## TMS 165/MSA350 Stochastic Calculus Part I Fall 2012 Exercise session 2

**Exercise 1.** Prove Equations 2.17 and 2.21 in Klebaner's book for conditional expectations.

**Solution.** It is an easy exercise to see that any constant random variable (that is, a non-random random variable) is measurable wrt. the trivial  $\sigma$ -field  $\{\emptyset, \Omega\}$ . In particular,  $\mathbf{E}\{X\}$  is  $\{\emptyset, \Omega\}$ -measurable. Further we have

$$\int_{\emptyset} \mathbf{E}\{X\} \, d\mathbf{P} = 0 = \int_{\emptyset} X \, d\mathbf{P} \quad \text{and} \quad \int_{\Omega} \mathbf{E}\{X\} \, d\mathbf{P} = \mathbf{P}\{\Omega\} \, \mathbf{E}\{X\} = \mathbf{E}\{X\} = \int_{\Omega} X \, d\mathbf{P}.$$

Hence  $\mathbf{E}\{X\}$  fulfill the defining properties on page 44 in the book of being the conditional expectation  $\mathbf{E}\{X | \{\emptyset, \Omega\}\}$ . This establishes (2.17).

As for (2.21), as  $\mathbf{E}\{X\}$  is  $\{\emptyset, \Omega\}$ -measurable it is measurable wrt. any other  $\sigma$ -field  $\mathcal{G} \subseteq \mathcal{F}$  (as any such  $\mathcal{G}$  must contain  $\{\emptyset, \Omega\}$ ). For X independent of  $\mathcal{G}$  we further have

$$\int_A X d\mathbf{P} = \mathbf{E}\{I_A X\} = \mathbf{E}\{I_A\} \mathbf{E}\{X\} = \mathbf{P}\{A\} \mathbf{E}\{X\} = \int_A \mathbf{E}\{X\} d\mathbf{P} \quad \text{for } A \in \mathcal{G}.$$

Hence  $\mathbf{E}\{X\}$  fulfill the defining properties of being the conditional expectation  $\mathbf{E}\{X|\mathcal{G}\}$ .

**Exercise 2.** Consider a finite sample space  $\Omega = \{1, ..., 2n\}$  equipped with the  $\sigma$ -field  $\mathcal{F}$  consisting of all subsets of  $\Omega$  together with the uniform probability measure  $\mathbf{P}$  on  $\Omega$  assigning probability 1/(2n) to each outcome  $\omega \in \Omega$ . Calculate  $\mathbf{E}\{X | \mathcal{G}\}$  for the random variable  $X(\omega) = \omega$  and the  $\sigma$ -field  $\mathcal{G} = \{\emptyset, A, A^c, \Omega\}$  where  $A = \{1, ..., n\}$ .

**Solution.** From intuitive reasoning we come up with the hypothesis that

$$\mathbf{E}\{X|\mathcal{G}\} = \begin{cases} (n+1)/2 & \text{for } \omega \in A \\ (3n+1)/2 & \text{for } \omega \in A^c \end{cases}.$$

That this really is correct follows from the fact that this random variable is  $\mathcal{G}$ -measurable and that by elementary calculations together with the uniformity of  $\mathbf{P}$  it satisfies

$$\int_{B} \mathbf{E}\{X | \mathcal{G}\} d\mathbf{P} = \int_{B} X d\mathbf{P} \quad \text{for } B \in \{\emptyset, A, A^{c}, \Omega\}.$$

**Exercise 3.** Show that among all zero-mean stochastic processes  $\{X(t)\}_{t\geq 0}$  with finite second moments  $\mathbf{E}\{X(t)^2\} < \infty$  for  $t\geq 0$ , the class of martingales contain all processes

with independent increments and are all included among processes with uncorrelated increments.

**Solution.** For X zero-mean with independent increments we have

$$\mathbf{E}\{X(t)|\mathcal{F}_{s}^{X}\} = \mathbf{E}\{X(t) - X(s)|\mathcal{F}_{s}^{X}\} + \mathbf{E}\{X(s)|\mathcal{F}_{s}^{X}\} = \mathbf{E}\{X(t) - X(s)\} + X(s) = X(s)$$

for  $s \leq t$ , where we use the independent increments and (2.21) together with the fact that X is adapted to the  $\sigma$ -field  $\{\mathcal{F}^X_t\}_{t\geq 0}$ . Hence X is a martingale.

On the other hand, for X a zero-mean martingale we have

$$\begin{split} \mathbf{E}\{(X(u)-X(t))\left(X(s)-X(r)\right)\} &= \mathbf{E}\big\{\mathbf{E}\{(X(u)-X(t))\left(X(s)-X(r)\right)|\mathcal{F}_s^X\big\}\big\} \\ &= \mathbf{E}\big\{(X(s)-X(r))\,\mathbf{E}\{X(u)-X(t)|\mathcal{F}_s^X\}\big\} \\ &= \mathbf{E}\{(X(s)-X(r))\left(X(s)-X(s)\right)\} \\ &= 0 \end{split}$$

for  $0 \le r \le s \le t \le u$ , where we made use of Equation 2.20 in Klebaner's book and the fact that X is adapted together with Equation 2.18 and the martingale property.

**Exercise 4.** Prove Equation 3.4 in Klebaner's book. (Note that it is assumed that  $0 < t_1 < \ldots < t_n$  in this formula.)

**Solution.** We prove (3.4) by induction. Note that the property (3.4) when n = 1 is just (3.3), which in turn is a rather elementary formula we proved during Lecture 4.

Now assume that (3.4) holds for n = k. Note that (3.4) for n = k in turn means that  $(B^x(t_1), \ldots, B^x(t_k))$  has probability density function

$$f_{(B^x(t_1),\dots,B^x(t_k))}(y_1,\dots,y_k) = p_{t_1}(x,y_1) \prod_{i=2}^k p_{t_i-t_{i-1}}(y_{i-1},y_i) \text{ for } (y_1,\dots,y_k) \in \mathbb{R}^k.$$

For the case when n = k+1 it therefore follows from conditioning on the value  $(y_1, \ldots, y_k)$  of  $(B^x(t_1), \ldots, B^x(t_k))$  and using independence of increments that

$$\mathbf{P}\left\{\bigcap_{i=1}^{k+1} \{B^{x}(t_{i}) \leq x_{i}\}\right\} \\
= \int_{-\infty}^{x_{1}} \dots \int_{-\infty}^{x_{k}} \mathbf{P}\left\{B^{x}(t_{k+1}) - B^{x}(t_{k}) + y_{k} \leq x_{k+1}\right\} f_{(B^{x}(t_{1}), \dots, B^{x}(t_{k}))}(y_{1}, \dots, y_{k}) dy_{1} \dots dy_{k} \\
= \int_{-\infty}^{x_{1}} \dots \int_{-\infty}^{x_{k}} \Phi\left(\frac{x_{k+1} - y_{k}}{\sqrt{t_{k+1} - t_{k}}}\right) p_{t_{1}}(x, y_{1}) \prod_{i=2}^{k} p_{t_{i} - t_{i-1}}(y_{i-1}, y_{i}) dy_{1} \dots dy_{k} \\
= \int_{-\infty}^{x_{1}} \dots \int_{-\infty}^{x_{k}} \int_{-\infty}^{x_{k+1}} p_{t_{1}}(x, y_{1}) \prod_{i=2}^{k+1} p_{t_{i} - t_{i-1}}(y_{i-1}, y_{i}) dy_{1} \dots dy_{k+1}$$

[as  $B^x(t_{k+1}) - B^x(t_k)$  is  $N(0, t_{k+1} - t_k)$ -distributed]. This proves (3.4) by induction.

**Exercise 5.** Let  $\xi$  and  $\eta$  be independent standard normal random variables. Show that the process  $\{X(t)\}_{t\in\{0,1\}}$  given by  $X(0) = \operatorname{sign}(\eta) \xi$  and  $X(1) = \operatorname{sign}(\xi) \eta$  is not Gaussian despite each of the process values X(0) and X(1) are standard Gaussian.

**Solution.** It is an elementary exercise to see that X(0) and X(1) are standard Gaussian (normal) distributed. Also note that

$$X(0) X(1) = sign(\eta) \xi sign(\xi) \eta = |\xi| |\eta| \ge 0.$$

However, if (X(0), X(1)) were bivariate standard Gaussian (as it must be if X is a Gaussian process), then the above non-negativity is possible if and only if X(0) and X(1) have perfect correlation 1. But this is not true as

$$\mathbf{Corr}\{X(0), X(1)\} = \mathbf{Cov}\{X(0), X(1)\} = \mathbf{E}\{X(0)X(1)\} = \mathbf{E}\{|\xi||\eta|\} = (\mathbf{E}\{|\xi|\})^2 = \frac{2}{\pi} < 1$$

by elementary calculations [where we used that X(0) and X(1) are standard Gaussian].

Exercise 6. Prove that the finite dimensional distributions of a zero-mean Gaussian stochastic process  $\{X(t)\}_{t\in T}$  are completely characterized by the covariance function of the process.

**Solution.** Given  $t_1, \ldots, t_n \in T$ , the distribution of the random variable  $(X(t_1), \ldots, X(t_n))$  is determined by its characteristic function (Fourier transform)

$$\mathbf{E}\left\{e^{i\sum_{j=1}^{n}a_{j}X(t_{j})}\right\} \quad \text{for } (a_{1},\ldots,a_{n}) \in \mathbb{R}^{n}.$$

As  $\sum_{j=1}^{n} a_j X(t_j)$  is a univariate zero-mean Gaussian random variable, that characteristic function in turn is equal to

$$\exp\left[-\frac{1}{2}\operatorname{Var}\left\{\sum_{j=1}^{n}a_{j}X(t_{j})\right\}\right] = \exp\left[-\frac{1}{2}\sum_{i=1}^{n}\sum_{j=1}^{n}a_{i}a_{j}\operatorname{Cov}\left\{X(t_{i}),X(t_{j})\right\}\right],$$

which in turn obviously is determined by the covariance function of X.