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ON STOCHASTIC OPTIMAL CONTROL PROBLEMS WITH SELECTION AMONG DIFFERENT COSTLY OBSERVATIONS

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Stéphane Lafortune

Memorandum No. UCB/ERL M85/99

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On Stochastic Optimal Control Problems

with Selection Among Different Costly Observations*

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Abstract

Optimizing observations is an important issue in the control of systems whose state is partially observed. This paper presents a general theorem for the computation of the optimal solution to discrete-time stochastic control problems, when the decision makers have the additional possibility of choosing at each step among different sets of observations on the system, each set incurring a different cost. Dynamic programming is employed to determine the optimal observations and controls. The result is applied to the special cases of: (i) finite-state controlled Markov chains, and (ii) linear Gaussian systems with a cost function quadratic in the states and controls.

Keywords: stochastic optimal control, observation selection, partial state information, dynamic programming, controlled Markov chains, LQG.

1. Problem statement

We consider the problem of the optimization of discrete-time stochastic systems, where at each step two consecutive decisions must be taken: (i) a decision on what type of observation to make on the system, and (ii) a decision on what control action to exert. The cost criterion depends on the state and these two control actions. We shall only consider finite-horizon problems.

Our aim is to present a general theorem for systems of the form:

$$x_{k+1} = f_k(x_k, m_k, u_k, w_k), (1.1a)$$

$$y_k = h_k(x_k, m_k, v_k),$$
 (1.1b)

for $k \ge 0$, with initial condition $x_0, x_0, w_0, \ldots, v_0, \ldots$, are mutually-independent

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random variables defined on an underlying probability space. Their probability distributions on R^n , R^n , and R^m , respectively, are known. $x_k \in R^n$ is the state. m_k and u_k are control variables taking values in $M \subset R^1$ and $U \subset R^m$. In particular, m_k parametrizes the observation equation (1.1b), where $y_k \in R^p$ is the observed process.

Let

$$I_k := \{ y_0, \ldots, y_k, m_0, \ldots, m_k, u_0, \ldots, u_k \}$$
 (1.2)

denote the information available to the decision maker at step k, $k \ge 0$. By convention, $I_{-1} = \emptyset$. At each step k, the values of the control actions m_k and u_k are determined by feedback in the following way:

$$m_k = g_k^{\ 1}(I_{k-1}) \in M \,, \tag{1.3}$$

$$u_k = g_k^2(I_{k-1}, m_k, y_k) \in U. (1.4)$$

Let the control strategy be denoted by $g = (g^1, g^2)$, where $g^i = \{g_0^i, \dots, g_{N-1}^i\}$, and let G denote the set of all admissible strategies. We define

$$J(g) := E^{g} \left[\sum_{k=0}^{N-1} c_{k}(x_{k}, m_{k}, u_{k}) + c_{N}(x_{N}) \right]$$
 (1.5)

to be the cost function associated with the control strategy g. The superscript g in the expectation emphasizes the fact that the stochastic processes x, m, u, and y become well-defined only when g is given. We want to find an optimal strategy $g^* \in G$, i.e., a strategy satisfying (a.s.)

$$J(g^*) = J^* := \inf_{g \in G} J(g)$$
. (1.6)

 m_k is an additional control variable parametrizing the observation equation. The choice among different sets of observations can be as simple as deciding whether or not to observe, in which case card(M) = 2. For the sake of generality, we shall also allow the

In the following, we will use the two notations $E^g c_k(x_k, m_k, u_k)$ and $E c_k(x_k^g, m_k^g, u_k^g)$ interchangeably.

possibility that the state equation depends on m_k . For example, in section 5, we consider linear Gaussian systems where the matrix C_k in $y_k = C_k x_k + v_k$ and the variances of the processes w and v depend on m_k . Another point is that at each step k, the decisions on m_k and u_k are made sequentially, and therefore u_k is allowed to depend on m_k , whereas the converse is not true. In short, this problem corresponds to optimizing the trade-off between the increased performance resulting from better observations (via better state estimates) and the higher cost of making better observations.

2. Information state for the system

We are dealing with a stochastic control problem with partial state information. We want to determine a suitable information state (in the terminology of Kumar and Varaiya [1]), or sufficient statistic (in that of Bertsekas [2]), for the system (1.1), i.e., a function of I_k that possesses the Markov property. For simplicity, we assume that densities exist. Let $p_0(x_0)$ denote the probability density (p.d.) of the initial condition x_0 , and, for a given control strategy g, let $p_{k+1}^g(x_k \mid I_{k-1})$ and $p_{k+1}^g(x_k \mid I_{k-1}, m_k, y_k)$ denote the conditional p.d. of x_k , given I_{k-1} and $I_{k-1} \cup \{m_k, y_k\}$, respectively.

Lemma 2.1: $p_{k+k-1}^g(\cdot \mid I_{k-1})$ is an information state for (1.1). It does not depend explicitly on g (and therefore we can drop the superscript g). There exists a function S_k such that

$$p_{k+1|k}(\cdot \mid I_k) = S_k[p_{k|k-1}(\cdot \mid I_{k-1}), m_k, y_k, u_k], \qquad (2.1)$$

with initial condition $p_{0l-1} = p_0$. S_k can be broken into two functions Φ_k and Ψ_k :

$$p_{k|k}(\cdot | I_{k-1}, m_k, y_k) = \Phi_k[p_{k|k-1}(\cdot | I_{k-1}), m_k, y_k]; \qquad (2.2)$$

$$p_{k+1|k}(\cdot | I_k) = \Psi_k[p_{k|k}(\cdot | I_{k-1}, m_k, y_k), m_k, u_k]. \tag{2.3}$$

² The organization of this paper and the proofs it contains were inspired by the treatment of standard stochastic systems (no m_k in (1.1)) in [1], chapters 2 to 6.

Proof: Given in the Appendix. \Box

The dynamics of the information state are illustrated in Figure 2.1.

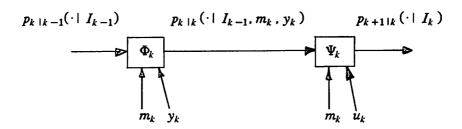


Figure 2.1 - Information State

For simplicity, we shall often denote $p_{k \mid k-1}(\cdot \mid I_{k-1})$ by $p_{k \mid k-1}$, and, similarly, $p_{k \mid k}(\cdot \mid I_{k-1}, m_k, y_k)$ by $p_{k \mid k}$. Observe that although the functions $p_{k \mid k-1}$ and $p_{k \mid k}$ do not depend explicitly on g, the processes x, m, u, and y, and consequently I_k , do depend on g. For this reason, it will sometimes be necessary to write I_{k-1}^g in the arguments of $p_{k \mid k-1}$ and $p_{k \mid k}$ to emphasize the strategy considered. Observe also that (2.2) and (2.3) imply that

$$p_{k+1|k+1}(\cdot \mid I_k, m_{k+1}, y_{k+1}) = T_{k+1}[p_{k|k}(\cdot \mid I_{k-1}, m_k, y_k), m_k, u_k, m_{k+1}, y_{k+1}],$$
 (2.4) and so $p_{k+1|k+1}$ is not an information state, because, due to the explicit dependence of x_k on m_k via (1.1a), m_k has to appear as an argument in T_{k+1} , even though it is already in $p_{k|k}$.

3. Optimal control

The sequentiality assumption for the decisions on m_k and u_k suggest that the optimal strategy g^* could be determined by a dynamic programming algorithm, where the dynamic programming equation (d.p.e.) would contain two nested minimizations. From Lemma 2.1, we expect that restricting attention to separated strategies is sufficient. Such strategies are of the form:

$$m_k = g_k^1(I_{k-1}) = g_k^1(p_{k|k-1}) \in M$$
; (3.1)

$$u_k = g_k^2(I_{k-1}, m_k, y_k) = g_k^2(p_{k+k-1}, m_k, y_k) = g_k^2(p_{k+k}, m_k) \in U.$$
 (3.2)

The following theorem shows that these claims are true. P denotes the set of all probability densities on \mathbb{R}^n . We define the cost-to-go from step k

$$J_{k}(g) := E\left[\sum_{j=k}^{N-1} c_{k}(x_{k}^{g}, m_{k}^{g}, u_{k}^{g}) + c_{N}(x_{N}^{g}) \mid I_{k-1}^{g}\right]. \tag{3.3}$$

Theorem 3.1: Define recursively the functions $V_k(p)$, $0 \le k \le N$, and $p \in P$, by:

$$V_N(p) := E[c_N(x_N) \mid p_{N+N-1} = p]; (3.4)$$

$$V_{k}(p) := \inf_{m \in M} E[\inf_{u \in U} E\{c_{k}(x_{k}, m, u) + V_{k+1}(\Psi[p_{k+k}, m, u]) \mid p_{k+k}, m\}$$

$$|p_{k|k-1} = p].$$
 (3.5)

(a) Consider any $g \in G$. Then

$$V_k(p_{k|k-1}(\cdot | I_{k-1}^g)) \le J_k(g) \text{ a.s., } 0 \le k \le N.$$
 (3.6)

(b) Let g^* be a separated policy such that for all $0 \le k \le N-1$ and for all $p \in P$, $g_k^*(p)$ achieves the infima in (3.5). Then

$$V_k(p_{k|k-1}(\cdot | I_{k-1}^*)) = J_k(g^*) \quad \text{a.s., } 0 \le k \le N,$$
(3.7)

and g^* is optimal. In particular, $V_0(p_0) = J^*$ a.s.

Proof: (a) The proof is by induction. Consider any $g \in G$. (3.6) is true with equality for k = N, because

$$J_{N}(g) = E[c_{N}(x_{N}^{g}) + I_{N-1}^{g}]$$

$$= \int_{x} c_{N}(x) p_{N+N-1}(x + I_{N-1}^{g}) dx$$

$$= V_{N}(p_{N+N-1}(\cdot + I_{N-1}^{g})). \tag{3.8}$$

by definition of p_{N+N-1} , and from (3.4). Now, suppose that (3.6) is true for k+1. We show that it is true for k, thus proving (a). Using successively the smoothing property of conditional expectations, (3.3), and the induction hypothesis, we get

$$J_{k}(g) = E^{g}[c_{k}(x_{k}, m_{k}, u_{k}) + E^{g}\{\sum_{j=k+1}^{N-1} c_{j}(x_{j}, m_{j}, u_{j}) + c_{N}(x_{N}) \mid I_{k}\} \mid I_{k-1}] \ a.s.$$

$$= E^{g}[c_{k}(x_{k}, m_{k}, u_{k}) + J_{k+1}(g) \mid I_{k-1}]$$

$$\geq E^{g}[c_{k}(x_{k}, m_{k}, u_{k}) + V_{k+1}(p_{k+1|k}(\cdot \mid I_{k})) \mid I_{k-1}] \ a.s.$$

$$= E^{g}[E^{g}\{c_{k}(x_{k}, m_{k}, u_{k}) + V_{k+1}(p_{k+1|k}(\cdot \mid I_{k}) \mid I_{k}\} \mid I_{k-1}] \ a.s.$$

$$(3.9)$$

But, by Lemma 2.1, we can replace the information sets by information states in (3.9):

$$J_{k}(g) \geqslant E^{g} \left[E^{g} \left\{ c_{k}(x_{k}, m_{k}, u_{k}) + V_{k+1}(\Psi_{k}[p_{k|k}(\cdot \mid I_{k-1}, m_{k}, y_{k}), m_{k}, u_{k}] \mid p_{k|k}, m_{k} \right\} \mid I_{k-1} \right] \text{ a.s.}$$

$$= E^{g} \left[E^{g} \left\{ c_{k}(x_{k}, m_{k}, u_{k}) + V_{k+1}(\Psi_{k}[p_{k|k}, m_{k}, y_{k}], m_{k}, u_{k}) \mid p_{k|k}, m_{k} \right\} \mid p_{k|k-1} \right]$$

$$\geqslant V_{k}(p_{k|k-1}(\cdot \mid I_{k-1}^{g})), \qquad (3.10)$$

the last inequality holding by (3.5).

(b) Again, we use induction to prove (3.7). First, we observe that (3.8) implies that (3.7) is true for k = N. Next, we repeat the development in (a), but with the given g^* in place of g. However, the two inequalities in (a) now become equalities: (3.9), by the induction hypothesis, and (3.10), because, by assumption, g_k^* achieves the infima in (3.5) for all $p \in P$. This proves (3.7). To show the optimality of g^* , we set k = 0 in (3.7) and (3.6) to get

$$J(g^*) = V_0(p(x_0)) \le J(g)$$
 a.s., for all $g \in G$. (3.11)

Remarks: (i) Observe that the V_{k+1} term could be removed from the inner conditional expectation in the d.p.e. (3.5).

(ii) (2.2) implies that, for each fixed m, the outer conditional expectation in (3.5) is an integral over y_k .

The argument of the value function V_k in (3.5) is a function, meaning that finding the optimal g^* is computationally difficult. In the next two sections, we consider two special cases where the problem is more amenable because the information state is finite-dimensional.

4. Special case I: finite-state controlled Markov chains

Consider a Markov chain whose state process x takes values in a finite set $S = \{1, 2, ..., S\}$, and whose transition-probability matrix P(m, u) can depend on two different controls m and u:

$$[\mathbf{P}(m,u)]_{i,j} = P_{ij}(m,u) := Prob(x_{k+1}=j \mid x_k=i, m_k=m, u_k=u).$$
 (4.1)

Let the observed process $y \in S$ be described by the output probability

$$P_j(i, m) := Prob(y_k = j \mid x_k = i, m_k = m).$$
 (4.2)

These probabilities do not depend on k. It is convenient to define the $S \times S$ matrix $\mathbf{D}(m, j)$ by

$$D(m, j) := diag[P_{j}(i, m)]_{i=1,...,S}.$$
(4.3)

Let $Prob_{k+k-1}(i+I_{k-1})$ and $Prob_{k+k}(i+I_{k-1},m_k,y_k)$ be the probabilities that $x_k=i$, given the respective information sets. Since the state space is finite, these probabilities are completely described by the $1\times S$ row-vectors:

$$\pi_{k|k-1}(I_{k-1}) := [Prob_{k|k-1}(1 \mid I_{k-1}), \dots, Prob_{k|k-1}(S \mid I_{k-1})]; \qquad (4.4)$$

$$\pi_{k|k}(I_{k-1}, m_k, y_k) := [Prob_{k|k}(1 \mid I_{k-1}, m_k, y_k), \dots, Prob_{k|k}(S \mid I_{k-1}, m_k, y_k)]. \quad (4.5)$$

To simplify the notation, we shall often omit writing the arguments of these two probabilities. Also, $\pi_{k+k-1}(j)$ will denote the jth component of π_{k+k-1} .

We write recursive relations for $\pi_{k|k}$ and $\pi_{k+1|k}$. The initial condition is $\pi_{0|-1} = \pi_0$, the given law of the initial state. It can be shown (cf: proof of Lemma 2.1, (A.8) and (A.11)) that the functions Φ_k and Ψ_k in (2.2) and (2.3) have the following

expression:

$$\pi_{k \mid k}(I_{k-1}, m_k, y_k) = \frac{\pi_{k \mid k-1}(I_{k-1}) D(m_k, y_k)}{\pi_{k \mid k-1}(I_{k-1}) D(m_k, y_k) \underline{1}};$$
(4.6)

$$\pi_{k+1|k}(I_k) = \pi_{k|k}(I_{k-1}, m_k, y_k) \mathbf{P}(m_k, u_k). \tag{4.7}$$

 $\underline{1}$ in (4.6) is the $S \times 1$ column-vector $(1, \ldots, 1)^T$.

We now write the complete expression of the d.p.e. (3.5). Consider $\pi \in \Pi$, the set of all $1 \times S$ probability row-vectors. Then, (3.4) and (3.5) become:

$$V_N(\pi) = \sum_{i \in S} c_N(i) \ \pi(i) \ , \tag{4.8}$$

$$V_{k}(\pi) = \inf_{m \in M} E \left[\inf_{u \in U} \left\{ \sum_{i \in S} c_{k}(i, m, u) \pi_{k|k}(i) + V_{k+1}(\pi_{k|k} | P(m, u)) \right\} \mid \pi_{k|k-1} = \pi \right]. (4.9)$$

Let $u_k^* = g_k^{2^*}(\pi_{k|k}, m) = g_k^{2^*}(\pi_{k|k-1}, m, y_k)$ achieve the inner infimum. We evaluate the conditional expectation in (4.9):

$$V_{k}(\pi) = \inf_{m \in M} E\left[\sum_{i \in S} c_{k}(i, m, u_{k}^{*}) \pi_{k \mid k}(i) + V_{k+1}(\pi_{k \mid k} | \mathbf{P}(m, u_{k}^{*})) \mid \pi_{k \mid k-1} = \pi\right]$$

$$=\inf_{m\in M}\int_{y_k} \left[\sum_{i\in S} c_k(i,m,u_k^*) \frac{\pi D(m,y_k)}{\pi D(m,y_k) \underline{1}}(i)\right]$$

+
$$V_{k+1}(\frac{\pi D(m, y_k)}{\pi D(m, y_k)} \mathbb{P}(m, u_k^*))] \pi D(m, y_k) \underline{1} dy_k$$
, (4.10)

since

$$Prob(y_k \mid \pi_{k \mid k-1} = \pi, m) = \sum_{i \in S} Prob(y_k \mid i, m) Prob(x_k = i \mid \pi_{k \mid k-1} = \pi)$$
 (4.11)

$$=\sum_{i\in S}P_{y_k}(i,m)\pi(i)$$

$$= \pi \mathbf{D}(m, y_k) \underline{1}, \qquad (4.12)$$

where in (4.11) we have used (3.1).

5. Special case II: linear Gaussian systems

Consider the case where (1.1) is of the form:

$$x_{k+1} = A_k x_k + B_k u_k + w_k , (5.1a)$$

$$y_k = C_k(m_k)x_k + v_k$$
, (5.1b)

with $x_0 \sim N(\bar{x}_0, \Sigma_0)$, $w_k \sim N(0, Q_k(m_k))$, and $v_k \sim N(0, R_k(m_k))$. Consider a cost-function quadratic in the states and in the controls u:

$$J(g) := E^g \left[\sum_{k=0}^{N-1} (x_k^T M_k x_k + u_k^T N_k u_k + c_k (m_k)) + x_N^T M_N x_N \right]. \tag{5.2}$$

(Here, we make the usual symmetry and positive (semi-)definiteness assumptions on M_k , N_k , Q_k , and R_k .) We mention at this point that Aoki and Li [3] have studied a version of this problem where the decision makers select the total *number* and the *spacings* between each observation, which is quite different from our formulation. Also, their model has no w term in (5.1a).

The derivation of the Kalman filter remains valid when the matrices A_k , B_k , and C_k are random, provided that they are measured at time k, i.e., that they are in I_k , and that they are independent of the noise variables (e.g., [1]). In our case, once m_k is chosen, all the parameters in (5.1) that depend on it can simply be regarded as time-varying, with the important difference that their time variation can be altered. However, that decision is based on past information, namely, I_{k-1} . It follows that the p.d. $p_{k+1|k}$ and $p_{k|k}$ defined in section 2 are Gaussian, and therefore the information state $p_{k+1|k}$ is two-dimensional. In fact, this remains true if A_k and B_k also depend on m_k .

Consider a fixed feedback strategy g and the corresponding processes x, m, u, and y. Then, using the notation

$$p_{k+1|k}(x_{k+1} | I_k) \sim N(\hat{x}_{k+1|k}, \Sigma_{k+1|k})$$
 (5.3)

³ For simplicity, we omit writing the superscript g for these processes.

$$p_{k|k}(x_k | I_{k-1}, m_k, y_k) \sim N(\hat{x}_{k|k}, \Sigma_{k|k}), \qquad (5.4)$$

the Kalman filter equations corresponding to (2.2) and (2.3) are:

$$\hat{x}_{k|k} = \hat{x}_{k|k-1} + L_k (y_k - C_k (m_k) \hat{x}_{k|k-1})$$
(5.5)

$$\hat{x}_{k+1|k} = A_k \, \hat{x}_{k|k} + B_k \, u_k \tag{5.6}$$

$$\Sigma_{k|k} = \Sigma_{k|k-1} - L_k C_k (m_k) \Sigma_{k|k-1}$$
 (5.7)

$$\Sigma_{k+1|k} = A_k \, \Sigma_{k|k} \, A_k^T + Q_k \, (m_k) \tag{5.8}$$

where

$$L_k := \sum_{k|k-1} C_k (m_k)^T [C_k (m_k) \sum_{k|k-1} C_k (m_k)^T + R_k (m_k)]^{-1}.$$
 (5.9)

The interesting feature of this special case is that, due to the quadratic form of the cost and the fact that only R_k in (5.1a) depends on m_k , the certainty-equivalence principle still holds and the value function has a partially-closed form. This is not true for linear Gaussian systems in general, and this was our motivation for these extra assumptions. More precisely, it can be shown, by substituting (5.10) in (3.5), that

$$V_k(\hat{x}_{k|k-1}, \Sigma_{k|k-1}) = \hat{x}_{k|k-1}^T P_k \hat{x}_{k|k-1} + W_k(\Sigma_{k|k-1}), \qquad (5.10)$$

 $0 \le k \le N$. P_k is determined by solving the standard backward Riccati equation

$$P_k = M_k + A_k^T P_{k+1} A_k - K_k^T (N_k + B_k^T P_{k+1} B_k) K_k , \qquad (5.11)$$

 $0 \le k < N$, with final condition $P_N = M_N$. K_k is the deterministic optimal control gain

$$K_k := -[N_k + B_k^T P_{k+1} B_k]^{-1} B_k^T P_{k+1} A_k , \qquad (5.12)$$

i.e., $u_k^* = K_k \hat{x}_{k \mid k} = g_k^{2^*}(\hat{x}_{k \mid k})$. P_k and K_k do not depend on m and that they can be completely determined beforehand. The other part of V_k has no closed-form solution and must be solved recursively as follows:

$$W_{k}(\Sigma) = \inf_{m \in M} [c_{k}(m) + Trace\{M_{k} \Sigma + l_{k}(\Sigma, m)(P_{k} - M_{k})\} + W_{k+1}(A_{k} \Sigma A_{k}^{T} + Q_{k}(m) - A_{k} l_{k}(\Sigma, m) A_{k}^{T})].$$
(5.13)

with final condition $W_N(\Sigma) = Trace\{M_N\Sigma\}$, where we have defined

$$l_k(\Sigma, m) := \Sigma C_k(m)^T [C_k(m)\Sigma C_k(m)^T + R_k(m)]^{-1} C_k(m)\Sigma.$$
 (5.14)

The optimal sequence m^* can be determined beforehand, but it depends on P_k , and consequently the Riccati equation must be solved first. If M is finite with card(M) = n, then at step k, the domain of Σ in (5.13) can contain up to n^k values.

As in the standard LQG problem, the control u has no learning role, but the control m has one, since it can influence the estimation covariance of the state. Clearly, if A_k or B_k were dependent on m_k , V_k would possess no separation property as (5.10) exhibits, even with a quadratic cost. Thus, u_k^* would in general also depend on Σ_{k+1} , meaning that it too would have a learning function.

Finally, we point out that the problem in this section was also treated by Deissenberg and Stöppler [4] for the special case when card(M) = 2, corresponding to the decision: observe / do not observe. However, the value function considered in that paper does not have the estimation covariance matrix Σ_{k+k-1} as an argument, and therefore the solution does not have the clear recursive form of (5.10) and (5.13), which also provides for more computational efficiency.

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Appendix - Proof of Lemma 2.1

The independence of all the noise variables in (1.1) and the fact that the values of m and u are measured imply that

$$p^{g}(x_{k+1} \mid x_{k}, I_{k}) = p(x_{k+1} \mid x_{k}, m_{k}, u_{k}), \qquad (A.1)$$

$$p^{g}(y_{k} \mid x_{k}, m_{k}, I_{k-1}) = p(y_{k} \mid x_{k}, m_{k}),$$
 (A.2)

where the densities on the right-hand sides do not depend on the strategy g, but only on the values of m and u.

We now establish the precise form of the recursive relations (2.2) and (2.3).

$$p_{k|k}^{g}(x_{k} \mid I_{k-1}, m_{k}, y_{k}) = \frac{p^{g}(y_{k} \mid x_{k}, I_{k-1}, m_{k}) p^{g}(x_{k}, I_{k-1}, m_{k})}{p^{g}(I_{k-1}, m_{k}, y_{k})}$$
(A.3)

$$= \frac{p(y_k \mid x_k, m_k) p^g(x_k, I_{k-1}, m_k)}{\int\limits_{x_k} p^g(x_k, I_{k-1}, m_k, y_k) dx_k}.$$
 (A.4)

where (A.3) follows from Bayes' rule, and (A.4) by using (A.2). But

$$p^{g}(x_{k}, I_{k-1}, m_{k}, y_{k}) = p^{g}(y_{k} \mid x_{k}, I_{k-1}, m_{k}) p^{g}(x_{k}, I_{k-1}, m_{k})$$
(A.5)

$$= p(y_k \mid x_k, m_k) p^{g}(x_k \mid I_{k-1}, m_k) p^{g}(I_{k-1}, m_k)$$
(A.6)

by (A.2). Substituting (A.6) in (A.4),

$$p_{k|k}^{g}(x_{k} | I_{k-1}, m_{k}, y_{k}) = \frac{p(y_{k} | x_{k}, m_{k}) p^{g}(x_{k} | I_{k-1}, m_{k})}{\int_{x_{k}} p(y_{k} | x_{k}, m_{k}) p^{g}(x_{k} | I_{k-1}, m_{k}) dx_{k}}$$
(A.7)

$$= \frac{p(y_k \mid x_k, m_k) p_{k+k-1}^g(x_k \mid I_{k-1})}{\int\limits_{x_k} p(y_k \mid x_k, m_k) p_{k+k-1}^g(x_k \mid I_{k-1}) dx_k}.$$
 (A.8)

because m_k is a function of I_{k-1} (see (1.3)) and x_k only depends on m_k via I_{k-1} . (A.8) is of the form given in (2.2). Next,

$$p_{k+1|k}^{g}(x_{k+1} + I_k) = \int_{x_k} p^{g}(x_{k+1} + x_k, I_k) p^{g}(x_k + I_k) dx_k$$
(A.9)

$$= \int_{x_k} p(x_{k+1} \mid x_k, m_k, u_k) p^{g}(x_k \mid I_{k-1}, m_k, y_k, u_k) dx_k$$
 (A.10)

$$= \int_{x_k} p(x_{k+1} \mid x_k, m_k, u_k) p_{k+k}^g(x_k \mid I_{k-1}, m_k, y_k) dx_k . \tag{A.11}$$

(A.10) is a consequence of (A.1). (A.11) is true because x_k does not depend explicitly on u_k but only through $I_{k-1} \cup \{m_k, y_k\}$, of which u_k is a function (see (1.4)). (A.11) corresponds to (2.3). $p_{k|k}$ and $p_{k+1|k}$ do not depend on g because the functions Φ and Ψ

of (A.8) and (A.11) do not, and the initial condition is $p_{0,-1} = p_0$ (recall that $I_{-1} = \emptyset$) and is therefore independent of g. \square

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