## MTH 9831 Assignment 1 (9/06 - 9/12).

Read Lecture 1. Some additional references for this material are:

- 1. J. Jacod, Ph. Protter, Probability Essentials, Chapter 16.
- 2. S. Shreve, Stochastic Calculus for Finance II, Chapter 3.
- 3. A. Etheridge, A Course in Financial Calculus, Chapter 3.

Rules of the game: When solving these problems you are allowed to use (with a proper reference) any statement from lecture notes, even if it was given without a proof. Citing a convenient theorem from some other source and deriving your solution from it is not allowed.

## Solve:

- (1) Let  $B(t) = (B_1(t), B_2(t)), t \ge 0$ , be a two dimensional Brownian motion (assume that the coordinates are independent). Find the distribution of the distance from B(t) to the origin.
- (2) Let  $B(t) = (B_1(t), B_2(t))$ ,  $t \ge 0$ , be a two dimensional Brownian motion (assume that the coordinates are independent), and  $\rho \in [-1, 1]$ . Is the process  $X(t) = \rho B_1(t) + \sqrt{1 \rho^2} B_2(t)$  a Brownian motion? What is the correlation between X(t) and  $B_1(t)$ ?
- (3) Let  $(B(t))_{t\geq 0}$  be a Brownian motion. Determine, which of the following processes are Brownian motions. Justify your answers.
  - (a)  $(-B(t))_{t>0}$ ;
  - (b)  $(cB(t/c^2)_{t>0})$ , where c>0 is a constant;
  - (c)  $(\sqrt{t}B(1))_{t\geq 0}$ ;
  - (d)  $(B(2t) B(t))_{t>0}$ ;
  - (e)  $(B(s) B(s-t))_{0 \le t \le s}$ , where s is fixed.
- (4) Let  $(B(t))_{t\geq 0}$  be a Brownian motion. Define  $B^*(t) := \max_{0\leq s\leq t} B(s)$ . For  $0\leq a\leq x$  calculate  $P(B^*(t)\geq a,\,B(t)\leq x)$ .

The next two problems are review problems for Gaussian random vectors. For your convenience I first recall some linear algebra. Let  $\mathbf{x}, \mathbf{y} \in \mathbb{R}^n$ . Then *inner product* (dot product) of  $\mathbf{x}$  and  $\mathbf{y}$  is

$$\mathbf{x} \cdot \mathbf{y} = \sum_{i=1}^{n} x_i y_i,$$

the *norm* (length) of  $\mathbf{x}$  is  $\|\mathbf{x}\| = \sqrt{\mathbf{x} \cdot \mathbf{x}}$ , and  $\mathbf{x}$  and  $\mathbf{y}$  are said to be orthogonal if  $\mathbf{x} \cdot \mathbf{y} = 0$ . Given (non-zero) vectors  $\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n \in \mathbb{R}^n$  we can perform the Gram-Schmidt orthogonalization by setting

$$\begin{split} y_1 &= x_1 \\ y_2 &= x_2 - \frac{(x_2 \cdot y_1)}{(y_1 \cdot y_1)} y_1 \\ y_3 &= x_3 - \frac{(x_3 \cdot y_1)}{(y_1 \cdot y_1)} y_1 - \frac{(x_3 \cdot y_2)}{(y_2 \cdot y_2)} y_2 \\ & \dots \\ y_n &= x_n - \frac{(x_n \cdot y_1)}{(y_1 \cdot y_1)} y_1 - \frac{(x_n \cdot y_2)}{(y_2 \cdot y_2)} y_2 - \dots - \frac{(x_n \cdot y_{n-1})}{(y_{n-1} \cdot y_{n-1})} y_{n-1}. \end{split}$$

Then  $\mathbf{y_i} \cdot \mathbf{y_j} = 0$  for all  $\mathbf{i} \neq \mathbf{j}$  and the collection of all non-zero vectors<sup>1</sup> from the set  $\{\mathbf{y_1}, \mathbf{y_2}, \dots, \mathbf{y_n}\}$  forms an orthogonal basis of the span of  $\mathbf{x_1}, \mathbf{x_2}, \dots, \mathbf{x_n}$ :

$$\begin{split} x_1 &= y_1 \\ x_2 &= \frac{(x_2 \cdot y_1)}{(y_1 \cdot y_1)} y_1 + y_2 \\ x_3 &= \frac{(x_3 \cdot y_1)}{(y_1 \cdot y_1)} y_1 + \frac{(x_3 \cdot y_2)}{(y_2 \cdot y_2)} y_2 + y_3 \\ &\dots \\ x_n &= \frac{(x_n \cdot y_1)}{(y_1 \cdot y_1)} y_1 + \frac{(x_n \cdot y_2)}{(y_2 \cdot y_2)} y_2 + \dots + \frac{(x_n \cdot y_{n-1})}{(y_{n-1} \cdot y_{n-1})} y_{n-1} + y_n. \end{split}$$

Now let's go back to probability. For any square-integrable mean zero random variables X and Y define their inner product to be  $\mathbb{E}(XY)$  and the  $L^2$ -norm  $||X||_2 := \sqrt{\mathbb{E}(X^2)} = \sigma(X)$ . Note that if (X,Y) is a mean zero Gaussian vector then orthogonality of X and Y (i.e.  $\mathbb{E}(XY) = 0$ ) is equivalent to independence of X and Y.

- (5) Let  $(X_1, X_2, ..., X_n)$  be a Gaussian vector with mean zero. Write down the Gram-Schmidt procedure for this vector.
  - (a) Use the result to give an alternative proof of Theorem 1.6 from Lecture 1.
  - (b) Show that for every  $1 < k \le n$ ,  $\mathbb{E}(X_k \mid X_1, X_2, \dots, X_{k-1})$  is a linear function of  $X_1, X_2, \dots, X_{k-1}$ .
  - (c) Show that for every 1 < k < n, the conditional distribution of

$$((X_{k+1}, X_{k+2}, \dots, X_n)|(X_1, X_2, \dots, X_k) = (x_1, x_2, \dots, x_k))$$

is Gaussian and that the dependence of the parameters on  $x_1, x_2, \ldots, x_k$  is linear.

(6) Let  $(\mathbf{X_1}, \mathbf{X_2})$  be a Gaussian vector, dim  $\mathbf{X_1} = k$  and dim  $\mathbf{X_2} = m$ . The mean vector of  $(\mathbf{X_1}, \mathbf{X_2})$  is  $(\mu_1, \mu_2)$  and the covariance matrix is

$$\begin{pmatrix} C_{11} & C_{12} \\ C_{21} & C_{22} \end{pmatrix}$$

where  $C_{11}$  the covariance matrix of  $\mathbf{X_1}$ ,  $C_{22}$  is the covariance matrix of  $\mathbf{X_2}$  and  $C_{12} = C_{21}^t = \|\operatorname{Cov}(X_1^i, X_2^j)\|_{1 \leq i \leq k, 1 \leq j \leq m}$ . Assume that  $C_{11}$  is non-degenerate and find the distribution of  $(\mathbf{X_2}|\mathbf{X_1} = \mathbf{x})$ .

*Hint:* the conditional distribution is Gaussian by the previous problem. One possible way to find the parameters is to use the idea presented in class, i.e. write

$$\mathbf{X_2} - \mu_2 = \Lambda(\mathbf{X_1} - \mu_1) + \mathbf{W},$$

where  $\Lambda$  is a non-random matrix and **W** is a mean zero Gaussian vector independent from  $\mathbf{X_1}^2$ . First multiply both sides by  $(\mathbf{X_1} - \mu_1)^t$ , compute the expectation, and find  $\Lambda$ . Then multiply both parts by  $(\mathbf{X_2} - \mu_2)^t = (\Lambda(\mathbf{X_1} - \mu_1) + \mathbf{W})^t$  and compute the expectation to find the covariance matrix of **W**. Finally perform the conditioning and write down the parameters.

<sup>&</sup>lt;sup>1</sup>Some vectors  $\mathbf{y_i}$  can be zero vectors as we did not assume that the original vectors were linearly independent. We agree to replace such terms in the Gram-Schmidt process simply by  $\mathbf{0}$ .

<sup>&</sup>lt;sup>2</sup>This is a version of Gram-Schmidt idea of orthogonalization: the vector  $\Lambda(\mathbf{X}_1 - \mu_1)$  is the orthogonal (in the  $L^2$  sense) projection of  $(\mathbf{X}_2 - \mu_2)$  on  $(\mathbf{X}_1 - \mu_1)$ . Subtracting it from  $(\mathbf{X}_2 - \mu_2)$  gives a vector  $\mathbf{W}$  which is orthogonal to  $(\mathbf{X}_1 - \mu_1)$ .