrangian

$$V_{t_0}^{mod} \equiv V_{t_0} + \beta^{-1} \varphi'_{t_0-1} g(y_{t_0-1}, \xi_{t_0-1}; y_{t_0}), \quad (4.1)$$

where  $\varphi_{t_0-1}$  is the vector of Lagrange multipliers associated with the initial constraint (2.4). This is a function that coincides (up to a constant) with the objective  $V_{t_0}$  in the case of policies satisfying the constraint (2.4), but that is defined more generally, and that is maximized over the broader class of feasible policies by the timelessly optimal policy. Hence an appropriate criterion to use in ranking alternative policies is the value of  $V_{t_0}^{mod}$  associated with each one. This criterion penalizes policies that

fail to satisfy the initial pre-commitment (2.4), by exactly the amount by which a previously *anticipated* deviation of that kind would have reduced the expected utility of the representative household.

In the case of any policy that satisfies the feasibility constraints (2.2)–(2.3) for all  $t \geq t_0$ , we observe that

$$\begin{aligned} V_{t_0}^{mod} &= \bar{\mathcal{L}}_{t_0} + \beta^{-1} \tilde{\varphi}'_{t_0-1} g(y_{t_0-1}, \xi_{t_0-1}; y_{t_0}) \\ &= V_{t_0}^Q + \beta^{-1} \tilde{\varphi}'_{t_0-1} \tilde{g}(y_{t_0-1}, \xi_{t_0-1}; y_{t_0}) + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3) \\ &= V_{t_0}^Q + \beta^{-1} \tilde{\varphi}'_{t_0-1} D_{\tilde{y}} g \cdot \tilde{y}_{t_0} + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \end{aligned}$$

This suggests that in the case of small enough shocks, the ranking of alternative policies in terms of  $V_{t_0}^{mod}$  will correspond to the ranking in terms of the welfare measure

$$W_{t_0} \equiv V_{t_0}^Q + \beta^{-1} \tilde{\varphi}'_{t_0-1} D_{\tilde{y}} g \cdot \tilde{y}_{t_0}. \quad (4.2)$$

Note that in this derivation we have assumed that  $\tilde{y}_t = \mathcal{O}(\|\xi\|)$ . This will be true in the equilibrium associated with any (sufficiently differentiable) policy rule that is *consistent with the optimal steady state* in the absence of random disturbances. We shall restrict attention to policy rules of this kind. Note that while this is an important restriction, it does not preclude consideration of extremely simple rules; and it is a property of the simple rules of greatest interest, *i.e.*, those that come closest to being optimal among rules of that degree of complexity.

In expression (4.1), and hence in (4.2),  $\varphi_{t_0-1}$  is the Lagrange multiplier associated with constraint (2.4) under the optimal policy. However, in order to evaluate  $W_{t_0}$  to second-order accuracy, it suffices to have a first-order approximation to this multiplier. Such an approximation is given by the multiplier  $\tilde{\varphi}_{t_0-1}$  associated with the constraint (2.25) of the LQ problem. Thus we need only solve the LQ problem, as discussed in the previous section — obtaining a value for  $\tilde{\varphi}_{t_0-1}$  along with our solution for the optimal evolution  $\{y_t\}$  — in order to determine the value of  $W_{t_0}$ .

Moreover, we observe that in the characterization given in the previous section of the solution to the LQ problem,  $\tilde{\varphi}_{t_0-1} = \mathcal{O}(\|\xi\|)$ .<sup>63</sup> Thus a solution for the equilibrium evolution  $\{\tilde{y}_t\}$  under a given policy that is accurate to first order suffices to evaluate the second term in (4.2) to second-order accuracy. Hence  $W_{t_0}$  inherits this

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<sup>63</sup>This follows from solution (3.17) for the Lagrange multiplier associated with the initial pre-commitment.

property of  $V_{t_0}^Q$ , and it suffices to compute a linear approximation to the equilibrium dynamics  $\{\tilde{y}_t\}$  under each candidate policy rule in order to evaluate  $W_{t_0}$  to second-order accuracy. We can therefore obtain an approximation solution for  $\{\tilde{y}_t\}$  under a given policy by solving the linearized structural equations (2.23)–(2.24), together with the policy rule, and use this solution in evaluating  $W_{t_0}$ . In this way welfare comparisons among alternative policies are possible, to second-order accuracy, using linear approximations to the model structural relations and a quadratic welfare objective.

Moreover, we can evaluate  $W_{t_0}$  to second-order accuracy using only a linear approximation to the policy rule. This has important computational advantages. For example, if we wish to find the optimal policy rule from among the family of simple rules of the form  $i_t = \phi(y_t)$ , where  $i_t$  is a policy instrument, and we are content to evaluate  $V_{t_0}^{mod}$  to second-order accuracy, then it suffices to search over the family of linear policy rules<sup>64</sup>

$$\tilde{i}_t = f' \tilde{y}_t,$$

parameterized by the vector of coefficients  $f$ . There are no possible second-order (or larger) welfare gains resulting from nonlinearities in the policy rule.

It is important to note that these conclusions obtain *only* because we evaluate welfare taking into account the welfare losses that would result from a violation of the initial pre-commitment if it were to have been anticipated. Some would prefer to evaluate alternative simple policy rules by computing the expected value of  $V_{t_0}$  (rather than  $V_{t_0}^{mod}$ ) associated with each rule (e.g., Schmitt-Grohé and Uribe, 2007). As noted above, this alternative criterion is one under which the optimal rule from a timeless perspective can be dominated by other rules, a point stressed by Blake (2001) and Jensen and McCallum (2002), among others. The alternative criterion is also one that cannot be evaluated to second-order accuracy using only a first-order solution for the equilibrium evolution under a given policy. For a general feasible policy — consistent with the optimal steady state, but not necessarily consistent

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<sup>64</sup>Here we restrict attention to rules that are consistent with the optimal steady state, so that the intercept term is zero when the rule is expressed in terms of deviations from steady-state values. Note that a rule without this property will result in lower welfare, in the case of any small enough disturbances.

with the initial pre-commitment (2.4) — we can show that<sup>65</sup>

$$V_{t_0} = V_{t_0}^Q - \beta^{-1} \bar{\varphi}' D_{\tilde{y}} g \cdot \tilde{y}_{t_0} + \text{t.i.p.} + \mathcal{O}(\|\xi\|^3). \quad (4.3)$$

The first term on the right-hand side of this expression is purely quadratic (has zero linear terms), but this is not true of the second term, if the initial pre-commitment is binding under the optimal policy. Evaluation of the second term to second-order accuracy requires a second-order approximation to the evolution  $\{y_t\}$  under the policy of interest; there is thus no alternative to the use of higher-order perturbation solution methods as illustrated by Schmitt-Grohé and Uribe, and nonlinear terms in the policy rule generally matter for welfare.<sup>66</sup>

In expression (4.2), the value of the multiplier  $\tilde{\varphi}_{t_0-1}$  depends on the economy's initial state and on the value of the initial pre-commitment  $\bar{y}_{t_0}$ . However, we wish to be able to rank alternative rules for an economy in which no such commitment may exist prior to the adoption of the policy rule. We can avoid having to make reference to any historically given pre-commitment by assuming a self-consistent constraint of the form (2.5).

If we define a new extended state vector

$$\hat{\mathbf{z}}_t \equiv \begin{bmatrix} \tilde{y}_{t-1} \\ \hat{h}(\xi_t, \xi_{t-1}) \\ \xi_t \\ \xi_{t-1} \end{bmatrix},$$

where<sup>67</sup>

$$\hat{h}(\xi_t, \xi_{t-1}) \equiv h_{t-1} - P_{22}^{-1} P_{23} (\xi_t - \Gamma \xi_{t-1}),$$

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<sup>65</sup>Here we use the more general form of (2.20) mentioned in footnote 42.

<sup>66</sup>Damjanovic *et al.* (2008) show that one can instead use an LQ approximation to evaluate time-invariant policy rules under an alternative criterion, which computes the expected value of  $V_{t_0}$  under a probability distribution for initial conditions that is independent of the policy rule considered, as in the calculations here, but rather under the ergodic distribution for the endogenous variables associated with the particular time-invariant policy that is to be evaluated. This criterion has the unappealing feature of giving a rule that leads to different long-run average values of an endogenous variable (e.g., the capital stock) “credit” for a higher initial average value of the variable as well. It also cannot be applied to evaluate non-stationary policies, or even time-invariant policies that imply non-stationary dynamics of endogenous variables, such as the optimal policy in Benigno and Woodford (2003).

<sup>67</sup>Here it should be recalled that  $h_{t-1}$  is a linear function of  $\xi_{t-1}$ , defined in (2.29).

then it follows from (3.19) that under the solution to the recursive policy problem,  $\mathbf{z}_t = \hat{\mathbf{z}}_t$  for each  $t \geq t_0 + 1$ . (However,  $\hat{\mathbf{z}}_t$ , unlike  $\mathbf{z}_t$ , is a function solely of  $\tilde{y}_{t-1}$  and the history of the exogenous disturbances.) Hence

$$\tilde{h}_{t_0} = \hat{h}(\xi_t, \xi_{t-1}) \quad (4.4)$$

is a self-consistent constraint of the form (2.5).

If we assume an initial pre-commitment specified in this way, it also follows from (3.17) that

$$\tilde{\psi}_t = [0 \ 0 \ -I] M^{-1} G \hat{\mathbf{z}}_t \quad (4.5)$$

is the Lagrange multiplier associated with the pre-commitment each period in the recursive problem. Moreover, because the only constraint on the way in which  $\tilde{h}_{t+1}(\xi_{t+1})$  can be chosen for the following period is given by the expected-value constraint (3.13), the first-order conditions for optimal policy imply that  $\tilde{\psi}_t = E_{t-1} \tilde{\psi}_t$  for each  $t \geq t_0 + 1$ ,<sup>68</sup> and hence that

$$\begin{aligned} \tilde{\psi}_t &= [0 \ 0 \ -I] M^{-1} G E_{t-1} \hat{\mathbf{z}}_t \\ &= \tilde{\psi}(\tilde{y}_{t-1}, \xi_{t-1}) \equiv [0 \ 0 \ -I] M^{-1} G \begin{bmatrix} \tilde{y}_{t-1} \\ h_{t-1} \\ \Gamma \xi_{t-1} \\ \xi_{t-1} \end{bmatrix}. \end{aligned}$$

Consistency of this result with (4.5) implies that the right-hand-side of (4.5) must be equivalent to  $\tilde{\psi}(\tilde{y}_{t-1}, \xi_{t-1})$ ; that is, that the coefficients multiplying  $\tilde{y}_{t-1}$ ,  $\xi_t$ , and  $\xi_{t-1}$  must be the same in both expressions. But since (4.5) must hold at  $t = t_0$  as well, in the case of an initial pre-commitment (4.4), and not only for  $t \geq t_0 + 1$ , it follows that under such a pre-commitment,

$$\tilde{\psi}_t = \tilde{\psi}(\tilde{y}_{t-1}, \xi_{t-1})$$

for all  $t \geq t_0$ . In the case that  $t = t_0$ , the multiplier  $\tilde{\psi}_{t_0}$  associated with the initial pre-commitment is the one that is denoted  $\beta^{-1} \tilde{\varphi}_{t_0-1}$  in (2.30) and in (4.2). Thus we can write

$$\tilde{\varphi}_{t_0-1} = \varphi^*(\mathbf{y}_{t_0-1}) \equiv \beta \tilde{\psi}(\tilde{y}_{t_0-1}, \xi_{t_0-1}). \quad (4.6)$$

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<sup>68</sup>In fact, one can show that  $\tilde{\psi}_t = \beta^{-1} \tilde{\varphi}_{t-1}$  for each  $t \geq t_0 + 1$ . This follows from differentiation of the value function  $V^Q(\mathbf{z}_{t+1})$  with respect to  $\tilde{h}_{t+1}$  using the envelope theorem, and comparison of the result with (3.18).

Then we can write<sup>69</sup>

$$W_{t_0} = W(\tilde{y}; \xi_{t_0}, \mathbf{y}_{t_0-1}) \equiv V_{t_0}^Q + \beta^{-1} \varphi^*(\mathbf{y}_{t_0-1})' D_{\tilde{y}} g \cdot \tilde{y}_{t_0}. \quad (4.7)$$

This gives us an expression for our welfare measure purely in terms of the history and subsequent evolution of the extended state vector.<sup>70</sup>

## 4.2 A Time-Invariant Criterion for Ranking Alternative Rules

Let us suppose that we are interested in evaluating a policy rule  $r$  that implies an equilibrium evolution of the endogenous variables of the form<sup>71</sup>

$$y_t = \phi_r(\xi_t, \mathbf{y}_{t-1}).$$

This (together with the law of motion for the exogenous disturbances) then implies a law of motion for the complete extended state vector

$$\mathbf{y}_t = \psi_r(\xi_t, \mathbf{y}_{t-1}). \quad (4.8)$$

Using this law of motion, we can evaluate (4.7), obtaining

$$W_{t_0} = W_r(\xi_{t_0}, \mathbf{y}_{t_0-1}).$$

We can do this for any rule  $r$  of the assumed type, and hence we can define an optimization problem

$$\max_{r \in \mathcal{R}} W_r(\xi_{t_0}, \mathbf{y}_{t_0-1}) \quad (4.9)$$

in order to determine the optimal rule from among the members of some family of rules  $\mathcal{R}$ .

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<sup>69</sup>In writing the function  $W(\cdot)$ , and others that follow, we suppress the argument  $\xi$ , as the evolution of the exogenous disturbances is the same in the case of each of the alternative policies under consideration.

<sup>70</sup>Note that it is possible to solve for the initial Lagrange multipliers  $\varphi^*(\mathbf{y}_{t_0-1})$  using only the values of  $\tilde{y}_{t_0-1}$  and of  $\xi_{t_0-1}$ . It is not necessary to simulate the optimal equilibrium dynamics over a lengthy “estimation period” prior to the date  $t_0$  at which the new policy is to commence, as proposed by Juillard and Pelgrin (2006).

<sup>71</sup>This assumption that  $y_t$  depends only on the state variables indicated is without loss of generality, as we can extend the vector  $\mathbf{y}_t$  if necessary in order for this to be so.



However, the solution to problem (4.9) may well depend on the initial conditions  $\mathbf{y}_{t_0-1}$  and  $\xi_{t_0}$  for which  $W_{t_0}$  is evaluated.<sup>72</sup> This leads to the possibility of an unappealing degree of arbitrariness of the choice that would be recommended from within some family of simple rules, as well as time inconsistency of the policy recommendation: a rule chosen at date  $t_0$  on the ground that it solves problem (4.9) need not be found to also solve the corresponding problem at some later date, though the calculation at date  $t_0$  assumes that rule  $r$  is to be followed forever. One way of avoiding this might be to assume that one should choose the rule that would be judged best in the case of initial conditions consistent with the optimal steady state, whether the economy's actual initial state is that one or not;<sup>73</sup> that is, one would choose the rule that solves the problem

$$\max_{r \in \mathcal{R}} W_r(0, \bar{\mathbf{y}}).$$

This choice would not be time-inconsistent, but the choice is still an arbitrary one. In particular, the decision to evaluate  $W_r$  assuming initial conditions consistent with the steady state — when in fact the state of the economy will fluctuate on both sides of the steady-state position — favors rules  $r$  for which  $W_r$  is a less concave function of the initial condition.

The criterion that we find most appealing is accordingly to integrate over a distribution of possible initial conditions, rather than evaluating  $W_r$  at the economy's actual state at the time of the choice, or at any other single state (such as the optimal steady state). Suppose that in the case of the optimal policy rule  $r^*$ , the law of motion (4.8) implies that the evolution of the extended state vector  $\{\mathbf{y}_t\}$  is *stationary*.<sup>74</sup> In this case, there exists a well-defined invariant (or unconditional) probability distribution  $\mu$  for the possible values of  $\mathbf{y}_t$  under the optimal policy.<sup>75</sup> Then we can define the optimal policy rule within some class of simple rules  $\mathcal{R}$  as the one that solves the problem

$$\max_{r \in \mathcal{R}} E_\mu[\bar{W}_r(\mathbf{y}_t)], \tag{4.10}$$

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<sup>72</sup>This is not a problem if the family of rules  $\mathcal{R}$  includes a fully optimal rule  $r^*$ , since the same rule  $r^*$  solves the problem (2.7) for all possible values of the initial conditions. But the result can easily depend on the initial conditions if we restrict attention to a family of suboptimal rules.

<sup>73</sup>This approach is proposed by Schmitt-Grohé and Uribe (2007), though they use  $V_{t_0}$  rather than  $V_{t_0}^{mod}$  as the criterion to be maximized.

<sup>74</sup>Benigno and Woodford (2005a) provide an example of an optimal monetary stabilization policy problem in which this is case.

<sup>75</sup>We discuss the computation of the relevant properties of this invariant measure in the Appendix.

where<sup>76</sup>

$$\bar{W}_r(\mathbf{y}_t) \equiv E_t W_r(\xi_{t+1}, \mathbf{y}_t). \quad (4.11)$$

Because of the linearity of our approximate characterization of optimal policy, the calculations required in order to evaluate  $E_\mu[W_r]$  to second-order accuracy are straightforward; these are illustrated in Benigno and Woodford (2005a, sec. 5).

The most important case in which the method just described cannot be applied is when some of the elements of  $\{\mathbf{y}_t\}$  possess unit roots, though all elements are at least difference-stationary (and some of the non-stationary elements may be cointegrated). Note that it is possible for even the equilibrium under optimal policy to have this property, consistent with our assumption of the bound (3.1).<sup>77</sup> There is a question in such a case whether our local approximation to the problem should remain an accurate approximation, but this is not a problem in the case that random disturbances occur in only a *finite* number of periods, so LQ problems of this kind may be of practical interest.

Let us suppose that those elements which possess unit roots are pure random walks (*i.e.*, with zero drift).<sup>78</sup> We can in such a case decompose the extended state vector as

$$\mathbf{y}_t = \mathbf{y}_t^{tr} + \mathbf{y}_t^{cyc},$$

where

$$\mathbf{y}_t^{tr} \equiv \lim_{T \rightarrow \infty} E_t \mathbf{y}_T$$

is the Beveridge-Nelson (1981) “trend” component, and the “cyclical” component  $\mathbf{y}_t^{cyc}$  will still be a stationary process. Moreover, the evolution of the cyclical component as a function of the exogenous disturbances under the optimal policy will be independent of the assumed initial value of the trend component (though not of the

<sup>76</sup>Recall that we assume that the exogenous disturbance process  $\{\xi_t\}$  is Markovian, and that  $\xi_t$  is included among the elements of  $\mathbf{y}_t$ . Hence  $\mathbf{y}_t$  contains all relevant elements of the period  $t$  information set for the calculation of this conditional expectation.

<sup>77</sup>Benigno and Woodford (2003) provide an example of an optimal stabilization policy problem in which the LQ approximate problem has this property. In this example, the unit root is associated with the dynamics of the level of real public debt, which display a unit root under optimal policy for the same reason as in the classic analysis of optimal tax smoothing by Barro (1979) and Sargent (1987, chap. XV).

<sup>78</sup>We may suppose that any deterministic trend under optimal policy has been eliminated by local expansion around a deterministic solution with constant trend growth, so that there is zero trend in the state variables  $\{\tilde{y}_t\}$  expressed as deviations from that deterministic solution.

initial value of the cyclical component). It follows that we can define an invariant distribution  $\mu$  for the possible values of  $\mathbf{y}_t^{cyc}$  under the optimal policy, that is independent of the assumed value for the trend component. Then for any assumed initial value for the trend component  $\mathbf{y}_{t_0-1}^{tr}$ , we can define the optimal policy rule within the class  $\mathcal{R}$  as the one that solves the problem

$$\max_{r \in \mathcal{R}} \Omega_r(\mathbf{y}_{t_0-1}^{tr}) \equiv E_\mu[\bar{W}_r(\mathbf{y}_{t_0-1})], \quad (4.12)$$

a generalization of (4.10).<sup>79</sup>

It might seem in this case that our criterion is again dependent on initial conditions, just as with the criterion (4.9) proposed first. The following result shows that this is not the case.

**Lemma 3** *Suppose that under optimal policy, the extended state vector  $\mathbf{y}_t$  consists entirely of components that are either (i) stationary, or (ii) pure random walks. Suppose also that the class of policy rules  $\mathcal{R}$  is such that each rule in the class implies convergence to the same long-run values of the state variables as under optimal policy, in the absence of stochastic disturbances, so that the initial value of the trend component  $\mathbf{y}_{t_0-1}^{tr}$  is the same regardless of the rule  $r$  that is considered. Then for any rule  $r \in \mathcal{R}$ , the objective  $\Omega_r(\mathbf{y}_{t_0-1}^{tr})$  defined in (4.12) can be decomposed into two parts,*

$$\Omega_r(\mathbf{y}_{t_0-1}^{tr}) = \Omega^1(\mathbf{y}_{t_0-1}^{tr}) + \Omega_r^2, \quad (4.13)$$

where the first component is the same for all rules in this class, while the second component is independent of the initial condition  $\mathbf{y}_{t_0-1}^{tr}$ .

Hence the criterion (4.12) establishes the same ranking of alternative rules, regardless of the initial condition. The proof of this result is given in the Appendix.

## 5 Applications

The approach expounded here has already proven fruitful in a number of applications to problems of optimal monetary and fiscal policy. Benigno and Woodford (2005a)

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<sup>79</sup>In the case that all elements of  $\mathbf{y}_t$  are stationary,  $\mathbf{y}_t^{tr}$  is simply a constant, and all variations in  $\mathbf{y}_t$  correspond to variations in  $\mathbf{y}_t^{cyc}$ . In this case, (4.12) is equivalent to the previous criterion (4.10).

use this method to derive an LQ approximation to the problem of optimal monetary stabilization policy in a DSGE model with monopolistic competition, Calvo-style staggered price-setting, and a variety of exogenous disturbances to preferences, technology, and fiscal policy. Unlike the LQ method used by Rotemberg and Woodford (1997) and Woodford (2002), the present method is applicable even in the case of (possibly substantial) distortions even in the absence of shocks, owing to market power or distorting taxes. The quadratic stabilization objective obtained is of the form

$$-\frac{1}{2}E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left[ q_{\pi} \pi_t^2 + q_y (\hat{Y}_t - \hat{Y}_t^*)^2 \right], \quad (5.1)$$

where  $\pi_t$  is the inflation rate between periods  $t - 1$  and  $t$ ,  $\hat{Y}_t$  is the log deviation of aggregate real output from trend,  $\hat{Y}_t^*$  is a target level of output that depends purely on the exogenous real disturbances,  $0 < \beta < 1$  is the representative household's discount factor, and the weights  $q_{\pi}, q_y$  are functions of model parameters (both positive if steady-state distortions are not severe). The single linear constraint corresponds to the familiar “new Keynesian Phillips curve,”

$$\pi_t = \kappa [\hat{Y}_t - \hat{Y}_t^*] + \beta E_t \pi_{t+1} + u_t, \quad (5.2)$$

where  $\kappa > 0$  is a function of model parameters and the “cost-push” term  $u_t$  is a linear function of the various exogenous real disturbances.

The resulting LQ problem is of a form that has already been extensively studied in the literature on optimal monetary stabilization policy,<sup>80</sup> and so the ways in which the parameterization of the objective and constraint shape the character of optimal policy is well understood once the problem is stated in this form. The analysis in Benigno and Woodford (2005a), however, explains the microeconomic determinants of these factors. For example, it provides an interpretation of the “cost-push” disturbances that play a crucial role in familiar discussions of the tradeoffs between inflation and output stabilization, and shows that the cost-push effects of most types of shocks are larger the more distorted is the economy's steady state; and it explains the relative weight that should be assigned to the output-gap stabilization objective, showing that this need not be positive in the case of a sufficiently distorted economy. (Indeed, if distortions are severe, the quadratic objective can fail to be concave, so that a small

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<sup>80</sup>See, e.g., , Clarida *et al.* (1999) and Woodford (2003, chap. 7).

amount of policy randomization can be welfare-improving.) Benigno and Woodford (2005b) extend the analysis to the case in which both wages and prices are sticky, obtaining a generalization of (5.1) in which a third quadratic loss term appears, proportional to squared deviations of nominal wage inflation from zero. This shows that the analysis by Erceg *et al.* (2000) of the tradeoff between stabilization of wage inflation and price inflation applies also to economies with distorted steady states, though the policy tradeoffs are complicated by the presence of cost-push terms that do not appear in those authors' analysis of the case of an undistorted steady state. Montoro (2007) extends the analysis to allow for real disturbances to the relative supply price of oil.

An important limitation of the LQ method of Rotemberg and Woodford (1997), that restricts attention to cases in which the utility gradient is zero in the steady state, is that it cannot easily be applied to analyses of optimal policy for open economies; for in an open economy, domestic production and consumption cannot be equated, and the marginal utility associated with a change in either individually will inevitably be non-zero in any reasonable case. The method proposed here instead allows LQ analyses of optimal policy also in the case of open economies.

Benigno and Benigno (2006) analyze policy coordination between two national monetary authorities which each seek to maximize the welfare of their own country's representative household, and show that it is possible to locally characterize each authority's aims by a quadratic stabilization objective. Previous LQ analyses of policy coordination have often assumed an objective of the form (5.1) for each national authority, but with the nation's own inflation rate and output being the arguments in each case. Benigno and Benigno instead show that household utility maximization would correspond to a quadratic objective for each authority with terms penalizing fluctuations in *both* domestic and foreign inflation (but with different weights on the two terms for the distinct national authorities), and similarly with terms penalizing fluctuations in both domestic and foreign output (again with different weights in the case of the two authorities). They also show that each authority's stabilization objective should contain a term penalizing departures of the terms of trade from a "target" level (that depends on exogenous disturbances), and show how both the weight placed on this additional objective and the nature of variation in the terms of trade "target" depend on underlying micro-foundations. De Paoli (2004) similarly shows how the analysis of Benigno and Woodford (2005a) can be extended to a

small open economy, requiring the addition of a terms-of-trade (or real-exchange-rate) stabilization objective to the two terms shown in (5.1).

Another advantage of the fact that the present method applies to economies with a distorted steady state is that it can be used to analyze optimal tax smoothing when only distorting taxes are available as sources of government revenue, after the fashion of Barro (1979) and Sargent (1987, chap. XV), and allows the theory of tax smoothing to be integrated with the theory of monetary stabilization policy. Benigno and Woodford (2003) extend the analysis of Benigno and Woodford (2005a) to the case of an economy with only distorting taxes, and show that the problem of choosing jointly optimal monetary and fiscal policies can also be treated within an LQ framework that nests standard analyses of tax smoothing (with flexible prices, so that real effects of monetary policy are ignored) and of monetary policy (with lump-sum taxes, so that fiscal effects of monetary policy can be ignored) as special cases. Notably, they find that allowing for tax distortions introduces no additional stabilization goals into the quadratic objective (5.1). Instead, the benefits of tax smoothing are represented by the penalty on squared departures of equilibrium output from its “target” level; tax variations can increase the average size of this term, because of the effects of the level of distorting taxes on equilibrium output (which occur due to a “cost-push” effect of tax rates in the generalized version of the constraint (5.2)). Benigno and De Paoli (2005) extend this analysis to treat optimal monetary and fiscal policy in a small open economy, while Ferrero (2005) analyzes optimal monetary and fiscal policy in a monetary union with separate national fiscal authorities. Berriel and Sinigaglia (2008) extend the analysis to the case of an economy with multiple sectors that differ in the degree of stickiness of prices.

All of the analyses just mentioned involve fairly simple DSGE models, in which it is possible to derive the coefficients of the LQ approximate policy problem by hand. In the case of larger (and more realistic) models of the kind that are now being estimated for use in practical policy analysis, such calculations are likely to be tedious. Nonetheless, it is an advantage of our method that it is straightforward to apply it even to fairly complex models and fairly general specifications of disturbances. Altissimo *et al.* (2005) describe computer code that executes the calculations explained above, for a general nonlinear problem with an arbitrary number of state variables, and demonstrate its application to two important extensions of the work described above, an analysis of optimal monetary policy in the presence of non-trivial frictions

of the kind that result in a transactions demand for money, and an analysis of optimal monetary policy for the empirical model of Smets and Wouters. Cúrdia (2007) illustrates the application of the methods proposed here to another fairly complex model, namely, a model of “sudden stops” in a small emerging-market economy; in particular, the method explained in section 4 is used to evaluate alternative simple policy rules for such a setting. We believe that it should similarly be practical to apply these methods to a wide variety of other models of interest to policy institutions.

# A Appendix: Proofs and Derivations

## A.1 Proposition 1

Recall that  $\mathcal{H}$  is the Hilbert space of (real-valued) stochastic processes  $\{\tilde{y}_t\}$  such that

$$E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \tilde{y}'_t \tilde{y}_t < \infty, \quad (\text{A.1})$$

and  $\mathcal{H}_1 \subset \mathcal{H}$  is the subspace of sequences  $\hat{y} \in \mathcal{H}$  that satisfy the additional constraints

$$C(L)\hat{y}_t = 0 \quad (\text{A.2})$$

$$E_t D(L)\hat{y}_{t+1} = 0 \quad (\text{A.3})$$

for each date  $t \geq t_0$ , along with the initial commitments

$$D(L)\hat{y}_{t_0} = 0, \quad (\text{A.4})$$

where we define  $\hat{y}_{t_0-1} \equiv 0$  in writing (A.2) for period  $t = t_0$  and in writing (A.4).

**Proposition 1** *For  $\{\tilde{y}_t\} \in \mathcal{H}$  to maximize the quadratic form (2.22), subject to the constraints (2.23) – (2.25) given initial conditions  $\tilde{y}_{t_0-1}$  and  $\bar{g}_{t_0}$ , it is necessary and sufficient that (i) there exist Lagrange multiplier processes<sup>81</sup>  $\tilde{\varphi}, \tilde{\lambda} \in \mathcal{H}$  such that the processes  $\{\tilde{y}_t, \tilde{\varphi}_t, \tilde{\lambda}_t\}$  satisfy (3.2) for each  $t \geq t_0$ ; and (ii)*

$$V^Q(\hat{y}) \equiv V_{t_0}^Q(\hat{y}; 0) = \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\hat{y}'_t A(L)\hat{y}_t] \leq 0 \quad (\text{A.5})$$

for all processes  $\hat{y} \in \mathcal{H}_1$ , where in evaluating (A.5) we define  $\hat{y}_{t_0-1} \equiv 0$ . A process  $\{\tilde{y}_t\}$  with these properties is furthermore uniquely optimal if and only if

$$V^Q(\hat{y}) < 0 \quad (\text{A.6})$$

for all processes  $\hat{y} \in \mathcal{H}_1$  that are non-zero almost surely.

PROOF: We have already remarked on the necessity of the first-order conditions (i). To prove the necessity of the second-order condition (ii) as well, let  $\{\tilde{y}_t\} \in \mathcal{H}$ , and consider the the perturbed process

$$y_t = \tilde{y}_t + \hat{y}_t \quad (\text{A.7})$$

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<sup>81</sup>Note that  $\tilde{\varphi}_t$  is also assumed to be defined for  $t = t_0 - 1$ .



for all  $t \geq t_0 - 1$ , where  $\{\hat{y}_t\}$  belongs to  $\mathcal{H}_1$  and we define  $\hat{y}_{t_0-1} \equiv 0$ . This construction guarantees that if the process  $\{\tilde{y}_t\}$  satisfies the constraints (2.23) – (2.25), so does the process  $\{y_t\}$ .

We note that

$$\begin{aligned} V_{t_0}^Q(y; \xi) &= V_{t_0}^Q(\tilde{y}; \xi) + \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\hat{y}'_t A(L) \tilde{y}_t + \tilde{y}'_t A(L) \hat{y}_t + 2\hat{y}'_t B(L) \xi_{t+1}] \\ &\quad + \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} [\hat{y}'_t A(L) \hat{y}_t]. \end{aligned}$$

The second term on the right-hand side is furthermore equal to

$$\begin{aligned} &\frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \hat{y}'_t \cdot \{ [A(L) + A'(\beta L^{-1})] \tilde{y}_t + 2B(L) \xi_{t+1} \} \\ &= -E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \hat{y}'_t \cdot \left\{ C'(\beta L^{-1}) \tilde{\lambda}_t + \beta^{-1} D'(\beta L^{-1}) \tilde{\varphi}_{t-1} \right\} \\ &= -E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \left\{ \tilde{\lambda}'_t C(L) \hat{y}_t + \beta^{-1} \tilde{\varphi}'_{t-1} D(L) \hat{y}_t \right\}, \end{aligned}$$

where we use the first-order conditions (3.2) to establish the first equality, and conditions (3.3) – (3.5) to establish the final equality.

Thus for any feasible process  $\tilde{y}$  and any perturbation (A.7) defined by a process  $\hat{y}$  belonging to  $\mathcal{H}_1$ ,

$$V_{t_0}^Q(y; \xi) = V_{t_0}^Q(\tilde{y}; \xi) + V^Q(\hat{y}). \quad (\text{A.8})$$

It follows that if there were to exist any  $\hat{y} \in \mathcal{H}_1$  for which  $V^Q(\hat{y}) > 0$ , the plan  $\tilde{y}$  could not be optimal. But as this is true regardless of what plan  $\tilde{y}$  may be, (A.5) is necessary for optimality. Furthermore, if there were to exist a non-zero  $\hat{y}$  for which  $V^Q(\hat{y}) = 0$ , it would be possible to construct a perturbation  $y$  (not equal to  $\tilde{y}$  almost surely at all dates) that would achieve an equally high level of welfare. Hence the stronger version of the second-order conditions (A.6) must hold for all  $\hat{y}$  not equal to zero almost surely, in order for  $\{\tilde{y}_t\}$  to be a unique optimum.

One easily sees from the same calculation that these conditions are also sufficient for an optimum. Let  $\{\tilde{y}_t\}$  be a process consistent with the constraints of the LQ problem. Then any alternative process  $\{y_t\}$  that is also consistent with those constraints can be written in the form (A.7), where  $\hat{y}$  is some element of  $\mathcal{H}_1$ . If the

first-order conditions (3.2) are satisfied by the process  $\{\tilde{y}_t\}$ , we can again establish (A.8). Condition (A.5) then implies that no alternative process is preferable to  $\{\tilde{y}_t\}$ , while (A.6) would imply that  $\{\tilde{y}_t\}$  is superior to any alternative that is not equal to  $\tilde{y}$  almost surely.

## A.2 Proposition 2

**Proposition 2** *Suppose that the exogenous disturbances have a law of motion of the form (3.8), where  $\Gamma$  is a matrix the eigenvalues of which all have modulus less than  $\beta^{-1/2}$ , and that the constraints satisfy the rank condition (3.14), where  $n_F + n_g < n_y$ . Then the LQ policy problem has a determinate solution, given by (3.21), if and only if (i) there exists a solution  $P_{11}$  to equations (3.25) such that for each of the minors of the matrix  $M$  defined in (3.16),  $\det M_r$  has the same sign as  $(-1)^r$ , for each  $n_F + n_g + 1 \leq r \leq n_y$ ; (ii) the eigenvalues of the matrix  $\Phi_{11}$  defined in (3.22) all have modulus less than  $\beta^{-1/2}$ ; and (iii) the matrix  $P_{22}$  defined in (3.26) is negative definite, i.e., is such that its  $r$ th principle minor has the same sign as  $(-1)^r$ , for each  $1 \leq r \leq n_g$ .*

PROOF: (1) The discussion in the text has already established the necessity of each of conditions (i)–(iii), so it remains only to show that they are also sufficient for the solution (3.21) to represent a solution to the original infinite-horizon optimal policy problem. We shall do this by establishing that conditions (i)–(iii) imply that the sufficient conditions of Proposition 1 are satisfied by this solution.

We begin by establishing that the processes  $\{\tilde{y}_t, \tilde{\lambda}_t, \tilde{\varphi}_t\}$  associated with the solution (3.21) satisfy the first-order conditions (3.2) for the infinite-horizon problem. We have already shown in the text that under conditions (i)–(iii), there exists a determinate solution (3.21) for the dynamics of  $\{\mathbf{z}_t\}$ , that it satisfies the bound (3.1) along with the constraints (2.23)–(2.25), and that associated with it are a unique system of Lagrange multipliers  $\{\tilde{\lambda}_t, \tilde{\psi}_t, \tilde{\varphi}_t\}$ , the solution for which has also been explained in the text. We wish to show that these processes must satisfy (3.2) for each  $t \geq t_0$ .

By construction, the processes  $\{y_t^\dagger\}$  satisfy the first-order conditions (3.15) for each  $t \geq t_0$ . Moreover, it follows from (3.23) that

$$\mathbf{P}_1 E_t \mathbf{z}_{t+1} = G_1' E_t y_{t+1}^\dagger.$$

Substituting this into (3.15), we obtain

$$\begin{aligned} \frac{1}{2}E_t\{[A(L) + A'(\beta L^{-1})]\tilde{y}_t\} + E_t[B(L)\xi_{t+1}] \\ + E_t[C'(\beta L^{-1})\tilde{\lambda}_t] + E_t[D'(\beta L^{-1})\tilde{\psi}_t] = 0 \end{aligned} \quad (\text{A.9})$$

for each  $t \geq t_0$ .

Differentiating  $\bar{V}^Q(\mathbf{z}_t)$  with respect to  $\tilde{h}_t$ , and using the envelope theorem as in the derivation of (3.23), we obtain  $\bar{V}_2^Q = -\tilde{\psi}_t$ , from which we conclude that

$$\mathbf{P}_2 \mathbf{z}_t = -\tilde{\psi}_t$$

for each  $t \geq t_0$ . Comparison with first-order condition (3.18) for the optimal choice of  $\tilde{h}_{t+1}$  in the recursive policy problem indicates that

$$\tilde{\psi}_t = \beta^{-1}\tilde{\varphi}_{t-1} \quad (\text{A.10})$$

for each  $t \geq t_0 + 1$ . We may assume (as a definition of  $\tilde{\varphi}_{t_0-1}$ <sup>82</sup>) that (A.10) holds when  $t = t_0$  as well. Then use of (A.10) to substitute for the process  $\{\tilde{\psi}_t\}$  in (A.9) yields (3.2), which accordingly must hold for each  $t \geq t_0$ . Hence the processes constructed to satisfy the first-order conditions of the recursive policy problem must satisfy the first-order conditions for the infinite-horizon policy problem characterized in section 3.1 as well.

(2) It remains to show that conditions (i)–(iii) also imply that the strict concavity condition (A.6) is satisfied. Let us consider an arbitrary process  $\tilde{y} \in \mathcal{H}_1$ , and associated with it define the process  $\tilde{h}$  by

$$\tilde{h}_t = D(L)\tilde{y}_t \quad (\text{A.11})$$

for each  $t \geq t_0 + 1$ , and by the stipulation that  $\tilde{h}_{t_0} = 0$ . We thus obtain a pair of processes satisfying

$$C(L)\tilde{y}_t = 0, \quad (\text{A.12})$$

$$D(L)\tilde{y}_t = \tilde{h}_t, \quad (\text{A.13})$$

$$E_t\tilde{h}_{t+1} = 0 \quad (\text{A.14})$$

---

<sup>82</sup>Note that  $\tilde{\varphi}_{t_0-1}$  has no other meaning in the analysis of the recursive policy problem presented in section 3.2.

for all  $t \geq t_0$ . These are furthermore an example of a process  $\{\mathbf{z}_t\}$  consistent with the constraints of the recursive policy problem, in the case that  $\xi_t = 0$  at all times and the initial precommitment is given by  $\tilde{h}_{t_0} = 0$ .

We note that the analysis given in the text of the single-period problem of maximizing (3.12), applied to the special case in which  $\xi_t = 0$  at all times,<sup>83</sup> implies that for any values of  $\tilde{y}_{t-1}$  and  $\tilde{h}_t$ , the maximum possible attainable value of the objective

$$\frac{1}{2}\tilde{y}'_t A(L)\tilde{y}_t + \frac{\beta}{2}E_t[\mathbf{z}'_{t+1}P\mathbf{z}_{t+1}]$$

consistent with constraints (A.12)–(A.14) is equal to

$$\frac{1}{2}\mathbf{z}'_t P\mathbf{z}_t;$$

and this value is attained only if

$$\mathbf{z}_{t+1} = \Phi \mathbf{z}_t$$

with certainty, which is to say, only if

$$\tilde{y}_t = \Phi_{11}\tilde{y}_{t-1} + \Phi_{12}\tilde{h}_t \tag{A.15}$$

and

$$\tilde{h}_{t+1} = 0 \tag{A.16}$$

in each possible state in period  $t + 1$ .

Thus the fact that the processes  $\{\tilde{y}_t, \tilde{h}_t\}$  satisfy (A.12)–(A.14) for all  $t \geq t_0$  implies that

$$\frac{1}{2}\tilde{y}'_t A(L)\tilde{y}_t + \frac{\beta}{2}E_t[\mathbf{z}'_{t+1}P\mathbf{z}_{t+1}] \leq \frac{1}{2}\mathbf{z}'_t P\mathbf{z}_t$$

for all  $t \geq t_0$ , and that the inequality is strict unless (A.15)–(A.16) hold. Now if conditions (A.15)–(A.16) hold for all  $t \geq t_0$ ,  $\tilde{y}_t = 0$  at all times. Thus in the case that  $\tilde{y}_t$  is not equal to zero almost surely for all  $t$ , there must be at least one date  $t_1$  such that at least one of these conditions is violated with positive probability when  $t = t_1$ . In that case, there must be some  $k > 0$  such that

$$E_{t_0} \left\{ \frac{1}{2}\tilde{y}'_{t_1} A(L)\tilde{y}_{t_1} + \frac{\beta}{2}\mathbf{z}'_{t_1+1}P\mathbf{z}_{t_1+1} \right\} \frac{1}{2} \leq E_{t_0}\mathbf{z}'_{t_1} P\mathbf{z}_{t_1} - k.$$

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<sup>83</sup>It follows from the usual principle of certainty equivalence for LQ problems that the matrices characterizing the solution to this problem do not depend on the value of the variance-covariance matrix  $\Sigma$  for the disturbances. In fact, it is easily observed that the derivations given in the text would apply equally to a problem in which  $\xi_t = 0$  at all times.

It then follows, by summing these inequalities (appropriately discounted) for successive periods, that

$$E_{t_0} \sum_{t=t_0}^T \beta^{t-t_0} \frac{1}{2} \tilde{y}'_t A(L) \tilde{y}_t + \frac{\beta^{T+1-t_0}}{2} E_{t_0} \mathbf{z}'_{T+1} P \mathbf{z}_{T+1} \leq \frac{1}{2} \mathbf{z}'_{t_0} P \mathbf{z}_{t_0} - k = -k, \quad (\text{A.17})$$

for all  $T \geq t_1$ .

As we have stipulated that the process  $\tilde{y}$  is an element of  $\mathcal{H}_1$ , and thus satisfies the bound (3.1), we necessarily have

$$\lim_{T \rightarrow \infty} \beta^{T+1} E_{t_0} \mathbf{z}'_{T+1} P \mathbf{z}_{T+1} = 0.$$

(Note that it follows from (A.11) that the elements of  $\tilde{h}$  cannot grow asymptotically at a faster rate than do the elements of  $\tilde{y}$ .) It then follows from (A.17) that

$$\limsup_{T \rightarrow \infty} E_{t_0} \sum_{t=t_0}^T \beta^{t-t_0} \frac{1}{2} \tilde{y}'_t A(L) \tilde{y}_t \leq -k. \quad (\text{A.18})$$

But since it follows from the assumption that  $\tilde{y}$  satisfies (3.1) that the series in (A.18) has a limit, this limit must be no greater than  $-k$ . Hence  $\tilde{y}$  satisfies (A.6), and all of the sufficient conditions of Proposition 1 have been verified. This establishes the proposition.

EXAMPLE: Suppose that  $y_t$  has two elements, that the objective of policy is to maximize

$$V_{t_0}^Q(\tilde{y}) \equiv \frac{1}{2} E_{t_0} \sum_{t=t_0}^{\infty} \beta^{t-t_0} \tilde{y}'_t A \tilde{y}_t, \quad (\text{A.19})$$

where  $A$  is a symmetric  $2 \times 2$  matrix, and that the only constraint on what policy can achieve is a single, forward-looking constraint

$$E_t[\delta \tilde{y}_{1,t} - \tilde{y}_{1,t+1}] = 0 \quad (\text{A.20})$$

for all  $t \geq t_0$ , where  $|\delta| < \beta^{-1/2}$ . There are no exogenous disturbances, but the expectations appear because we wish to consider the possibility of (arbitrarily) randomized policies. We assume an initial pre-commitment of the form

$$\tilde{y}_{1,t_0} = \delta \tilde{y}_{1,t_0-1} + \tilde{h}_{t_0}, \quad (\text{A.21})$$

for some quantity  $\tilde{h}_{t_0}$ .

In the case that policy is restricted to be deterministic, the constraint completely determines the path of  $\{\tilde{y}_{1t}\}$ ; the only (perfect foresight) sequence consistent with the initial pre-commitment and the forward-looking constraint is the one in which

$$\tilde{y}_{1,t} = [\delta\tilde{y}_{1,t_0-1} + \tilde{h}_{t_0}]\delta^{t-t_0}$$

for all  $t \geq t_0$ . The problem then reduces to the choice of a sequence  $\{\tilde{y}_{2,t}\}$ , constrained only by the bound (3.1), so as to maximize the objective. This is obviously a concave problem if and only if  $\tilde{y}'A\tilde{y}$  is a concave function of  $\tilde{y}_2$  for a given value of  $\tilde{y}_1$ . This in turn is true if and only if  $A_{22} < 0$ ; the other elements of  $A$  are irrelevant.

If instead we allow random policies, the condition just derived is no longer sufficient for concavity (though still necessary). One can show that the problem is concave if and only if  $A$  is a negative definite matrix. This is obviously a sufficient condition (as it implies that (A.19) is concave for arbitrary sequences). To show that it is also necessary, suppose instead that it is not true. Then there exists a vector  $v \neq 0$  such that  $v'Av \geq 0$ . Now let  $\{\bar{y}_t\}$  be any process satisfying the constraints (3.1), (A.20), and (A.21), and consider the alternative process  $\{\tilde{y}_t\}$  generated by the law of motion

$$\tilde{y}_t = \bar{y}_t + \delta(\tilde{y}_{t-1} - \bar{y}_{t-1}) + v\epsilon_t$$

for each  $t \geq t_0 + 1$ , starting from the initial condition (A.21), where  $\{\epsilon_t\}$  is a (scalar-valued) martingale-difference sequence satisfying the bound (3.1). One can easily show that the process  $\{\tilde{y}_t\}$  satisfies (3.1), (A.20), and (A.21) as well; moreover, the value of the objective in the case of this process satisfies

$$\begin{aligned} V_{t_0}^Q(\tilde{y}) &= V_{t_0}^Q(\bar{y}) + (1 - \beta\delta^2)^{-1} v'Av E_{t_0} \sum_{t=t_0+1}^{\infty} \beta^t \epsilon_t^2 \\ &\geq V_{t_0}^Q(\bar{y}). \end{aligned}$$

Since we can construct an alternative policy that is at least as good in the case of *any* policy, there is no uniquely optimal policy in such a case; and in addition, we have shown that arbitrary randomization of policy is possible without welfare loss.

Let us examine how these results compare with the conditions stated in Proposition 2. In this example, condition (3.25) states that

$$P_{11} = \alpha \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix},$$

where

$$\alpha = -\delta^2 [M^{-1}]_{33}.$$

This form for  $P_{11}$  implies in turn that  $M$  is invertible as long as  $A_{22} \neq 0$ , and that in that case,

$$[M^{-1}]_{33} = -\alpha\beta - \frac{|A|}{A_{22}}.$$

Hence we obtain a unique solution,

$$\alpha = \frac{\delta^2}{1 - \beta\delta^2} \frac{|A|}{A_{22}}.$$

Since  $n_F = 0, n_g = 1, n_y = 2$ , condition (i) of the proposition holds if and only if  $\det M_2 = \det M > 0$ , and under the above solution for  $P_{11}$ ,  $\det M = -A_{22}$ ; hence condition (i) reduces to the requirement that  $A_{22} < 0$ .

This solution for  $P_{11}$ , and hence for  $M$ , also implies that

$$\Phi_{11} = \begin{bmatrix} \delta & 0 \\ -\delta A_{21}/A_{22} & 0 \end{bmatrix}.$$

Hence the eigenvalues of  $\Phi_{11}$  are 0 and  $\delta$ . Thus under our assumption about  $\delta$ , condition (ii) is necessarily satisfied, as long as  $A_{22} \neq 0$  (so that  $\Phi_{11}$  exists). We observe that both conditions (i) and (ii) hold if and only if  $A_{22} < 0$ , which is just the concavity condition derived above for the deterministic policy problem.

The solution for  $P_{11}$  similarly implies that

$$P_{22} = -G'_2 M^{-1} G_2 = -[M^{-1}]_{33} = \frac{1}{1 - \beta\delta^2} \frac{|A|}{A_{22}}.$$

Since the numerator in this last expression is positive, condition (iii) holds (in addition to the other two conditions) if and only if we also have  $\det A > 0$ . Since  $A$  is negative definite if and only if  $A_{22} < 0$  and  $\det A > 0$ , we can alternatively state that condition (iii) holds (in addition to the other two) if and only if  $A$  is also negative definite. This is the additional condition derived above for concavity in the case of stochastic policies.

### A.3 Lemma 3

**Lemma 3** *Suppose that under optimal policy, the extended state vector  $\mathbf{y}_t$  consists entirely of components that are either (i) stationary, or (ii) pure random walks. Suppose also that the class of policy rules  $\mathcal{R}$  is such that each rule in the class implies*

convergence to the same long-run values of the state variables as under optimal policy, in the absence of stochastic disturbances, so that the initial value of the trend component  $\mathbf{y}_{t_0-1}^{tr}$  is the same regardless of the rule  $r$  that is considered. Then for any rule  $r \in \mathcal{R}$ , the objective

$$\Omega_r(\mathbf{y}_{t_0-1}^{tr}) \equiv E_\mu[\bar{W}_r(\mathbf{y}_{t_0-1})], \quad (\text{A.22})$$

can be decomposed into two parts,

$$\Omega_r(\mathbf{y}_{t_0-1}^{tr}) = \Omega^1(\mathbf{y}_{t_0-1}^{tr}) + \Omega_r^2, \quad (\text{A.23})$$

where the first component is the same for all rules in this class, while the second component is independent of the initial condition  $\mathbf{y}_{t_0-1}^{tr}$ .

PROOF: We restrict attention to a class of rules  $\mathcal{R}$  with the property that each rule in the class implies convergence to the same long-run values of the state variables as under optimal policy, in the absence of stochastic disturbances. Because we analyze the dynamics under a given policy using a linearized version of the structural relations, certainty-equivalence obtains, and it follows that the limiting behavior (as  $T \rightarrow \infty$ ) of the long-run forecast  $E_{t_0}[\mathbf{y}_T]$  must also be the same under any rule  $r \in \mathcal{R}$ , given the initial conditions  $\mathbf{y}_{t_0-1}$ . Thus given these initial conditions, the decomposition of the initial extended state vector into components  $\mathbf{y}_{t_0-1}^{tr}$  and  $\mathbf{y}_{t_0-1}^{cyc}$  is the same under any rule  $r \in \mathcal{R}$ .

Let us consider the decomposition

$$\tilde{y}_t = \bar{y}_t + \hat{y}_t,$$

where  $\{\bar{y}_t\}$  is the deterministic sequence

$$\bar{y}_t \equiv E_{t_0-1}\tilde{y}_t$$

and  $\hat{y}_t$  is the component of  $\tilde{y}_t$  that is unforecastable as of date  $t_0 - 1$ . Then if we evaluate

$$\bar{W}(\tilde{y}; \mathbf{y}_{t_0-1}) \equiv E_{t_0-1}W(\tilde{y}; \xi_{t_0}, \mathbf{y}_{t_0-1}),$$

where  $W$  is the quadratic form defined in (4.7), under the evolution implied by any rule  $r$ , we find that

$$\bar{W}(\tilde{y}; \mathbf{y}_{t_0-1}) = \bar{W}(\bar{y}; \mathbf{y}_{t_0-1}) + \bar{W}(\hat{y}; \mathbf{y}_{t_0-1}). \quad (\text{A.24})$$



Here all the cross terms in the quadratic form have conditional expectation zero because  $\bar{y}$  is deterministic while  $\hat{y}$  is unforecastable.

Moreover, under any rule  $r$ , the value of  $\hat{y}_t$  is a linear function of the sequence of unexpected shocks between periods  $t_0$  and  $t$ , that is independent of the initial state. (This independence follows from the linearity of the law of motion (4.8), under the linear approximation that we use to solve for the equilibrium dynamics under a given policy rule.) Hence the second term on the right-hand side of (A.24),<sup>84</sup>

$$\bar{W}(\hat{y}; \mathbf{y}_{t_0-1}) = E_{t_0-1} V_{t_0}^Q(\hat{y}),$$

is independent of the initial state  $\mathbf{y}_{t_0-1}$  as well. Let  $\bar{W}_r^2$  denote the value of this expression associated with a given rule  $r$ .

Instead, the value of  $\bar{y}_t$  will be a linear function of  $\mathbf{y}_{t_0-1}$ , again as a result of the linearity of (4.8). And in our LQ problem with a self-consistent initial pre-commitment, the function (4.6) is linear as well. It follows that the first term on the right-hand side of (A.24) is a quadratic function of  $\mathbf{y}_{t_0-1}$ ,

$$\bar{W}(\bar{y}; \mathbf{y}_{t_0-1}) = \mathbf{y}'_{t_0-1} \Xi_r \mathbf{y}_{t_0-1},$$

where the subscript  $r$  indicates that the matrix of coefficients  $\Xi_r$  can depend on the policy rule that is chosen. Then substituting  $\mathbf{y}_{t_0-1}^{tr} + \mathbf{y}_{t_0-1}^{cyc}$  for  $\mathbf{y}_{t_0-1}$  in the above expression, and integrating over possible initial values of the cyclical component, for a given initial value of the trend component, we observe that

$$E_\mu[\bar{W}(\bar{y}; \mathbf{y}_{t_0-1})] = \mathbf{y}_{t_0-1}^{tr'} \Xi_r \mathbf{y}_{t_0-1}^{tr} + E_\mu[\mathbf{y}^{cyc'} \Xi_r \mathbf{y}^{cyc}], \quad (\text{A.25})$$

using the fact that  $E_\mu[\mathbf{y}^{cyc}] = 0$ .

Finally, we observe that under any rule  $r$ , the linearity of the law of motion (4.8) implies that conditional forecasts of the evolution of the endogenous variables take the form

$$E_{t_0-1} \mathbf{y}_T = \mathbf{y}_{t_0-1}^{tr} + B_{T+1-t_0} \mathbf{y}_{t_0-1}^{cyc},$$

where the sequence of matrices  $\{B_j\}$  may depend on the rule  $r$ , but the first term on the right-hand side is the same for all rules in the class  $\mathcal{R}$ . Using this solution for the sequence  $\bar{y}$  to evaluate  $\bar{W}(\bar{y}; \mathbf{y}_{t_0-1})$ , we find that the first term in (A.25) must

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<sup>84</sup>Here the expected value of the second term on the right-hand side of (4.7) vanishes because of the unforecastability of  $\hat{y}_{t_0}$ .

be a quadratic function of  $y_{t_0-1}^{tr}$  that is the same for all rules  $r$ , that can be denoted  $y_{t_0-1}^{tr'} \bar{\Xi} y_{t_0-1}^{tr}$ . Thus if we integrate (A.24) over the invariant distribution  $\mu$ , we obtain

$$E_\mu[\bar{W}_r(\mathbf{y}_{\mathbf{t}_0-1})] = y_{t_0-1}^{tr'} \bar{\Xi} y_{t_0-1}^{tr} + E_\mu[\mathbf{y}^{cyc'} \bar{\Xi}_r \mathbf{y}^{cyc}] + \bar{W}_r^2,$$

which is precisely a decomposition of the asserted form (A.23). This proves that the criterion (A.22) establishes the same ranking of alternative rules, regardless of the initial condition.

## A.4 Computing the Invariant Measure $\mu$

We need to know the invariant distribution  $\mu$  over possible initial conditions under optimal policy, in order to compute the proposed welfare criterion (4.12). Because  $\bar{W}_r(\cdot)$  is a quadratic function, we only need to compute the unconditional mean and variance-covariance matrix of  $\mathbf{y}_{\mathbf{t}}^{cyc}$  under optimal policy.

Substituting (3.19) for the pre-commitment  $\tilde{h}_{t+1}$  in the solution (3.17) for the optimal choice of  $\tilde{y}_{t+1}$ , we observe that under the solution to the recursive policy problem (and hence under the solution to the original problem as well),  $\tilde{y}_{t+1}$  is a linear function of  $\tilde{y}_t, \xi_{t+1}$ , and  $\xi_t$ , for each  $t \geq t_0$ . This solution together with the process (3.8) for the exogenous disturbances imply a law of motion of the form

$$\mathbf{y}_{t+1} = \bar{\Phi} \mathbf{y}_t + \bar{\Psi} \epsilon_{t+1} \tag{A.26}$$

for the extended state vector

$$\mathbf{y}_t \equiv \begin{bmatrix} \tilde{y}_t \\ \xi_t \end{bmatrix}. \tag{A.27}$$

Under this law of motion, the trend component of the extended state vector is given by  $\mathbf{y}_{\mathbf{t}}^{tr} = \Pi \mathbf{y}_{\mathbf{t}}$ , where  $\Pi$  is the matrix<sup>85</sup>

$$\Pi \equiv \lim_{j \rightarrow \infty} \bar{\Phi}^j,$$

and the cyclical component is correspondingly given by  $\mathbf{y}_{\mathbf{t}}^{cyc} = [I - \Pi] \mathbf{y}_{\mathbf{t}}$ . It then follows that the law of motion for the cyclical component is

$$\mathbf{y}_{\mathbf{t}+1}^{cyc} = \bar{\Phi} \mathbf{y}_{\mathbf{t}}^{cyc} + [I - \Pi] \bar{\Psi} \epsilon_{t+1}. \tag{A.28}$$

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<sup>85</sup>Under the assumption (made in the text) that the extended state vector is difference-stationary, this limit must be well-defined.

We note furthermore that (A.28) describes a jointly stationary set of processes, since the matrix  $\bar{\Phi}$  is stable on the subspace of vectors  $\mathbf{v}$  of the form  $\mathbf{v} = [I - \Pi]\mathbf{y}$  for some vector  $\mathbf{y}$ .<sup>86</sup> Hence there exist a well-defined vector of unconditional means  $\mathbf{E}$  and an unconditional variance-covariance matrix  $\mathbf{V}$ . The unconditional means are all zero, while the matrix  $V$  is given by the solution to the linear equation system

$$\mathbf{V} = \bar{\Phi}\mathbf{V}\bar{\Phi}' + [I - \Pi]\bar{\Psi}\Sigma\bar{\Psi}'[I - \Pi'].$$

In the case of some policy rules, it may be necessary to include additional lags of  $\tilde{y}_t$  or  $\xi_t$  in the extended state vector  $\mathbf{y}_t$ , in order for the equilibrium dynamics under the rule  $r$  to have a representation of the form (4.8). However, in this case, the additional elements of  $\mathbf{y}_t^{cyc}$  will all be lags of elements of the vector considered above. Hence the law of motion (A.28) can be used to derive the relevant unconditional moments in this case as well (though we omit the algebra).

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<sup>86</sup>When restricted to this subspace, the operator  $\bar{\Phi}$  has eigenvalues consisting of those eigenvalues of  $\bar{\Phi}$  that are less than one in modulus; these are in turn a subset of the eigenvalues of  $\Phi$  that are less than one in modulus (some zero eigenvalues have been dropped).

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