Dynamic Models with Latent Variables

Jean-Michel Zakoian

CREST-ENSAE

Markov-switching models

 $\begin{array}{c} {\rm Hidden\ Markov\ Model} \\ {\rm MS-ARMA}(p,q) \ {\rm process} \\ {\rm Estimation\ of\ MS-AR\ models} \end{array}$

- 1 Hidden Markov Model
- 2 MS-ARMA(p,q) process
- 3 Estimation of MS-AR models

Time series and breaks

- Many economic time series occasionally exhibit dramatic breaks in their behavior. Such breaks may concern the level, the variability, the serial correlations..
- They are often associated with events such as financial crises or abrupt changes in government policy. For instance, many economic variables tend to behave differently during economic downturns.
- When models are fitted over sub-periods of long times series, one often detects significant changes in the estimated parameter values.

How can we model dramatic changes?

Until the end of the 80's, the main approach consisted in fitting different models over different subperiods.

For instance, for AR(1) models:

$$y_t = \phi_1 y_{t-1} + \epsilon_t$$
, if $t \le t_0$,
 $y_t = \phi_2 y_{t-1} + \epsilon_t$, if $t > t_0$.

Drawbacks:

- Sub-models are not related
- Arbitrary choice of the break date t_0
- Assumption of non stationarity

A simple way to model dramatic changes

AR(1) model with random AR coefficient:

$$y_t = \phi(\Delta_t)y_{t-1} + \epsilon_t,$$

where Δ_t is a discrete random variable.

A complete description of the probability law governing the data would then require a probabilistic model for Δ_t .

The simplest such specification is that Δ_t is the realization of a finite-state Markov chain.

Such a model is called Markov-Switching (MS) AR(1) model.

References

The introduction of **MS models** in the econometrics literature is due to Hamilton

 Hamilton, J. D. (1989) A new approach to the economic analysis of nonstationary time series and the business cycle. Econometrica 57, 357-384.

In the statistics literature, related models called Hidden Markov Models (HMM) had been studied much earlier:

- Baum, L. E., et T. Petrie (1966) Statistical inference for probabilistic functions of finite state Markov chains. *Annals of Mathematical Statistics* 30, 1554-1563.
- Baum, L. E., Petrie, T., Soules, G. et Weiss, N. (1970) A maximization technique in the statistical analysis of probabilistic functions of Markov chains. *Annals of Mathematical statistics* 41, 164-171.

- 1 Hidden Markov Model
 - Finite-state Markov chain:
 - Properties of Hidden Markov Models
- 2 MS-ARMA(p,q) process
- 3 Estimation of MS-AR models

Definition

Let $\Delta_0, \Delta_1, \dots \in \mathcal{E} = \{1, \dots, d\}$ a sequence of random variables. The sequence is a *Markov chain* if

$$P(\Delta_t = j | \Delta_{t-1} = i)$$
= $P(\Delta_t = j | \Delta_{t-1} = i, \Delta_{t-2} = e_{t-2}, ..., \Delta_0 = e_0) = p(i, j)$

for any t and any $(i,j,e_{t-2},\ldots,e_0)\in\mathcal{E}^{t+1}$ for which $P(\Delta_{t-1}=i,\Delta_{t-2}=e_{t-2},\ldots,\Delta_0=e_0)>0$.

The set $\mathscr E$ is called the state space of the process and the p(i,j) are the transition probabilities.

Law

The law of the Markov chain is entirely characterized by

1 the initial probabilities

$$\pi_0(i) = P(\Delta_0 = i), \quad \pi_0(i) \ge 0, \qquad i = 1, \dots, d, \quad \sum_{i=1}^d \pi_0(i) = 1$$

2 the transition probability matrix $P = (p(i,j))_{1 \le i,j \le d}$.

Higher-order transitions

The kth power of the transition matrix, $P^k = (p^{(k)}(i,j))_{1 \le i,j \le d}$, provides the k-step transition probabilities :

$$p^{(k)}(i,j) = P(\Delta_t = j | \Delta_{t-k} = i), \quad i, j \in \mathcal{E}, k \ge 0.$$

Let

$$\pi_0 = \left(\begin{array}{c} \pi_0(1) \\ \vdots \\ \pi_0(d) \end{array} \right) \quad \text{and} \quad \pi_n = \left(\begin{array}{c} P(\Delta_n = 1) \\ \vdots \\ P(\Delta_n = d) \end{array} \right).$$

We have

$$\pi_n = P' \pi_{n-1}, \quad \pi_n = P^{'n} \pi_0, \quad n \ge 0.$$

Invariant probability

A probability π on $\mathscr E$ is called invariant probability if

$$\pi = P'\pi$$
, $\pi' \mathbf{1} = 1$ (with $\mathbf{1}' = (1, ..., 1)$).

- If the limit law $\pi_{\infty} := \lim_{n \to \infty} \pi_n$ exists then it is an invariant probability.
- An invariant probability always exists (for a finite state space).
- If $\pi_0 = \pi$ where π is an invariant probability, then $\pi_n = \pi$ for all $n \ge 0$.

Irreducibility, aperiodicity

It is possible for a chain starting in i to reach j if and only if

$$p^{(n)}(i,j) > 0$$
, for some n .

If it is true for all i and j, the Markov chain is called irreducible.

A state i is called aperiodic if

$$1 = \gcd\{n; \quad p^{(n)}(i, i) > 0\}.$$

If all states verify this condition, the chain is aperiodic.

Exponential convergence to the stationary law and ergodicity

Proposition

If the chain is irreducible and aperiodic, there is a stationary distribution π and there exists $K \ge 0$ and $0 < \rho < 1$ such that

$$|p^{(n)}(i,j) - \pi(j)| \le K\rho^n,$$

for all states i and j.

Under these conditions, we have the ergodic property:

$$\frac{1}{n}\sum_{t=1}^{n}f(\Delta_{t})\stackrel{n\to\infty}{\to}E_{\pi}\{f(\Delta_{t})\}=\sum_{i=1}^{d}\pi_{i}f(i),\quad a.s.$$

for any function f.

An irreducible, aperiodic and stationary Markov chain is called **ergodic**.

Definition

A process $(X_t)_{t\geq 0}$ follows a HMM if

- ① conditionally to a (hidden) Markov chain (Δ_t) , the variables X_0, X_1, \ldots are independent;
- 2 the conditional law of X_s given (Δ_t) only depends on Δ_s .

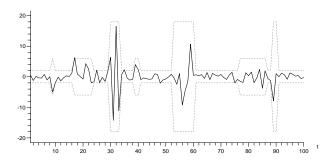
Canonical HMM:

$$\epsilon_t = \sigma(\Delta_t)\eta_t$$

where

- $0 < \sigma(1) < ... < \sigma(d)$,
- (η_t) is an iid sequence of variables, $E(\eta_t) = 0$, $var(\eta_t) = 1$,
- (Δ_t) is an **ergodic Markov chain** on $\mathscr{E} = \{1, ..., d\}$.
- the sequences (η_t) and (Δ_t) are independent.

Simulations of length 100 of a HMM



3 regimes Model with $\eta_t \sim \mathcal{N}(0,1)$ and

$$\sigma(1) = 1$$
, $\sigma(2) = 3$, $\sigma(3) = 9$, $P = \begin{pmatrix} 0.85 & 0.1 & 0.05 \\ 0.3 & 0.7 & 0 \\ 0.3 & 0 & 0.7 \end{pmatrix}$.

Unconditional law

If $\eta_t \sim \mathcal{N}(0,1)$, the law of ϵ_t is a mixture of centered Gaussian distributions : its density is given by

$$f(x) = \sum_{i=1}^{d} \pi(i) \frac{1}{\sigma(i)} \phi\left(\frac{x}{\sigma(i)}\right).$$

The law is **not Gaussian** (except when $\sigma(i) = \sigma$ for all i).

For any law of η_t , the marginal moments of ϵ_t can be obtained : for any r > 0,

$$E(\boldsymbol{\epsilon}_t^r) = E\boldsymbol{\sigma}^r(\boldsymbol{\Delta}_t)E(\boldsymbol{\eta}_t^r) = \sum_{i=1}^d \boldsymbol{\sigma}^r(i)\pi(i)E(\boldsymbol{\eta}_t^r).$$

In particular, $E(\epsilon_t) = 0$.

Correlations

We have

$$Corr(\epsilon_t, \epsilon_{t-k}) = 0$$
, for all $k > 0$.

Thus (ϵ_t) is a white noise with variance

$$E\epsilon_t^2 = \sum_{i=1}^d \sigma^2(i)\pi(i).$$

The dependence cannot be seen on the 2nd order structure.

Correlations of squares : case d=2

- Eigenvalues of P: 1 and $\lambda = p(1,1) + p(2,2) 1$.
- $-1 < \lambda < 1$ because the chain is irreducible and aperiodic.
- ullet The entries of P^k have the form

$$p^{(k)}(i,j) = a_1(i,j) + a_2(i,j)\lambda^k, \quad k \ge 0.$$

- We have $a_1(i,j) = \pi(j)$, and $a_1(i,j) + a_2(i,j) = \mathbf{1}_{\{i=j\}}$.
- Thus, for j = 1, 2 and $i \neq j$

$$p^{(k)}(i,j) = \pi(j)(1-\lambda^k), \quad p^{(k)}(j,j) = \pi(j) + \lambda^k(1-\pi(j)),$$

and, for $i, j = 1, 2, k \ge 0$

$$p^{(k)}(i,j) - \pi(j) = \lambda^k \left[\left\{ 1 - \pi(j) \right\} \mathbf{1}_{\{i=j\}} - \pi(j) \mathbf{1}_{\{i\neq j\}} \right].$$

Correlations of squares : case d=2

We have, for k > 0,

$$\begin{split} \mathsf{Cov}(\epsilon_t^2, \epsilon_{t-k}^2) &= E\left\{\sigma^2(\Delta_t)\sigma^2(\Delta_{t-k})\right\} - \{E\sigma^2(\Delta_t)\}^2 \\ &= \sum_{i,j=1}^2 p^{(k)}(i,j)\pi(i)\sigma^2(i)\sigma^2(j) - \left\{\sum_i^2 \pi(i)\sigma^2(i)\right\}^2 \\ &= \sum_{i,j=1}^2 \{p^{(k)}(i,j) - \pi(j)\}\pi(i)\sigma^2(i)\sigma^2(j) \\ &= \lambda^k \left\{\sum_{j=1}^2 (1 - \pi(j))\pi(j)\sigma^4(j) - \sum_{i\neq j} \pi(i)\pi(j)\sigma^2(i)\sigma^2(j)\right\} \\ &= \lambda^k \{\sigma^2(1) - \sigma^2(2)\}^2 \pi(1)\pi(2). \end{split}$$

Correlations of squares : case d=2

$$\mathsf{Cov}(\epsilon_t^2,\epsilon_{t-k}^2) \quad = \quad \lambda^k \{\sigma^2(1) - \sigma^2(2)\}^2 \pi(1)\pi(2), \qquad k > 0.$$

- Covariances have the sign of λ^k .
- The dependence increases with $|\lambda| = |p(1,1) + p(2,2) 1|$ and with $|\sigma^2(1) \sigma^2(2)|$.
- (ϵ_t^2) is an ARMA(1,1) process, since

$$\gamma_{\epsilon^2}(k) = \lambda \gamma_{\epsilon^2}(k-1), \quad k > 1.$$

A similar computation shows that

$$\mathsf{Var}(\epsilon_t^2) = \{\sigma^2(1) - \sigma^2(2)\}^2 \pi(1)\pi(2) + \{\sigma^4(1)\pi(1) + \sigma^4(2)\pi(2)\} \mathsf{Var}(\eta_t^2).$$

Correlations of squares: general case

The transition matrix P may not be diagonalizable but 1 remains an eigenvalue, with corresponding eigenspace of dimension 1 (otherwise there would exist several invariant distributions). Let $\lambda_1, \dots, \lambda_m$ the eigenvalues different from 1 and n_1, \dots, n_m the dimensions of the corresponding eigenspace $(n_1 + \cdots + n_m = d - 1)$. We use the Jordan representation

$$P = SJS^{-1}$$
, where S is nonsingular,

$$P = SJS^{-1}, \qquad \text{where } S \text{ is nonsingular},$$

$$J = \begin{pmatrix} J_{n_1}(\lambda_1) & 0 & \dots & 0 \\ 0 & J_{n_2}(\lambda_2) & 0 & & \vdots \\ \vdots & & \ddots & \ddots & \vdots \\ \vdots & & & \ddots & \ddots & \vdots \\ \vdots & & & & \ddots & \ddots & \vdots \\ \vdots & & & & \ddots & J_{n_m}(\lambda_m) & 0 \\ 0 & & \dots & & 0 & 1 \end{pmatrix}, \qquad J_l(\lambda) = \lambda I_l + N_l(1)$$

where $N_l(i)$ the $l \times l$ matrix whose elements are null, except the i-th superdiagonal filled with 1.

Correlations of squares : general case

We have $J_l^k(\lambda) = \lambda^k P^{(l)}(k)$, for some polynomial $P^{(l)}$ of degree l-1,

$$\Rightarrow P^{k} = S \begin{pmatrix} \lambda_{1}^{k} P^{(n_{1})}(k) & 0 & \dots & 0 \\ 0 & \lambda_{2}^{k} P^{(n_{2})}(k) & 0 & & \vdots \\ \vdots & & \ddots & \ddots & \vdots \\ \vdots & & & \lambda_{m}^{k} P^{(n_{m})}(k) & 0 \\ 0 & \dots & & 0 & 1 \end{pmatrix} S^{-1}.$$

It follows that

$$p^{(k)}(i,j) = \pi(j) + \sum_{l=1}^{m} \lambda_l^k p_{i,j}^{(n_l)}(k), \qquad d^o(p_{i,j}^{(n_l)}) = n_l - 1$$

and, for some polynomials $q^{(n_l)}(k)$ of degree n_l-1 ,

$$\mathsf{Cov}(\epsilon_t^2, \epsilon_{t-k}^2) = \sum_{i,j=1}^d \sum_{l=1}^m \lambda_l^k p_{i,j}^{(n_l)}(k) := \sum_{l=1}^m \lambda_l^k q^{(n_l)}(k), \quad k > 0.$$

Stationarity of the MS-AR(1) model Stationarity of the MS-ARMA(p,q) model Examples

- 1 Hidden Markov Model
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 - Examples
- 3 Estimation of MS-AR models

Definition

The MS-ARMA(p,q) model is defined by

$$\left\{ \begin{array}{lcl} X_t & = & c(\Delta_t) + \sum_{i=1}^p a_i(\Delta_t) X_{t-i} + \epsilon_t + \sum_{j=1}^q b_j(\Delta_t) \epsilon_{t-j}, \\ \\ \epsilon_t & = & \epsilon_t(\Delta_t) = \sigma(\Delta_t) \eta_t, \qquad (\eta_t) \stackrel{iid}{\sim} (0,1), \end{array} \right.$$

where (Δ_t) is an ergodic (aperiodic, irreducible, stationary) Markov chain on $\mathscr{E} = \{1, 2, ..., d\}$, which is independent of (η_t) , $a_i(\cdot), b_i(\cdot), c(\cdot) \in \mathbb{R}, \ \sigma(\cdot) > 0$.

This model contains the HMM as a particular case.

Except in the case p = 0 (MS-MA(q)), the existence of stationary solutions require additional conditions.

Notations

For any function $f: \mathscr{E} \to \mathscr{M}_{n \times n'}(\mathbb{R})$, where $\mathscr{M}_{n \times n'}(\mathbb{R})$ is the space of real $n \times n'$ matrices, and for all positive integers i, n et n', let

$$\mathbb{P}^{(i)}(f) = \begin{pmatrix} p^{(i)}(1,1)f(1) & \cdots & p^{(i)}(d,1)f(1) \\ \vdots & & \vdots \\ p^{(i)}(1,d)f(d) & \cdots & p^{(i)}(d,d)f(d) \end{pmatrix}, \qquad \Pi(f) = \begin{pmatrix} \pi(1)f(1) \\ \vdots \\ \pi(d)f(d) \end{pmatrix}$$

- When i = 1, write $\mathbb{P}(f) = \mathbb{P}^{(1)}(f)$,
- for $f \equiv 1$, let $\mathbb{P} = \mathbb{P}(1) = (p(j, i))$, the transpose of the transition matrix.

A useful result for computing expectations

Lemma

Let $f_0, ..., f_k$ functions $\mathcal{E} \mapsto \mathcal{M}_{n \times n}(\mathbb{R})$. For k > 0,

$$E\{f_0(\Delta_t)f_1(\Delta_{t-1})\dots f_k(\Delta_{t-k})\} = \boldsymbol{I}\mathbb{P}(f_0)\dots \mathbb{P}(f_{k-1})\Pi(f_k)$$

where $I = (I_n, ..., I_n)$ is a $n \times nd$ matrix and I_n is the identity matrix of size n.

Proof (k=1). We have

$$E\{f_{0}(\Delta_{t})f_{1}(\Delta_{t-1})\} = E[E\{f_{0}(\Delta_{t}) \mid \Delta_{t-1}\}f_{1}(\Delta_{t-1})]$$

$$= \sum_{j=1}^{d} E\{f_{0}(\Delta_{t}) \mid \Delta_{t-1} = j\}f_{1}(j)\pi(j)$$

$$= \sum_{j=1}^{d} \sum_{i=1}^{d} f_{0}(i)p(j,i)f_{1}(j)\pi(j) = I\mathbb{P}(f_{0})\Pi(f_{1}).$$

MS-AR(1) without intercept

$$X_t = a(\Delta_t)X_{t-i} + \sigma(\Delta_t)\eta_t$$

Problems:

- Existence of a strictly stationary solution.
- Existence of a 2nd-order stationary solution.

We are interested in **non-anticipative solutions**, i.e. solutions of the form

$$X_t = f(\eta_t, \Delta_t, \eta_{t-1}, \Delta_{t-1}, \ldots).$$

MS-AR(1) without intercept

By successive replacements, for $k \ge 1$,

$$X_t = a(\Delta_t) \dots a(\Delta_{t-k+1}) X_{t-k}$$

+
$$\sum_{i=0}^{k-1} a(\Delta_t) \dots a(\Delta_{t-i+1}) \sigma(\Delta_{t-i}) \eta_{t-i}$$

(with by convention $a(\Delta_t) \dots a(\Delta_{t-i+1}) = 1$ if i = 0).

Solutions should be given by

$$\tilde{X}_t = \sum_{n=0}^{\infty} a(\Delta_t) \dots a(\Delta_{t-n+1}) \sigma(\Delta_{t-n}) \eta_{t-n},$$

provided the series converges almost surely.

Cauchy root test

To derive an absolute convergence condition, we will use the nth root (Cauchy) test.*

Let $u_n = a(\Delta_t) \dots a(\Delta_{t-n+1}) \sigma(\Delta_{t-n}) \eta_{t-n}$. We have

$$|u_n|^{1/n} = \exp\left\{\frac{1}{n}\sum_{k=1}^n \log|a(\Delta_{t-k+1})| + \frac{1}{n}\log\{\sigma(\Delta_{t-n})|\eta_{t-n}|\}\right\}.$$

Because $\overline{\lim} \ n^{-1} \log \{\sigma(\Delta_{t-n}) | \eta_{t-n}| \} = 0$ a.s., we have, by the ergodic theorem

$$\overline{\lim} |u_n|^{1/n} = \exp\{E\log|a(\Delta_t)|\}.$$

 \Rightarrow Condition : $E\log|a(\Delta_t)| < 0$.

^{*.} For a non-negative sequence (u_n) , the series $\sum u_n$ converges if $\overline{\lim} \sqrt[n]{u_n} < 1$.

Uniqueness

Under the previous condition,

$$\tilde{X}_t = \sum_{n=0}^{\infty} a(\Delta_t) \dots a(\Delta_{t-n+1}) \sigma(\Delta_{t-n}) \eta_{t-n}, \quad a.s$$

is well defined. Moreover it is a solution of the MS-AR(1) model. Now suppose there exists another strictly stationary solution (X_t^*) .

$$X_t^* = a(\Delta_t) \dots a(\Delta_{t-k+1}) X_{t-k}^* + \sum_{i=0}^{k-1} a(\Delta_t) \dots a(\Delta_{t-i+1}) \sigma(\Delta_{t-i}) \eta_{t-i}.$$

Hence

$$|X_t - X_t^*| \leq |a(\Delta_t) \dots a(\Delta_{t-k+1})||X_{t-k}^*| + r_{t,k},$$

where $r_{t,k} \to 0$ and $|a(\Delta_t) \dots a(\Delta_{t-k+1})| \to 0$ a.s. as $k \to \infty$. Hence $X_t = X_t^*$ a.s.

Strict stationarity condition

Proposition

There exists a unique strictly stationary solution if

$$E\log|a(\Delta_t)| = \sum_{i=1}^d \log|a(i)|\pi(i) < 0.$$

Moreover, this solution is non anticipative.

Strict stationarity condition

Remarks:

0

$$\sum_{i=1}^{d} \log |a(i)| \pi(i) < 0 \quad \Leftrightarrow \quad \prod_{i=1}^{d} |a(i)|^{\pi(i)} < 1.$$

- The stationarity condition of an AR(1), |a| < 1, has to be satisfied "in average" over the different regimes. Regimes with |a(i)| > 1 are allowed (but they must not be visited too often).
- Only depends on the stationary probabilities of the Markov chain, not on the transition probabilities.
- If a(i) = 0 for at least one regime i, there is a stationary solution.

Second-order stationarity of the MS-AR(1) model

Problem: existence in L^2 of

$$\tilde{X}_t = \sum_{n=0}^{\infty} r_{t,n}, \qquad \text{where} \quad r_{t,n} = a(\Delta_t) \dots a(\Delta_{t-n+1}) \sigma(\Delta_{t-n}) \eta_{t-n}.$$

Norm on L^2 : $||X|| = {E(X^2)}^{1/2}$.

We have

$$\|\tilde{X}_t\|_{L^2} \leq \sum_{n=0}^{\infty} \|r_{t,n}\|_{L^2},$$

and

$$Er_{t,n}^2 = Ea^2(\Delta_t) \dots a^2(\Delta_{t-n+1})\sigma^2(\Delta_{t-n})$$
$$= (1,\dots,1)\mathbb{P}^n(a^2)\Pi(\sigma^2).$$

Second-order stationarity of the MS-AR(1) model

Proposition

lf

$$\rho\{\mathbb{P}(a^2)\} = \rho \left(\begin{array}{ccc} p(1,1)a^2(1) & \cdots & p(d,1)a^2(1) \\ \vdots & & \vdots \\ p(1,d)a^2(d) & \cdots & p(d,d)a^2(d) \end{array} \right) < 1,$$

 (\tilde{X}_t) is the unique 2nd-order stationary and non anticipative solution of the MS-AR(1) model.

Conversely, if $\rho\{\mathbb{P}(a^2)\} \ge 1$ there is no 2nd-order stationary and non anticipative solution.

Proof

• **Sufficient part**: We have, for some K > 0,

$$(1,\ldots,1)\mathbb{P}^n(a^2)\Pi(\sigma^2) \le K[\rho\{\mathbb{P}(a^2)\}]^i$$

- Uniqueness: same arguments as for the strict stationarity.
- Necessary part : Let (X_t) a 2nd-order stationary, non anticipative solution. We have

$$EX_{t}^{2} = E\{a(\Delta_{t}) \dots a(\Delta_{t-k+1}) X_{t-k}\}^{2} + E\left\{\sum_{n=0}^{k-1} r_{t,n}\right\}^{2}$$

$$\geq \sum_{n=0}^{k-1} E\{a(\Delta_{t}) \dots a(\Delta_{t-n+1}) \sigma(\Delta_{t-n})\}^{2} E\eta_{t}^{2}, \quad \forall k > 0.$$

Therefore, since $EX_t^2 < \infty$,

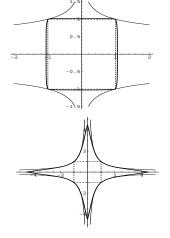
$$\sum_{n=0}^{\infty} E\{a(\Delta_t) \dots a(\Delta_{t-n+1}) \sigma(\Delta_{t-n})\}^2 = \sum_{i=0}^{\infty} (1, \dots, 1) \mathbb{P}^i(a^2) \Pi(\sigma^2) < \infty.$$

Remarks

- The condition involves the transition probabilities p(i,j) (not only the stationary probabilities $\pi(i)$).
- When the Markov chain is independent, we have p(i,j) = p(k,j) for all i,j,k. Thus, the columns of $\mathbb{P}(a^2)$ are identical. The only non-zero eigenvalue is $\sum_{i=1}^d \pi(i) a^2(i)$ and the condition reduces to

$$\sum_{i=1}^{d} \pi(i) a^{2}(i) < 1.$$

Stationarity regions when d=2 in the (a(1), a(2)) plane



$$p(1,1) = 0.8, p(2,2) = 0.95$$

$$p(1,1) = p(2,2) = 0.05$$

Autocovariance function

Under the 2nd-order stationarity condition, the autocovariance of (X_t) can be explicitly computed. We have

$$EX_t = \sum_{i=0}^{\infty} E\{a(\Delta_t) \dots a(\Delta_{t-i+1}) \sigma(\Delta_{t-i})\} E\eta_{t-i} = 0,$$

and

$$EX_t^2 = \sum_{i=0}^{\infty} E\{a(\Delta_t) \dots a(\Delta_{t-i+1})\sigma(\Delta_{t-i})\}^2$$
$$= \sum_{i=0}^{\infty} (1, \dots, 1) \mathbb{P}^i(a^2) \Pi(\sigma^2) = (1, \dots, 1) \{I_d - \mathbb{P}(a^2)\}^{-1} \Pi(\sigma^2).$$

Autocovariance function

For
$$h > 0$$

$$EX_{t}X_{t-h} = \sum_{i=0}^{\infty} E\{a(\Delta_{t}) \dots a(\Delta_{t-h+1})a^{2}(\Delta_{t-h}) \dots a^{2}(\Delta_{t-h-i+1})\sigma^{2}(\Delta_{t-h-i})\}$$

$$= \sum_{i=0}^{\infty} (1, \dots, 1)\mathbb{P}^{h}(a)\mathbb{P}^{i}(a^{2})\Pi(\sigma^{2})$$

$$= (1, \dots, 1)\mathbb{P}^{h}(a)\{I_{d} - \mathbb{P}(a^{2})\}^{-1}\Pi(\sigma^{2}).$$

Example with d=2

$$X_t = \begin{cases} \epsilon_t & \text{f} \quad \Delta_t = 1, \\ aX_{t-1} + \sigma\epsilon_t & \text{if} \quad \Delta_t = 2. \end{cases}$$

The model admits a strictly stationary solution whatever a. We have

$$\mathbb{P}(a^2) = \begin{pmatrix} 0 & 0 \\ p(1,2)a^2 & p(2,2)a^2 \end{pmatrix},$$

thus $\rho(\mathbb{P}(a^2)) = p(2,2)a^2$ and the 2nd-order stationarity condition is

$$p(2,2)a^2 < 1.$$

Under this condition

$$EX_t^2 = \frac{\{1 - a^2 p(2, 2) + a^2 p(1, 2)\}\pi(1) + \sigma^2 \pi(2)}{1 - a^2 p(2, 2)}.$$

Example with d=2

We have $\mathbb{P}^h(a) = \{ap(2,2)\}^{h-1}\mathbb{P}(a)$ for all h > 0. Therefore

$$\gamma(h) := EX_t X_{t-h} = \frac{p(1,2)\pi(1) + \sigma^2 p(2,2)\pi(2)}{1 - a^2 p(2,2)} a^h p^{h-1}(2,2), \quad h > 0$$

Note that

$$\gamma(h) = ap(2,2)\gamma(h-1), \qquad h > 1,$$

which entails that X_t admits an ARMA representation, of the form

$$X_t - ap(2,2)X_{t-1} = u_t - \theta u_{t-1}$$

where (u_t) is a white noise and θ is a coefficient (which can be obtained from $\gamma(0)$ and $\gamma(1)$).

This property, existence of an ARMA representation, is general to MS-ARMA processes.

Markov representation : $\mathbf{Z}_t = \mathbf{A}_t \mathbf{Z}_{t-1} + \mathbf{B}_t$

$$\boldsymbol{B}_t = \boldsymbol{C}_t + \underline{\boldsymbol{\epsilon}}_t = \boldsymbol{C}_t + \boldsymbol{\Sigma}_t \boldsymbol{\eta}_t,$$

$$\boldsymbol{C}_{t} = \begin{pmatrix} c(\Delta_{t}) \\ 0 \\ \vdots \\ 0 \\ 0 \\ \vdots \\ 0 \end{pmatrix} \in \mathbb{R}^{p+q}, \quad \boldsymbol{Z}_{t} = \begin{pmatrix} X_{t} \\ X_{t-1} \\ \vdots \\ X_{t-p+1} \\ \varepsilon_{t} \\ \varepsilon_{t-1} \\ \vdots \\ \varepsilon_{t-q+1} \end{pmatrix} \in \mathbb{R}^{p+q}, \quad \boldsymbol{\Sigma}_{t} = \begin{pmatrix} \sigma(\Delta_{t}) \\ 0 \\ \vdots \\ 0 \\ \sigma(\Delta_{t}) \\ 0 \\ \vdots \\ 0 \end{pmatrix},$$

$$\boldsymbol{A}_t = \begin{pmatrix} \boldsymbol{a}(\Delta_t) & \boldsymbol{b}(\Delta_t) \\ \mathbf{0} & \boldsymbol{J} \end{pmatrix} \in \mathcal{M}_{(p+q)\times(p+q)}(\mathbb{R})$$

Markov representation : $Z_t = A_t Z_{t-1} + B_t$

$$m{a}(\Delta_t) = egin{bmatrix} 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \ddots & \ddots & \vdots \\ 0 & & \cdots & 1 & 0 \end{pmatrix},$$
 $m{b}(\Delta_t) = egin{bmatrix} b_1(\Delta_t) & \cdots & b_q(\Delta_t) \\ 0 & 0 & \cdots & 0 \\ 0 & 0 & \cdots & 0 \\ \vdots & \ddots & \ddots & \vdots \\ 0 & & \cdots & 1 & 0 \end{pmatrix}, \quad m{J} = egin{bmatrix} 0 & 0 & \cdots & 0 \\ 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \ddots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 & 0 \end{pmatrix}$

Top-Lyapunov exponent

The top-Lyapunov exponent of the sequence $\{a(\Delta_t)\}$ is defined by

$$\gamma = \inf_{t>0} E\left(\frac{1}{t}\log\|\boldsymbol{a}(\Delta_{t})\boldsymbol{a}(\Delta_{t-1})\cdots\boldsymbol{a}(\Delta_{1})\|\right)$$

$$\stackrel{a.s.}{=} \lim_{t\to\infty} \frac{1}{t}\log\left\|\prod_{i=1}^{t}\boldsymbol{a}(\Delta_{t-i})\right\|$$

for any norm on $\mathscr{M}_{p \times p}(\mathbb{R})$.

Strict stationarity condition

Proposition

Suppose $\gamma_a < 0$. Then, for any $t \in \mathbb{Z}$, the series

$$\boldsymbol{Z}_t = \boldsymbol{B}_t + \sum_{k=1}^{\infty} \boldsymbol{A}_t \cdots \boldsymbol{A}_{t-k+1} \boldsymbol{B}_{t-k}$$

converges a.s. and the process (X_t) , defined as the first component of (Z_t) , is the unique strictly stationary solution of the MS-ARMA(p,q) model.

Remarks

- The strict stationarity condition only depends on the coefficients $a_i(k)$ (as for standard ARMA).
- In particular for a MS-ARMA(1,q), the condition reduces to

$$\sum_{i=1}^{d} \pi(i) \log |a_1(i)| < 0.$$

• The condition is only sufficient in general. The coefficients $b_j(k)$ and c(k) may matter for the strict stationarity.

Second-order stationarity

Let \otimes the Kronecker product. For any matrix function f defined on $\{1,\ldots,d\}$, let $f^{\otimes 2}(k)=f(k)\otimes f(k)$.

Proposition

Ιf

$$\rho\{\mathbb{P}(a^{\otimes 2})\}<1,$$

the strict stationary solution is also second-order stationary.

It can be shown that this condition is also necessary in the case p=q=1. However this is not general.

Example 1 : A two-regime MS-AR(2)

$$X_t = \begin{cases} \eta_t & \text{if } \Delta_t = 1\\ \eta_t + aX_{t-2} & \text{if } \Delta_t = 2 \end{cases}$$

Then

$$\mathbb{P}(A^{\otimes 2}) = \begin{pmatrix} 0_2 & 0_2 & 0_2 & 0_2 \\ P_1 & 0_2 & P_2 & 0_2 \\ 0_2 & aP_3 & 0_2 & aP_4 \\ P_3 & 0_2 & P_4 & 0_2 \end{pmatrix}$$

where 0_2 is the null 2×2 matrix and

$$P_{1} = \begin{pmatrix} 0 & 0 \\ p(1,1) & 0 \end{pmatrix}, \quad P_{2} = \begin{pmatrix} 0 & 0 \\ p(2,1) & 0 \end{pmatrix},$$

$$P_{3} = \begin{pmatrix} 0 & ap(1,2) \\ p(1,2) & 0 \end{pmatrix}, \quad P_{4} = \begin{pmatrix} 0 & ap(2,2) \\ p(2,2) & 0 \end{pmatrix}.$$

Example 1 : A two-regime MS-AR(2)

We have

$$\rho(\{\mathbb{P}(a^{\otimes 2})\}) = \max\left\{|a|p(2,2)^{1/2}, |a|\left\{p(2,2)^2 + p(1,2)p(2,1)\right\}^{1/2}\right\}.$$

But the expansion

$$X_t = \eta_t + \sum_{k=1}^{\infty} a^k \eta_{t-2k} \mathbf{1}_{\Delta_t = 2, \dots, \Delta_{t-2k+2} = 2}$$

shows that the necessary and sufficient 2nd-order stationarity condition is simply

$$|a| \{p(2,2)^2 + p(1,2)p(2,1)\}^{1/2} < 1,$$

which shows that the condition $\rho(\{\mathbb{P}(a^{\otimes 2})\}) < 1$ may be too strong.

Example 2 : local stationarity is neither necessary nor sufficient

It is clear (see for instance Example 1) that the **local stationarity** (i.e. within each regime) is not necessary to ensure the (global) 2nd-order stationarity.

The next example also shows that it is not sufficient.

$$X_{t} = \begin{cases} a_{1}(1)X_{t-1} + a_{2}(1)X_{t-2} + \eta_{t} & \text{if } \Delta_{t} = 1\\ a_{1}(2)X_{t-1} + \eta_{t} & \text{if } \Delta_{t} = 2 \end{cases}$$

where (η_t) is i.i.d. $\mathcal{N}(0,1)$.

If (X_t) is 2nd-order stationary then $E(X_t^2|\Delta_t=1,\Delta_{t-1}=2)$ exists and is independent of t.

Example 2 (continued)

Thus

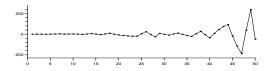
$$\begin{split} &E(X_t^2|\Delta_t=1,\Delta_{t-1}=2)\\ &=& E\Big(\big[\{a_1(1)a_1(2)+a_2(1)\}X_{t-2}+\eta_t+a_1(1)\eta_{t-1}\big]^2\,|\Delta_t=1,\Delta_{t-1}=2\Big)\\ &\geq& \{a_1(1)a_1(2)+a_2(1)\}^2\,E\Big(X_{t-2}^2|\Delta_t=1,\Delta_{t-1}=2\Big)\\ &\geq& \{a_1(1)a_1(2)+a_2(1)\}^2\,E\Big(X_{t-2}^2|\Delta_t=1,\Delta_{t-1}=2,\Delta_{t-2}=1,\Delta_{t-3}=2\Big)\\ &\times P(\Delta_{t-2}=1,\Delta_{t-3}=2|\Delta_t=1,\Delta_{t-1}=2)\\ &=& \{a_1(1)a_1(2)+a_2(1)\}^2\,P(2,1)P(1,2)E\Big(X_{t-2}^2|\Delta_{t-2}=1,\Delta_{t-3}=2\Big) \end{split}$$

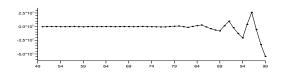
This is not possible when

$${a_1(1)a_1(2) + a_2(1)}^2 p(2,1)p(1,2) > 1.$$

But this condition is compatible with the local stationarity.

Example 2 (continued) : simulation





For instance if $a_1(1) = 1.8$, $a_2(1) = -0.9$, $a_1(2) = -0.2$, p(1,1) = 0.2 and p(2,2) = 0.1.

- 1 Hidden Markov Model
- 2 MS-ARMA(p,q) process
- 3 Estimation of MS-AR models
 - Computing the likelihood
 - Maximizing the likelihood
 - Illustration

MS-AR(p) model

$$X_t = \sum_{i=1}^p a_i(\Delta_t) X_{t-i} + \sigma(\Delta_t) \eta_t, \qquad (\eta_t) \stackrel{iid}{\sim} \mathcal{N}(0,1),$$

Parameter vector:

$$\boldsymbol{\theta} = (p(1,1), \dots, p(1,d-1), p(2,1), \dots, p(d,d-1), \sigma(1), \dots, \sigma(d))'.$$

The likelihood can be written by conditioning with respect to all possible paths

$$(e_1,\ldots,e_n)$$

of the Markov chain, where $e_i \in \mathcal{E} = \{1, \ldots, d\}$. The probability of this path is

$$P(e_1, ..., e_n) = P(\Delta_1 = e_1, ..., \Delta_n = e_n) = \pi(e_1)p(e_1, e_2)...p(e_{n-1}, e_n).$$

Likelihood

For any path, we a have a likelihood given by

$$L^{(e_1,\ldots,e_n)}(X_1,\ldots,X_n) = \prod_{t=1}^n \phi_{e_t} \left(X_t - \sum_{i=1}^p a_i(e_t) X_{t-i} \right),$$

where $\phi_i(\cdot)$ is the density of $\mathcal{N}\left\{0,\sigma^2(i)\right\}$.

Finally the likelihood of the observations is

$$L_{\theta}(X_1,...,X_n) = \sum_{(e_1,...,e_n) \in \mathscr{E}^n} L_{\sigma}^{(e_1,...,e_n)}(X_1,...,X_n) P(e_1,...,e_n).$$

However, such a formula cannot be used in practice.

Several algorithms can be used to compute the likelihood.

Idea : decompose the likelihood by conditioning on the last state

$$L_{\boldsymbol{\theta}}(X_1,\ldots,X_n) = \sum_{i=1}^d L_{\boldsymbol{\theta}}(X_1,\ldots,X_n \mid \Delta_n = i) \pi(i).$$

Let

$$F_k(i) = g_k(X_1, \dots, X_k | \Delta_k = i) \pi(i)$$

where $g_k(\cdot|\Delta_k=i)$ is the density of (X_1,\ldots,X_k) given $\{\Delta_k=i\}$.

We have (with initial values $X_0 = X_{-1} = \dots = X_{1-p} = 0$)

$$F_1(i) = g_1(X_1 \mid \Delta_1 = i)\pi(i) = \phi_i(X_1)\pi(i).$$

$$\begin{split} F_k(i) &= g(X_k|X_1,\dots,X_{k-1},\Delta_k=i)g_{k-1}(X_1,\dots,X_{k-1}|\Delta_k=i)\pi(i) \\ &= \phi_i \left(X_k - \sum_{\ell=1}^p a_\ell(i)X_{k-\ell}\right) \sum_{j=1}^d g_{k-1}(X_1,\dots,X_{k-1}|\Delta_{k-1}=j,\Delta_k=i) \\ &\qquad \qquad \times P(\Delta_{k-1}=j\mid \Delta_k=i)\pi(i) \\ &= \frac{1}{\sigma(i)}\phi\{\eta_k(i)\} \sum_{j=1}^d F_{k-1}(j)p(j,i) \end{split}$$

where

$$\eta_k(i) = \frac{1}{\sigma(i)} \left(X_k - \sum_{\ell=1}^p a_\ell(i) X_{k-\ell} \right).$$

Finally, the algorithm is given by

$$F_{1}(i) = \phi_{i}(X_{1})\pi(i), \quad i = 1, ..., d,$$

$$F_{k}(i) = \frac{1}{\sigma(i)}\phi\{\eta_{k}(i)\}\sum_{j=1}^{d}F_{k-1}(j)p(j,i), \quad i = 1, ..., d, \quad k > 1$$

and the likelihood is obtained as

$$L_{\theta}(X_1,\ldots,X_n)=\sum_{i=1}^d F_n(i).$$

In matrix form:

$$F_k := (F_k(1), \dots, F_k(d))' = M(X_k, \dots, X_{k-p})F_{k-1},$$

where

$$M(X_k, \dots, X_{k-p}) = \left(\begin{array}{ccc} p(1,1) \frac{\phi(\eta_k(1))}{\sigma(1)} & \cdots & p(d,1) \frac{\phi(\eta_k(1))}{\sigma(1)} \\ \vdots & & \vdots \\ p(1,d) \frac{\phi(\eta_k(d))}{\sigma(d)} & \cdots & p(d,d) \frac{\phi(\eta_k(d))}{\sigma(d)} \end{array} \right).$$

Hence

$$L_{\theta}(X_1,...,X_n) = \mathbf{1}' M(X_n,...,X_{n-p}) M(X_{n-1},...,X_{n-1-p}) \cdots M(X_2,...,X_{2-p}) F_1.$$

which is numerically tractable $(O(d^2n)$ multiplications).

However, the algorithm can be unstable if n is large (underflows due to a very small likelihood).

Forward-Backward Algorithm [Baum (1972)]

Let

$$B_k(i) = g_{n-k}(X_{k+1}, ..., X_n | \Delta_k = i, X_1, ..., X_k).$$

We have

$$L_{\boldsymbol{\theta}}(X_1, \dots, X_n) = \sum_{i=1}^d L_{\boldsymbol{\theta}}(X_1, \dots, X_n \mid \Delta_k = i) \pi(i)$$
$$= \sum_{i=1}^d B_k(i) F_k(i).$$

Forward formulas allow to compute $F_k(i)$ for k = 1, 2, ...

Backward formulas allow to compute $B_k(i)$ for k = n - 1, n - 2,...

Forward-Backward Algorithm

$$\begin{split} B_{k}(i) &= \sum_{j=1}^{d} g_{n-k}(X_{k+1}, \dots, X_{n} | \Delta_{k+1} = j, \Delta_{k} = i, X_{1}, \dots, X_{k}) P(\Delta_{k+1} = j | \Delta_{k} = i) \\ &= \sum_{j=1}^{d} g_{n-k-1}(X_{k+2}, \dots, X_{n} | \Delta_{k+1} = j, X_{1}, \dots, X_{k}, \underbrace{X_{k+1}}) \\ &\qquad \qquad \times g(X_{k+1} | \Delta_{k+1} = j, X_{1}, \dots, X_{k}) P(\Delta_{k+1} = j | \Delta_{k} = i) \\ &= \sum_{j=1}^{d} B_{k+1}(j) \frac{\phi \{ \eta_{k+1}(j) \}}{\sigma(j)} p(i, j) \end{split}$$

Forward-Backward Algorithm

Backward formulas:

$$B_n(i) = 1,$$

$$B_k(i) = \sum_{j=1}^d B_{k+1}(j) \frac{\phi \{ \eta_{k+1}(j) \}}{\sigma(j)} p(i,j), \quad k < n.$$

Compare with the Forward formulas:

$$\begin{split} F_1(i) &= \phi_i(X_1)\pi(i), \\ F_k(i) &= \frac{\phi\{\eta_k(i)\}}{\sigma(i)} \sum_{j=1}^d F_{k-1}(j) p(j,i), \quad k > 1. \end{split}$$

For any k = 1, ..., n, the likelihood is given by

$$L_{\boldsymbol{\theta}}(X_1,\ldots,X_n) = \sum_{i=1}^d B_k(i)F_k(i).$$

Introduced by Hamilton (1989, Econometrica).

Main differences with respect to the previous algorithms :

- Based on the log-likelihood
- Provides filtered probabilities of the regimes

We have (neglecting the distribution of X_1):

$$\log L_{\theta}(X_1,...,X_n) = \sum_{t=1}^n \log f_t(X_t \mid X_{t-1},...,X_1),$$

where

$$f_t(X_t \mid X_{t-1}, \dots, X_1) = \sum_{j=1}^d f_t(X_t \mid X_{t-1}, \dots, X_1, \Delta_t = j) P(\Delta_t = j \mid X_{t-1}, \dots, X_1).$$

Let

$$\pi_{t|t-1}(j) = P(\Delta_t = j \mid X_{t-1}, ..., X_1),$$

 $\pi_{t|t}(j) = P(\Delta_t = j \mid X_t, ..., X_1).$

We have (here g denotes a generic density)

$$\pi_{t+1|t}(j) = \sum_{j=1}^{d} P(\Delta_{t+1} = j \mid \Delta_{t} = i, X_{t}, \dots, X_{1}) \pi_{t|t}(i) = \sum_{j=1}^{d} p(i, j) \pi_{t|t}(i),$$

$$\pi_{t|t}(j) = \frac{g(X_{t}, \dots, X_{1} \mid \Delta_{t} = j) \pi(j)}{g(X_{t}, \dots, X_{1})},$$

using the formula $P(A|X=x) = \frac{f(x|A)}{f(x)}P(A)$.

$$\begin{split} \pi_{t|t}(j) &= \frac{g(X_t, \dots, X_1 \mid \Delta_t = j)\pi(j)}{g(X_t, \dots, X_1)} \\ &= \frac{\frac{\phi\{\eta_t(j)\}}{\sigma(j)}g(X_{t-1}, \dots, X_1 \mid \Delta_t = j)\pi(j)}{g(X_t, \dots, X_1)} \\ &= \frac{\frac{\phi\{\eta_t(j)\}}{\sigma(j)}P(\Delta_t = j \mid X_{t-1}, \dots, X_1)g(X_{t-1}, \dots, X_1)}{g(X_t \mid X_{t-1}, \dots, X_1)g(X_{t-1}, \dots, X_1)} \\ &= \frac{\frac{\phi\{\eta_t(j)\}}{\sigma(j)}\pi_{t|t-1}(j)}{g(X_t \mid X_{t-1}, \dots, X_1)} \\ &= \frac{\frac{\phi\{\eta_t(j)\}}{\sigma(j)}\pi_{t|t-1}(j)}{\sum_{l=1}^d \frac{\phi\{\eta_t(l)\}}{\sigma(l)}\pi_{t|t-1}(l)} \end{split}$$

Finally, the sequences

$$\pi_{t|t-1}(j) = P(\Delta_t = j \mid X_{t-1}, \dots, X_1), \quad \pi_{t|t}(j) = P(\Delta_t = j \mid X_t, \dots, X_1)$$

are obtained recursively from

$$\begin{cases} \pi_{t|t}(j) &= \frac{\frac{\phi(\eta_t(j))}{\sigma(j)}\pi_{t|t-1}(j)}{\sum_{i=1}^d \frac{\phi(\eta_t(i))}{\sigma(i)}\pi_{t|t-1}(i)}, \\ \\ \pi_{t+1|t}(j) &= \sum_{j=1}^d p(i,j)\pi_{t|t}(i), \end{cases}$$

with initial values $\pi_{1|0}(i) = \pi(i)$ (or any other distribution).

Numerically, the filter generates less underflows and tends to perform better than the (standard) forward-backward algorithm.

Hamilton filter in matrix form

Let

$$\boldsymbol{\pi}_{t|t} = \left(\begin{array}{c} \boldsymbol{\pi}_{t|t}(1) \\ \vdots \\ \boldsymbol{\pi}_{t|t}(d) \end{array} \right), \quad \boldsymbol{\pi}_{t+1|t} = \left(\begin{array}{c} \boldsymbol{\pi}_{t+1|t}(1) \\ \vdots \\ \boldsymbol{\pi}_{t+1|t}(d) \end{array} \right), \quad \boldsymbol{\Phi}_{t} = \left(\begin{array}{c} \frac{\boldsymbol{\phi}\{\eta_{t}(1)\}}{\boldsymbol{\sigma}(1)} \\ \vdots \\ \frac{\boldsymbol{\phi}\{\eta_{t}(d)\}}{\boldsymbol{\sigma}(d)} \end{array} \right).$$

⊙ : Hadamard product (componentwise).

We have

$$\pi_{t|t} = \frac{\pi_{t|t-1} \odot \Phi_t}{(1, \dots, 1) \left\{ \pi_{t|t-1} \odot \Phi_t \right\}}, \qquad \pi_{t+1|t} = \mathbb{P} \pi_{t|t}$$

EM algorithm

Maximisation of the (log-)likelihood can be achieved either using a classical optimization procedure, or using the EM (Expectation–Maximization) algorithm.

The EM algorithm, proposed by Dempster, Laird and Rubin (1977)[†] is used to find ML estimators of parameters in general statistical models involving latent variables.

^{†.} Maximum Likelihood from incomplete data via the EM algorithm. *J. Roy. Statist. Soc. B* 39, 1–38.

Parameter of interest : $oldsymbol{ heta}$

Likelihood inference based on the observed data, Y, is intractable.

Likelihood inference based on a completed data set, (X, Y), becomes tractable.

The log-likelihood of the observation Y can be decomposed as

$$\log \ell_{\boldsymbol{\theta}}(Y) = \log \ell_{\boldsymbol{\theta}}(X, Y) - \log \ell_{\boldsymbol{\theta}}(X|Y).$$

Multiplying both sides by $\ell_{\tilde{\boldsymbol{\theta}}}(X|Y)$ and integrating w.r.t. X yields

$$\log \ell_{\boldsymbol{\theta}}(Y) = Q(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) - H(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}})$$

where

$$Q(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) = E_{\tilde{\boldsymbol{\theta}}} \left[\log \ell_{\boldsymbol{\theta}}(X, Y) | Y \right] \quad \text{and} \quad H(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) = E_{\tilde{\boldsymbol{\theta}}} \left[\log \ell_{\boldsymbol{\theta}}(X | Y) | Y \right].$$

The difference

$$H(\tilde{\boldsymbol{\theta}}, \tilde{\boldsymbol{\theta}}) - H(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) = -E_{\tilde{\boldsymbol{\theta}}} \left[\log \frac{\ell_{\boldsymbol{\theta}}(X|Y)}{\ell_{\tilde{\boldsymbol{\theta}}}(X|Y)} |Y| \right]$$

is the Kullback-Leibler divergence between the conditional distributions $\ell_{\boldsymbol{\theta}}(X|Y)$ and $\ell_{\tilde{\boldsymbol{\theta}}}(X|Y)$.

Thus,

$$\log \ell_{\boldsymbol{\theta}}(Y) - \log \ell_{\tilde{\boldsymbol{\theta}}}(Y) \ge Q(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) - Q(\tilde{\boldsymbol{\theta}}, \tilde{\boldsymbol{\theta}}),$$

where the inequality is strict unless if $\ell_{\tilde{\boldsymbol{\theta}}}(\cdot|Y) = \ell_{\boldsymbol{\theta}}(\cdot|Y)$, a.e. Moreover, under regularity conditions,

$$\frac{\partial}{\partial \boldsymbol{\theta}} \log \ell_{\boldsymbol{\theta}}(Y) = \frac{\partial}{\partial \boldsymbol{\theta}} Q(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) \mid_{\tilde{\boldsymbol{\theta}} = \boldsymbol{\theta}}.$$

These results suggest that $\theta \mapsto Q(\theta, \tilde{\theta})$ can be used as a surrogate for the log-likelihood function $\log \ell_{\theta}(Y)$.

The EM algorithm uses this idea to maximize the (incomplete) likelihood $\log \ell_{\boldsymbol{\theta}}(Y)$, by iteratively maximizing the auxiliary function $\boldsymbol{\theta} \mapsto Q(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}})$.

Starting from an initial value θ_0 , the procedure iterates between 2 steps for computing a new parameter value $\theta^{(k)}$ from $\theta^{(k-1)}$:

- **(E)** Compute $\theta \mapsto Q(\theta, \theta^{(k-1)})$;
- (M) Compute $\theta^{(k)} = \operatorname{argmax}_{\theta \in \Theta} Q(\theta, \theta^{(k-1)})$.

In view of

$$\log \ell_{\boldsymbol{\theta}}(Y) - \log \ell_{\tilde{\boldsymbol{\theta}}}(Y) \ge Q(\boldsymbol{\theta}, \tilde{\boldsymbol{\theta}}) - Q(\tilde{\boldsymbol{\theta}}, \tilde{\boldsymbol{\theta}}),$$

the sequence $\{\log \ell_{\boldsymbol{\theta}^{(k)}}(Y)\}$ of log-likelihoods is non decreasing.

Under appropriate regularity conditions, it can be shown that the sequence $(\theta^{(k)})$ converges to some value θ^* (though it may be a local maximum of the likelihood).

EM algorithm for the HMM

We only consider the **HMM case**:

$$\epsilon_t = \sigma(\Delta_t)\eta_t$$
.

Observations : $\epsilon_1, ..., \epsilon_n$

Parameters : $\boldsymbol{\theta}$ and the initial distribution $\boldsymbol{\pi}_0 = (\pi_0(1), \dots, \pi_0(d))$.

Recall that likelihood of the observations is

$$L_{\boldsymbol{\theta},\boldsymbol{\pi}_0}(\epsilon_1,\ldots,\epsilon_n) = \sum_{(e_1,\ldots,e_n)\in\mathcal{E}^n} L_{\sigma}^{(e_1,\ldots,e_n)}(\epsilon_1,\ldots,\epsilon_n) P(e_1,\ldots,e_n).$$

If, in addition, one could observe $(\Delta_1, ..., \Delta_n)$, $\boldsymbol{\theta}$ and $\boldsymbol{\pi}_0$ could be easily estimated by ML.

Maximizing the likelihood as if the chain was observed

Indeed

$$\begin{split} & \log L_{\boldsymbol{\theta}, \pi_0}(\epsilon_1, \dots, \epsilon_n, \Delta_1, \dots, \Delta_n) \\ &= & \sum_{t=1}^n \log \phi_{\Delta_t}(\epsilon_t) + \log \pi_0(\Delta_1) + \sum_{t=2}^n \log p(\Delta_{t-1}, \Delta_t) = a_1 + a_2 + a_3 \end{split}$$

where

$$\begin{aligned} a_1 &= a_1(\sigma^2) = \sum_{i=1}^d \sum_{t=1}^n \log \phi_i(\epsilon_t) \mathbf{1}_{\{\Delta_t = i\}}, \\ a_2 &= a_2(\boldsymbol{\pi}_0) = \log \boldsymbol{\pi}_0(\Delta_1), \\ a_3 &= a_3(P) = \sum_{i=1}^d \sum_{j=1}^d \log p(i,j) \sum_{t=2}^n \mathbf{1}_{\{\Delta_{t-1} = i, \Delta_t = j\}}. \end{aligned}$$

Maximizing the likelihood as if the chain was observed

Maximization of

$$a_1(\sigma^2) = \sum_{i=1}^d \sum_{t=1}^n \log \left\{ \frac{1}{\sigma(i)} \phi \left(\frac{\epsilon_t}{\sigma(i)} \right) \right\} \mathbf{1}_{\{\Delta_t = i\}}$$

with respect to the $\sigma^2(i)$, yields the « estimators »

$$\tilde{\sigma}^{2}(i) = \frac{1}{\sum_{t=1}^{n} \mathbf{1}_{\{\Delta_{t}=i\}}} \sum_{t=1}^{n} \epsilon_{t}^{2} \mathbf{1}_{\{\Delta_{t}=i\}}.$$

• Maximization of a_2 in $\pi_0(1), \ldots, \pi_0(d)$, under the constraint $\sum_{i=1}^d \pi_0(i) = 1$, yields

$$\tilde{\pi}_0(i)=\mathbf{1}_{\{\Delta_1=i\}}.$$

Maximizing the likelihood as if the chain was observed

• For $i=1,\ldots,d$: maximization in $p(i,1),\ldots,p(i,d)$, under the constraint $\sum_{i=1}^d p(i,j)=1$, of

$$\sum_{j=1}^{d} \log p(i,j) \frac{\sum_{t=2}^{n} \mathbf{1}_{\{\Delta_{t-1}=i,\Delta_{t}=j\}}}{\sum_{t=2}^{n} \mathbf{1}_{\{\Delta_{t-1}=i\}}}$$

yields

$$\tilde{p}(i,j) = \frac{1}{\sum_{t=2}^{n} \mathbf{1}_{\{\Delta_t = i\}}} \sum_{t=2}^{n} \mathbf{1}_{\{\Delta_{t-1} = i, \Delta_t = j\}}.$$

Indeed, if p_1,\ldots,p_n are positive numbers such that $\sum_i p_i=1$, under the constraint $\sum_{i=1}^d \pi_i=1$, the global maximum of the function $(\pi_1,\ldots,\pi_d)\to \sum_i p_i\log\pi_i$ is reached at $(\pi_1,\ldots,\pi_d)=(p_1,\ldots,p_d)$.

EM algorithm : E-step

In practice, (Δ_t) is not observed and the previous estimators cannot be used.

The idea is to replace the quantities depending on the MC by their expectation.

Suppose that at some step k, we have an estimator $(\theta^{(k)}, \pi_0^{(k)})$.

The unknown likelihood is approximated by its expectation given the observations $(\epsilon_1, \ldots, \epsilon_n)$, computed under the law of parameter $(\theta^{(k)}, \pi_0^{(k)})$.

EM algorithm : E-step

We get the criterion

$$Q(\theta, \pi_0 | \theta^{(k)}, \pi_0^{(k)}) = E_{\theta^{(k)}, \pi_0^{(k)}} \left\{ \log L_{\theta, \pi_0}(\epsilon_1, \dots, \epsilon_n, \Delta_1, \dots, \Delta_n) | \epsilon_1, \dots, \epsilon_n \right\}$$
$$= A_1(\sigma) + A_2(\pi_0) + A_3(P),$$

where

$$\begin{split} A_1(\sigma) &= \sum_{i=1}^d \sum_{t=1}^n \log \phi_i(\epsilon_t) P_{\theta^{(k)}, \pi_0^{(k)}} \{ \Delta_t = i | \epsilon_1, \dots, \epsilon_n \}, \\ A_2(\pi_0) &= \sum_{i=1}^d \log \pi_0(i) P_{\theta^{(k)}, \pi_0^{(k)}} \{ \Delta_1 = i | \epsilon_1, \dots, \epsilon_n \}, \\ A_3(P) &= \sum_{i,i} \log p(i,j) \sum_{t=2}^n P_{\theta^{(k)}, \pi_0^{(k)}} \{ \Delta_{t-1} = i, \Delta_t = j | \epsilon_1, \dots, \epsilon_n \}. \end{split}$$

M-step

0

0

0

In this step we solve

$$\max_{(\theta,\pi_0)} Q(\theta,\pi_0|\theta^{(k)},\pi_0^{(k)}).$$

$$\hat{\sigma}^2(i) = \frac{\sum_{t=1}^n \epsilon_t^2 P_{\theta^{(k)}, \pi_0^{(k)}} \{ \Delta_t = i | \epsilon_1, \dots, \epsilon_n \}}{\sum_{t=1}^n P_{\theta^{(k)}, \pi_0^{(k)}} \{ \Delta_t = i | \epsilon_1, \dots, \epsilon_n \}}.$$

$$\hat{\pi}_0(i) = P_{\theta^{(k)}, \pi_0^{(k)}} \{ \Delta_1 = i | \epsilon_1, \dots, \epsilon_n \},$$

$$\hat{p}(i,j) = \frac{\sum_{t=2}^{n} P_{\theta^{(k)},\pi_0^{(k)}} \left\{ \Delta_{t-1} = i, \Delta_t = j | \epsilon_1, \dots, \epsilon_n \right\}}{\sum_{t=2}^{n} P_{\theta^{(k)},\pi_0^{(k)}} \left\{ \Delta_{t-1} = i | \epsilon_1, \dots, \epsilon_n \right\}}.$$

M-step

Starting from an initial value $(\theta^{(0)}, \pi_0^{(0)})$, these formulas allow to build a sequence $(\theta^{(k)}, \pi_0^{(k)})_k$ which increases the likelihood.

Require to compute the smoothed probabilities

$$\pi_{t|n} = (P\{\Delta_t = i|\epsilon_1, \dots, \epsilon_n\})'_{1 \le i \le d} \in \mathbb{R}^d$$

and

$$\pi_{t-1,t|n} = \left(P\left\{\Delta_{t-1} = i, \Delta_t = j|\epsilon_1, \dots, \epsilon_n\right\}\right)'_{1 \leq i,j \leq d} \in \mathbb{R}^d \times \mathbb{R}^d,$$

(still evaluated with the parameters $heta^{(k)}, \pi_0^{(k)}$).

Such probabilities are obtained from Hamilton's algorithm.

Smoothed probabilities

The Markov property entails that given Δ_t , the observations $\epsilon_t, \epsilon_{t+1}, \ldots$ do not convey information on Δ_{t-1} . Thus

$$P(\Delta_{t-1} = i | \Delta_t = j, \epsilon_1, \dots, \epsilon_n) = P(\Delta_{t-1} = i | \Delta_t = j, \epsilon_1, \dots, \epsilon_{t-1})$$

and

$$\pi_{t-1,t|n}(i,j) = P(\Delta_{t-1} = i|\Delta_t = j, \epsilon_1, \dots, \epsilon_{t-1})\pi_{t|n}(j)
= \frac{p(i,j)\pi_{t-1|t-1}(i)\pi_{t|n}(j)}{\pi_{t|t-1}(j)}.$$

Moreover, for $t = n, n-1, \ldots, 2$,

$$\pi_{t-1|n}(i) \quad = \quad \sum_{j=1}^d \pi_{t-1,t|n}(i,j) = \sum_{j=1}^d \frac{p(i,j)\pi_{t-1|t-1}(i)\pi_{t|n}(j)}{\pi_{t|t-1}(j)}$$

EM algorithm

Starting from initial values for the parameters π_0 , p(i,j), σ_i , the EM algorithm consists in repeating until convergence the steps:

① Set $\pi_{1|0} = \pi_0$ and

$$\pi_{t|t} = \frac{\pi_{t|t-1} \odot \Phi_t}{(1, \dots, 1) \left\{ \pi_{t|t-1} \odot \Phi_t \right\}}, \quad \pi_{t+1|t} = \mathbb{P}\pi_{t|t}, \quad \text{for } t = 1, \dots, n.$$

2 Compute the smoothed probabilities $\pi_{t|n}(i)$ and $\pi_{t-1,t|n}(i,j)$:

$$\begin{array}{rcl} \pi_{t-1|n}(i) & = & \displaystyle \sum_{j=1}^d \frac{p(i,j)\pi_{t-1|t-1}(i)\pi_{t|n}(j)}{\pi_{t|t-1}(j)} & \text{for } t=n,n-1,\ldots,2, \\ \\ \pi_{t-1,t|n}(i,j) & = & \displaystyle \frac{p(i,j)\pi_{t-1|t-1}(i)\pi_{t|n}(j)}{\pi_{t|t-1}(j)}. \end{array}$$

3 Replace the previous parameter values by $\pi_0 = \pi_{1|n}$,

$$p(i,j) = \frac{\sum_{t=2}^{n} \pi_{t-1,t|n}(i,j)}{\sum_{t=2}^{n} \pi_{t-1|n}(i)} \quad \text{and} \quad \sigma^{2}(i) = \frac{\sum_{t=1}^{n} \varepsilon_{t}^{2} \pi_{t|n}(i)}{\sum_{t=1}^{n} \pi_{t|n}(i)}.$$

CAC40 and SP500 series

Daily series of the CAC40 and SP500 over the period : March 1st 1990 to December 29 2006.

HMM model with 4 regimes : 60 iterations of the EM algorithm yield

SP500

$$\hat{\omega}_{SP} = \begin{pmatrix} 0.26 \\ 0.62 \\ 1.28 \\ 4.8 \end{pmatrix}, \quad \hat{P}_{SP} = \begin{pmatrix} 0.981 & 0.019 & 0.000 & 0.000 \\ 0.018 & 0.979 & 0.003 & 0.000 \\ 0.000 & 0.003 & 0.986 & 0.011 \\ 0.000 & 0.000 & 0.055 & 0.945 \end{pmatrix}$$

CAC40

$$\hat{\omega}_{CAC} = \left(\begin{array}{c} 0.51 \\ 1.19 \\ 2.45 \\ 8.4 \end{array} \right), \quad \hat{P}_{CAC} = \left(\begin{array}{cccc} 0.993 & 0.003 & 0.002 & 0.002 \\ 0.003 & 0.991 & 0.003 & 0.003 \\ 0.000 & 0.020 & 0.977 & 0.003 \\ 0.004 & 0.000 & 0.032 & 0.963 \end{array} \right).$$

CAC40 and SP500 series

Estimated probabilities of the 4 regimes :

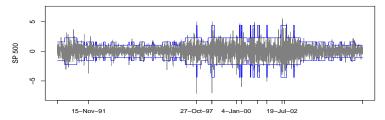
$$\hat{\pi}_{SP} = (0.30, 0.32, 0.32, 0.06)', \qquad \hat{\pi}_{CAC} = (0.26, 0.49, 0.19, 0.06)',$$

Expected duration of the different regimes (equal to $1/\{1-p(i,i)\}$):

$$D_{SP} = (53, 48, 71, 18)', \qquad D_{CAC} = (140, 107, 43, 27)'.$$

CAC40 and SP500 series

Rendements du SP 500



Rendements du CAC 40

